Untouchability is a sin
Untouchability is a crime
Untouchability is inhuman
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Preface

The most important and crucial stage of school education is the higher secondary level. This is the transition level from a generalised curriculum to a discipline-based curriculum.

In order to pursue their career in basic sciences and professional courses, students take up Physics as one of the subjects. To provide them sufficient background to meet the challenges of academic and professional streams, the Physics textbook for Std. XI has been reformed, updated and designed to include basic information on all topics.

Each chapter starts with an introduction, followed by subject matter. All the topics are presented with clear and concise treatments. The chapters end with solved problems and self evaluation questions.

Understanding the concepts is more important than memorising. Hence it is intended to make the students understand the subject thoroughly so that they can put forth their ideas clearly. In order to make the learning of Physics more interesting, application of concepts in real life situations are presented in this book.

Due importance has been given to develop in the students, experimental and observation skills. Their learning experience would make them to appreciate the role of Physics towards the improvement of our society.

The following are the salient features of the text book.

- The data has been systematically updated.
- Figures are neatly presented.
- Self-evaluation questions (only samples) are included to sharpen the reasoning ability of the student.
- As Physics cannot be understood without the basic knowledge of Mathematics, few basic ideas and formulae in Mathematics are given.

**While preparing for the examination, students should not restrict themselves, only to the questions/problems given in the self evaluation. They must be prepared to answer the questions and problems from the text/syllabus.**

Sincere thanks to Indian Space Research Organisation (ISRO) for providing valuable information regarding the Indian satellite programme.

— Dr. S. Gunasekaran
Chairperson
SYLLABUS (180 periods)

UNIT – 1 Nature of the Physical World and Measurement (7 periods)

Physics – scope and excitement – physics in relation to technology and society.

Forces in nature – gravitational, electromagnetic and nuclear forces (qualitative ideas)

Measurement – fundamental and derived units – length, mass and time measurements.

Accuracy and precision of measuring instruments, errors in measurement – significant figures.

Dimensions - dimensions of physical quantities - dimensional analysis – applications.

UNIT – 2 Kinematics (29 periods)

Motion in a straight line – position time graph – speed and velocity – uniform and non-uniform motion – uniformly accelerated motion – relations for uniformly accelerated motions.

Scalar and vector quantities – addition and subtraction of vectors, unit vector, resolution of vectors - rectangular components, multiplication of vectors – scalar, vector products.

Motion in two dimensions – projectile motion – types of projectile – horizontal and oblique projectile.

Force and inertia, Newton’s first law of motion.


Newton’s third law of motion – law of conservation of linear momentum and its applications.

Equilibrium of concurrent forces – triangle law, parallelogram law and Lami’s theorem – experimental proof.


Work done by a constant force and a variable force – unit of work.
Energy – Kinetic energy, work – energy theorem – potential energy – power.

Collisions – Elastic and in-elastic collisions in one dimension.

**UNIT – 3 Dynamics of Rotational Motion (14 periods)**

Centre of a two particle system – generalization – applications – equilibrium of bodies, rigid body rotation and equations of rotational motion. Comparison of linear and rotational motions.

Moment of inertia and its physical significance – radius of gyration – Theorems with proof, Moment of inertia of a thin straight rod, circular ring, disc cylinder and sphere.

Moment of force, angular momentum. Torque – conservation of angular momentum.

**UNIT – 4 Gravitation and Space Science (16 periods)**

The universal law of gravitation; acceleration due to gravity and its variation with the altitude, latitude, depth and rotation of the Earth – mass of the Earth. Inertial and gravitational mass.


**UNIT – 5 Mechanics of Solids and Fluids (18 periods)**

States of matter- inter-atomic and inter-molecular forces.


Pressure due to a fluid column – Pascal’s law and its applications (hydraulic lift and hydraulic brakes) – effect of gravity on fluid pressure.
Surface energy and surface tension, angle of contact – application of surface tension in (i) formation of drops and bubbles (ii) capillary rise (iii) action of detergents.


UNIT – 6 Oscillations (12 periods)

Periodic motion – period, frequency, displacement as a function of time.

Simple harmonic motion – amplitude, frequency, phase – uniform circular motion as SHM.

Oscillations of a spring, liquid column and simple pendulum – derivation of expression for time period – restoring force – force constant.

Energy in SHM. kinetic and potential energies – law of conservation of energy.

Free, forced and damped oscillations. Resonance.

UNIT – 7 Wave Motion (17 periods)

Wave motion- longitudinal and transverse waves – relation between \( v, n, \lambda \).

Speed of wave motion in different media – Newton’s formula – Laplace’s correction.

Progressive wave – displacement equation – characteristics.


Doppler effect (quantitative idea) – applications.

UNIT – 8 Heat and Thermodynamics (17 periods)

Kinetic theory of gases – postulates – pressure of a gas – kinetic energy and temperature – degrees of freedom (mono atomic, diatomic and triatomic) – law of equipartition of energy – Avogadro’s number.

Thermal equilibrium and temperature (zeroth law of thermodynamics) Heat, work and internal energy. Specific heat – specific
heat capacity of gases at constant volume and pressure. Relation between $C_p$ and $C_v$.


UNIT – 9 Ray Optics (16 periods)

Reflection of light – reflection at plane and curved surfaces.

Total internal refelction and its applications – determination of velocity of light – Michelson’s method.


UNIT – 10 Magnetism (10 periods)

Earth’s magnetic field and magnetic elements. Bar magnet - magnetic field lines

Magnetic field due to magnetic dipole (bar magnet) along the axis and perpendicular to the axis.

Torque on a magnetic dipole (bar magnet) in a uniform magnetic field.

Tangent law – Deflection magnetometer - Tan A and Tan B positions.

Magnetic properties of materials – Intensity of magnetisation, magnetic susceptibility, magnetic induction and permeability

Dia, Para and Ferromagnetic substances with examples.

Hysteresis.
EXPERIMENTS  \(12 \times 2 = 24\) periods

1. To find the density of the material of a given wire with the help of a screw gauge and a physical balance.

2. Simple pendulum - To draw graphs between (i) \(L\) and \(T\) and (ii) \(L\) and \(T^2\) and to decide which is better. Hence to determine the acceleration due to gravity.

3. Measure the mass and dimensions of (i) cylinder and (ii) solid sphere using the vernier calipers and physical balance. Calculate the moment of inertia.

4. To determine Young’s modulus of the material of a given wire by using Searles’ apparatus.

5. To find the spring constant of a spring by method of oscillations.

6. To determine the coefficient of viscosity by Poiseuille’s flow method.

7. To determine the coefficient of viscosity of a given viscous liquid by measuring the terminal velocity of a given spherical body.

8. To determine the surface tension of water by capillary rise method.

9. To verify the laws of a stretched string using a sonometer.

10. To find the velocity of sound in air at room temperature using the resonance column apparatus.

11. To determine the focal length of a concave mirror

12. To map the magnetic field due to a bar magnet placed in the magnetic meridian with its (i) north pole pointing South and (ii) north pole pointing North and locate the null points.
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(Unit 6 to 10 continues in Volume II)
1. Nature of the Physical World and Measurement

The history of humans reveals that they have been making continuous and serious attempts to understand the world around them. The repetition of day and night, cycle of seasons, volcanoes, rainbows, eclipses and the starry night sky have always been a source of wonder and subject of thought. The inquiring mind of humans always tried to understand the natural phenomena by observing the environment carefully. This pursuit of understanding nature led us to today’s modern science and technology.

1.1 Physics

The word science comes from a Latin word “scientia” which means ‘to know’. Science is nothing but the knowledge gained through the systematic observations and experiments. Scientific methods include the systematic observations, reasoning, modelling and theoretical prediction. Science has many disciplines, physics being one of them. The word physics has its origin in a Greek word meaning ‘nature’. Physics is the most basic science, which deals with the study of nature and natural phenomena. Understanding science begins with understanding physics. With every passing day, physics has brought to us deeper levels of understanding of nature.

Physics is an empirical study. Everything we know about physical world and about the principles that govern its behaviour has been learned through observations of the phenomena of nature. The ultimate test of any physical theory is its agreement with observations and measurements of physical phenomena. Thus physics is inherently a science of measurement.

1.1.1 Scope of Physics

The scope of physics can be understood if one looks at its various sub-disciplines such as mechanics, optics, heat and thermodynamics, electrodynamics, atomic physics, nuclear physics, etc.
Mechanics deals with motion of particles and general systems of particles. The working of telescopes, colours of thin films are the topics dealt in optics. Heat and thermodynamics deals with the pressure - volume changes that take place in a gas when its temperature changes, working of refrigerator, etc. The phenomena of charged particles and magnetic bodies are dealt in electrodynamics. The magnetic field around a current carrying conductor, propagation of radio waves etc. are the areas where electrodynamics provide an answer. Atomic and nuclear physics deals with the constitution and structure of matter, interaction of atoms and nuclei with electrons, photons and other elementary particles.

Foundation of physics enables us to appreciate and enjoy things and happenings around us. The laws of physics help us to understand and comprehend the cause-effect relationships in what we observe. This makes a complex problem to appear pretty simple.

Physics is exciting in many ways. To some, the excitement comes from the fact that certain basic concepts and laws can explain a range of phenomena. For some others, the thrill lies in carrying out new experiments to unravel the secrets of nature. Applied physics is even more interesting. Transforming laws and theories into useful applications require great ingenuity and persistent effort.

1.1.2 Physics, Technology and Society

Technology is the application of the doctrines in physics for practical purposes. The invention of steam engine had a great impact on human civilization. Till 1933, Rutherford did not believe that energy could be tapped from atoms. But in 1938, Hann and Meitner discovered neutron-induced fission reaction of uranium. This is the basis of nuclear weapons and nuclear reactors. The contribution of physics in the development of alternative resources of energy is significant. We are consuming the fossil fuels at such a very fast rate that there is an urgent need to discover new sources of energy which are cheap. Production of electricity from solar energy and geothermal energy is a reality now, but we have a long way to go. Another example of physics giving rise to technology is the integrated chip, popularly called as IC. The development of newer ICs and faster processors made the computer industry to grow leaps and bounds in the last two decades. Computers have become affordable now due to improved production techniques
and low production costs.

The legitimate purpose of technology is to serve people. Our society is becoming more and more science-oriented. We can become better members of society if we develop an understanding of the basic laws of physics.

1.2 Forces of nature

Sir Issac Newton was the first one to give an exact definition for force.

"Force is the external agency applied on a body to change its state of rest and motion".

There are four basic forces in nature. They are gravitational force, electromagnetic force, strong nuclear force and weak nuclear force.

Gravitational force

It is the force between any two objects in the universe. It is an attractive force by virtue of their masses. By Newton’s law of gravitation, the gravitational force is directly proportional to the product of the masses and inversely proportional to the square of the distance between them. Gravitational force is the weakest force among the fundamental forces of nature but has the greatest large-scale impact on the universe. Unlike the other forces, gravity works universally on all matter and energy, and is universally attractive.

Electromagnetic force

It is the force between charged particles such as the force between two electrons, or the force between two current carrying wires. It is attractive for unlike charges and repulsive for like charges. The electromagnetic force obeys inverse square law. It is very strong compared to the gravitational force. It is the combination of electrostatic and magnetic forces.

Strong nuclear force

It is the strongest of all the basic forces of nature. It, however, has the shortest range, of the order of $10^{-15}$ m. This force holds the protons and neutrons together in the nucleus of an atom.
Weak nuclear force

Weak nuclear force is important in certain types of nuclear process such as $\beta$-decay. This force is not as weak as the gravitational force.

1.3 Measurement

Physics can also be defined as the branch of science dealing with the study of properties of materials. To understand the properties of materials, measurement of physical quantities such as length, mass, time etc., are involved. The uniqueness of physics lies in the measurement of these physical quantities.

1.3.1 Fundamental quantities and derived quantities

Physical quantities can be classified into two namely, fundamental quantities and derived quantities. Fundamental quantities are quantities which cannot be expressed in terms of any other physical quantity. For example, quantities like length, mass, time, temperature are fundamental quantities. Quantities that can be expressed in terms of fundamental quantities are called derived quantities. Area, volume, density etc. are examples for derived quantities.

1.3.2 Unit

To measure a quantity, we always compare it with some reference standard. To say that a rope is 10 metres long is to say that it is 10 times as long as an object whose length is defined as 1 metre. Such a standard is called a unit of the quantity.

Therefore, unit of a physical quantity is defined as the established standard used for comparison of the given physical quantity.

The units in which the fundamental quantities are measured are called fundamental units and the units used to measure derived quantities are called derived units.

1.3.3 System International de Units (SI system of units)

In earlier days, many system of units were followed to measure physical quantities. The British system of foot-pound–second or fps system, the Gaussian system of centimetre – gram – second or cgs system, the metre-kilogram – second or the mks system were the three
systems commonly followed. To bring uniformity, the General Conference on Weights and Measures in the year 1960, accepted the SI system of units. This system is essentially a modification over mks system and is, therefore rationalised mksA (metre kilogram second ampere) system. This rationalisation was essential to obtain the units of all the physical quantities in physics.

In the SI system of units there are seven fundamental quantities and two supplementary quantities. They are presented in Table 1.1.

<table>
<thead>
<tr>
<th>Physical quantity</th>
<th>Unit</th>
<th>Symbol</th>
</tr>
</thead>
<tbody>
<tr>
<td><strong>Fundamental quantities</strong></td>
<td></td>
<td></td>
</tr>
<tr>
<td>Length</td>
<td>metre</td>
<td>m</td>
</tr>
<tr>
<td>Mass</td>
<td>kilogram</td>
<td>kg</td>
</tr>
<tr>
<td>Time</td>
<td>second</td>
<td>s</td>
</tr>
<tr>
<td>Electric current</td>
<td>ampere</td>
<td>A</td>
</tr>
<tr>
<td>Temperature</td>
<td>kelvin</td>
<td>K</td>
</tr>
<tr>
<td>Luminous intensity</td>
<td>candela</td>
<td>cd</td>
</tr>
<tr>
<td>Amount of substance</td>
<td>mole</td>
<td>mol</td>
</tr>
<tr>
<td><strong>Supplementary quantities</strong></td>
<td></td>
<td></td>
</tr>
<tr>
<td>Plane angle</td>
<td>radian</td>
<td>rad</td>
</tr>
<tr>
<td>Solid angle</td>
<td>steradian</td>
<td>sr</td>
</tr>
</tbody>
</table>

### 1.3.4 Uniqueness of SI system

The SI system is logically far superior to all other systems. The SI units have certain special features which make them more convenient in practice. Permanence and reproduceability are the two important characteristics of any unit standard. The SI standards do not vary with time as they are based on the properties of atoms. Further SI system of units are coherent system of units, in which the units of derived quantities are obtained as multiples or submultiples of certain basic units. Table 1.2 lists some of the derived quantities and their units.
### 1.3.5 SI standards

#### Length

Length is defined as the distance between two points. The SI unit of length is metre.

One standard metre is equal to 1 650 763.73 wavelengths of the orange – red light emitted by the individual atoms of krypton – 86 in a krypton discharge lamp.

#### Mass

Mass is the quantity of matter contained in a body. It is independent of temperature and pressure. It does not vary from place to place.

---

<table>
<thead>
<tr>
<th>Physical Quantity</th>
<th>Expression</th>
<th>Unit</th>
</tr>
</thead>
<tbody>
<tr>
<td>Area</td>
<td>length × breadth</td>
<td>m²</td>
</tr>
<tr>
<td>Volume</td>
<td>area × height</td>
<td>m³</td>
</tr>
<tr>
<td>Velocity</td>
<td>displacement / time</td>
<td>m s⁻¹</td>
</tr>
<tr>
<td>Acceleration</td>
<td>velocity / time</td>
<td>m s⁻²</td>
</tr>
<tr>
<td>Angular velocity</td>
<td>angular displacement / time</td>
<td>rad s⁻¹</td>
</tr>
<tr>
<td>Angular acceleration</td>
<td>angular velocity / time</td>
<td>rad s⁻²</td>
</tr>
<tr>
<td>Density</td>
<td>mass / volume</td>
<td>kg m⁻³</td>
</tr>
<tr>
<td>Momentum</td>
<td>mass × velocity</td>
<td>kg m s⁻¹</td>
</tr>
<tr>
<td>Moment of inertia</td>
<td>mass × (distance)²</td>
<td>kg m²</td>
</tr>
<tr>
<td>Force</td>
<td>mass × acceleration</td>
<td>kg m s⁻² or N</td>
</tr>
<tr>
<td>Pressure</td>
<td>force / area</td>
<td>N m⁻² or Pa</td>
</tr>
<tr>
<td>Energy (work)</td>
<td>force × distance</td>
<td>N m or J</td>
</tr>
<tr>
<td>Impulse</td>
<td>force × time</td>
<td>N s</td>
</tr>
<tr>
<td>Surface tension</td>
<td>force / length</td>
<td>N m⁻¹</td>
</tr>
<tr>
<td>Moment of force (torque)</td>
<td>force × distance</td>
<td>N m</td>
</tr>
<tr>
<td>Electric charge</td>
<td>current × time</td>
<td>A s</td>
</tr>
<tr>
<td>Current density</td>
<td>current / area</td>
<td>A m⁻²</td>
</tr>
<tr>
<td>Magnetic induction</td>
<td>force / (current × length)</td>
<td>N A⁻¹ m⁻¹</td>
</tr>
</tbody>
</table>
to place. The SI unit of mass is kilogram.

The kilogram is equal to the mass of the international prototype of the kilogram (a plantinum–iridium alloy cylinder) kept at the International Bureau of Weights and Measures at Sevres, near Paris, France.

An atomic standard of mass has not yet been adopted because it is not yet possible to measure masses on an atomic scale with as much precision as on a macroscopic scale.

**Time**

Until 1960 the standard of time was based on the mean solar day, the time interval between successive passages of the sun at its highest point across the meridian. It is averaged over an year. In 1967, an atomic standard was adopted for second, the SI unit of time.

One standard second is defined as the time taken for $9\,192\,631\,770$ periods of the radiation corresponding to unperturbed transition between hyperfine levels of the ground state of cesium–133 atom. Atomic clocks are based on this. In atomic clocks, an error of one second occurs only in 5000 years.

**Ampere**

The ampere is the constant current which, flowing through two straight parallel infinitely long conductors of negligible cross-section, and placed in vacuum 1 m apart, would produce between the conductors a force of $2 \times 10^{-7}$ newton per unit length of the conductors.

**Kelvin**

The Kelvin is the fraction of $\frac{1}{273.16}$ of the thermodynamic temperature of the triple point of water*.

**Candela**

The candela is the luminous intensity in a given direction due to a

---

* Triple point of water is the temperature at which saturated water vapour, pure water and melting ice are all in equilibrium. The triple point temperature of water is 273.16 K.
source, which emits monochromatic radiation of frequency \(540 \times 10^{12}\) Hz and of which the radiant intensity in that direction is \(\frac{1}{683}\) watt per steradian.

**Mole**

The mole is the amount of substance which contains as many elementary entities as there are atoms in 0.012 kg of carbon-12.

1.3.6 *Rules and conventions for writing SI units and their symbols*

1. The units named after scientists are not written with a capital initial letter.
   
   For example: newton, henry, watt

2. The symbols of the units named after scientist should be written by a capital letter.
   
   For example: N for newton, H for henry, W for watt

3. Small letters are used as symbols for units not derived from a proper name.
   
   For example: m for metre, kg for kilogram

4. No full stop or other punctuation marks should be used within or at the end of symbols.
   
   For example: 50 m and not as 50 m.

5. The symbols of the units do not take plural form.
   
   For example: 10 kg not as 10 kgs

6. When temperature is expressed in kelvin, the degree sign is omitted.
   
   For example: 273 K not as 273° K
   
   (If expressed in Celsius scale, degree sign is to be included. For example 100° C and not 100 C)

7. Use of solidus is recommended only for indicating a division of one letter unit symbol by another unit symbol. Not more than one solidus is used.
   
   For example: m s\(^{-1}\) or m/s, J/K mol or J K\(^{-1}\) mol\(^{-1}\) but not J/K/mol.
8. Some space is always to be left between the number and the symbol of the unit and also between the symbols for compound units such as force, momentum, etc.

For example, it is not correct to write 2.3m. The correct representation is 2.3 m; kg m s\(^{-2}\) and not as kgms\(^{-2}\).

9. Only accepted symbols should be used.

For example: ampere is represented as A and not as amp. or am; second is represented as s and not as sec.

10. Numerical value of any physical quantity should be expressed in scientific notation.

For an example, density of mercury is \(1.36 \times 10^4\) kg m\(^{-3}\) and not as 13600 kg m\(^{-3}\).

### 1.4 Expressing larger and smaller physical quantities

Once the fundamental units are defined, it is easier to express larger and smaller units of the same physical quantity. In the metric (SI) system these are related to the fundamental unit in multiples of 10 or 1/10. Thus 1 km is 1000 m and 1 mm is 1/1000 metre. Table 1.3 lists the standard SI prefixes, their meanings and abbreviations.

In order to measure very large distances, the following units are used.

**(i) Light year**

Light year is the distance travelled by light in one year in vacuum.

<table>
<thead>
<tr>
<th>Power of ten</th>
<th>Prefix</th>
<th>Abbreviation</th>
</tr>
</thead>
<tbody>
<tr>
<td>(10^{-15})</td>
<td>femto</td>
<td>f</td>
</tr>
<tr>
<td>(10^{-12})</td>
<td>pico</td>
<td>p</td>
</tr>
<tr>
<td>(10^{-9})</td>
<td>nano</td>
<td>n</td>
</tr>
<tr>
<td>(10^{-6})</td>
<td>micro</td>
<td>µ</td>
</tr>
<tr>
<td>(10^{-3})</td>
<td>milli</td>
<td>m</td>
</tr>
<tr>
<td>(10^{-2})</td>
<td>centi</td>
<td>c</td>
</tr>
<tr>
<td>(10^{-1})</td>
<td>deci</td>
<td>d</td>
</tr>
<tr>
<td>(10^1)</td>
<td>deca</td>
<td>da</td>
</tr>
<tr>
<td>(10^2)</td>
<td>hecto</td>
<td>h</td>
</tr>
<tr>
<td>(10^3)</td>
<td>kilo</td>
<td>k</td>
</tr>
<tr>
<td>(10^6)</td>
<td>mega</td>
<td>M</td>
</tr>
<tr>
<td>(10^9)</td>
<td>giga</td>
<td>G</td>
</tr>
<tr>
<td>(10^{12})</td>
<td>tera</td>
<td>T</td>
</tr>
<tr>
<td>(10^{15})</td>
<td>peta</td>
<td>P</td>
</tr>
</tbody>
</table>
Distance travelled = velocity of light × 1 year
\[
\therefore 1 \text{ light year} = 3 \times 10^8 \text{ m s}^{-1} \times 1 \text{ year (in seconds)}
\]
\[
= 3 \times 10^8 \times 365.25 \times 24 \times 60 \times 60
\]
\[
= 9.467 \times 10^{15} \text{ m}
\]
\[
1 \text{ light year} = 9.467 \times 10^{15} \text{ m}
\]

\text{(ii) Astronomical unit}

Astronomical unit is the mean distance of the centre of the Sun from the centre of the Earth.

\[
1 \text{ Astronomical unit (AU)} = 1.496 \times 10^{11} \text{ m}
\]

1.5 Determination of distance

For measuring large distances such as the distance of moon or a planet from the Earth, special methods are adopted. Radio-echo method, laser pulse method and parallax method are used to determine very large distances.

\text{Laser pulse method}

The distance of moon from the Earth can be determined using laser pulses. The laser pulses are beamed towards the moon from a powerful transmitter. These pulses are reflected back from the surface of the moon. The time interval between sending and receiving of the signal is determined very accurately.

If \( t \) is the time interval and \( c \) the velocity of the laser pulses, then

\[
\text{the distance of the moon from the Earth is } d = \frac{ct}{2}.
\]

1.6 Determination of mass

The conventional method of finding the mass of a body in the laboratory is by physical balance. The mass can be determined to an accuracy of 1 mg. Now–a–days, digital balances are used to find the mass very accurately. The advantage of digital balance is that the mass of the object is determined at once.

1.7 Measurement of time

We need a clock to measure any time interval. Atomic clocks provide better standard for time. Some techniques to measure time interval are given below.
Quartz clocks

The *piezo-electric property* of a crystal is the principle of quartz clock. These clocks have an accuracy of one second in every $10^9$ seconds.

Atomic clocks

These clocks make use of periodic vibration taking place within the atom. Atomic clocks have an accuracy of 1 part in $10^{13}$ seconds.

1.8 Accuracy and precision of measuring instruments

All measurements are made with the help of instruments. The accuracy to which a measurement is made depends on several factors. For example, if length is measured using a metre scale which has graduations at 1 mm interval then all readings are good only upto this value. The error in the use of any instrument is normally taken to be half of the smallest division on the scale of the instrument. Such an error is called instrumental error. In the case of a metre scale, this error is about 0.5 mm.

Physical quantities obtained from experimental observation always have some uncertainty. Measurements can never be made with absolute precision. Precision of a number is often indicated by following it with $\pm$ symbol and a second number indicating the maximum error likely.

For example, if the length of a steel rod = 56.47 $\pm$ 3 mm then the true length is unlikely to be less than 56.44 mm or greater than 56.50 mm. If the error in the measured value is expressed in fraction, it is called fractional error and if expressed in percentage it is called percentage error. For example, a resistor labelled “470 $\Omega$, 10%” probably has a true resistance differing not more than 10% from 470 $\Omega$. So the true value lies between 423 $\Omega$ and 517 $\Omega$.

1.8.1 Significant figures

The digits which tell us the number of units we are reasonably sure of having counted in making a measurement are called significant figures. Or in other words, the number of meaningful digits in a number is called the number of significant figures. A choice of change of different units does not change the number of significant digits or figures in a measurement.

*When pressure is applied along a particular axis of a crystal, an electric potential difference is developed in a perpendicular axis.*
For example, 2.868 cm has four significant figures. But in different units, the same can be written as 0.02868 m or 28.68 mm or 28680 µm. All these numbers have the same four significant figures.

From the above example, we have the following rules.

i) All the non-zero digits in a number are significant.

ii) All the zeroes between two non-zeroes digits are significant, irrespective of the decimal point.

iii) If the number is less than 1, the zeroes on the right of decimal point but to the left of the first non-zero digit are not significant. (In 0.02868 the underlined zeroes are not significant).

iv) The zeroes at the end without a decimal point are not significant. (In 23080 µm, the trailing zero is not significant).

v) The trailing zeroes in a number with a decimal point are significant. (The number 0.07100 has four significant digits).

**Examples**

i) 30700 has three significant figures.

ii) 132.73 has five significant figures.

iii) 0.00345 has three and

iv) 40.00 has four significant figures.

### 1.8.2 Rounding off

Calculators are widely used now–a–days to do the calculations. The result given by a calculator has too many figures. In no case the result should have more significant figures than the figures involved in the data used for calculation. The result of calculation with number containing more than one uncertain digit, should be rounded off. The technique of rounding off is followed in applied areas of science.

A number 1.876 rounded off to three significant digits is 1.88 while the number 1.872 would be 1.87. The rule is that if the insignificant digit (underlined) is more than 5, the preceeding digit is raised by 1, and is left unchanged if the former is less than 5.

If the number is 2.845, the insignificant digit is 5. In this case, the convention is that if the preceeding digit is even, the insignificant digit is simply dropped and, if it is odd, the preceeding digit is raised by 1. Following this, 2.845 is rounded off to 2.84 where as 2.815 is rounded off to 2.82.
**Examples**

1. Add 17.35 kg, 25.8 kg and 9.423 kg.

   Of the three measurements given, 25.8 kg is the least accurately known.

   \[\therefore 17.35 + 25.8 + 9.423 = 52.573 \text{ kg}\]

   Correct to three significant figures, 52.573 kg is written as 52.6 kg

2. Multiply 3.8 and 0.125 with due regard to significant figures.

\[3.8 \times 0.125 = 0.475\]

   The least number of significant figure in the given quantities is 2. Therefore the result should have only two significant figures.

\[\therefore 3.8 \times 0.125 = 0.475 = 0.48\]

1.8.3 Errors in Measurement

The uncertainty in the measurement of a physical quantity is called error. It is the difference between the true value and the measured value of the physical quantity. Errors may be classified into many categories.

(i) **Constant errors**

   It is the same error repeated every time in a series of observations. Constant error is due to faulty calibration of the scale in the measuring instrument. In order to minimise constant error, measurements are made by different possible methods and the mean value so obtained is regarded as the true value.

(ii) **Systematic errors**

   These are errors which occur due to a certain pattern or system. These errors can be minimised by identifying the source of error. Instrumental errors, personal errors due to individual traits and errors due to external sources are some of the systematic errors.

(iii) **Gross errors**

   Gross errors arise due to one or more than one of the following reasons.

   (1) Improper setting of the instrument.
Wrong recordings of the observation.

(3) Not taking into account sources of error and precautions.

(4) Usage of wrong values in the calculation.

Gross errors can be minimised only if the observer is very careful in his observations and sincere in his approach.

(iv) Random errors

It is very common that repeated measurements of a quantity give values which are slightly different from each other. These errors have no set pattern and occur in a random manner. Hence they are called random errors. They can be minimised by repeating the measurements many times and taking the arithmetic mean of all the values as the correct reading.

The most common way of expressing an error is percentage error. If the accuracy in measuring a quantity $x$ is $\Delta x$, then the percentage error in $x$ is given by $\frac{\Delta x}{x} \times 100\%$.

1.9 Dimensional Analysis

Dimensions of a physical quantity are the powers to which the fundamental quantities must be raised.

We know that velocity $= \frac{\text{displacement}}{\text{time}}$

$= \frac{[L]}{[T]}$

$= [M^0L^1T^{-1}]$

where [M], [L] and [T] are the dimensions of the fundamental quantities mass, length and time respectively.

Therefore velocity has zero dimension in mass, one dimension in length and $-1$ dimension in time. Thus the dimensional formula for velocity is $[M^0L^1T^{-1}]$ or simply $[LT^{-1}]$. The dimensions of fundamental quantities are given in Table 1.4 and the dimensions of some derived quantities are given in Table 1.5.
### Table 1.4 Dimensions of fundamental quantities

<table>
<thead>
<tr>
<th>Fundamental quantity</th>
<th>Dimension</th>
</tr>
</thead>
<tbody>
<tr>
<td>Length</td>
<td>L</td>
</tr>
<tr>
<td>Mass</td>
<td>M</td>
</tr>
<tr>
<td>Time</td>
<td>T</td>
</tr>
<tr>
<td>Temperature</td>
<td>K</td>
</tr>
<tr>
<td>Electric current</td>
<td>A</td>
</tr>
<tr>
<td>Luminous intensity</td>
<td>cd</td>
</tr>
<tr>
<td>Amount of substance</td>
<td>mol</td>
</tr>
</tbody>
</table>

### Table 1.5 Dimensional formulae of some derived quantities

<table>
<thead>
<tr>
<th>Physical quantity</th>
<th>Expression</th>
<th>Dimensional formula</th>
</tr>
</thead>
<tbody>
<tr>
<td>Area</td>
<td>length × breadth</td>
<td>[L²]</td>
</tr>
<tr>
<td>Density</td>
<td>mass / volume</td>
<td>[ML⁻³]</td>
</tr>
<tr>
<td>Acceleration</td>
<td>velocity / time</td>
<td>[LT⁻²]</td>
</tr>
<tr>
<td>Momentum</td>
<td>mass × velocity</td>
<td>[MLT⁻¹]</td>
</tr>
<tr>
<td>Force</td>
<td>mass × acceleration</td>
<td>[MLT⁻²]</td>
</tr>
<tr>
<td>Work</td>
<td>force × distance</td>
<td>[ML²T⁻²]</td>
</tr>
<tr>
<td>Power</td>
<td>work / time</td>
<td>[ML²T⁻³]</td>
</tr>
<tr>
<td>Energy</td>
<td>work</td>
<td>[ML²T⁻²]</td>
</tr>
<tr>
<td>Impulse</td>
<td>force × time</td>
<td>[MLT⁻¹]</td>
</tr>
<tr>
<td>Radius of gyration</td>
<td>distance</td>
<td>[L]</td>
</tr>
<tr>
<td>Pressure</td>
<td>force / area</td>
<td>[ML⁻¹T⁻²]</td>
</tr>
<tr>
<td>Surface tension</td>
<td>force / length</td>
<td>[MT⁻²]</td>
</tr>
<tr>
<td>Frequency</td>
<td>1 / time period</td>
<td>[T⁻¹]</td>
</tr>
<tr>
<td>Tension</td>
<td>force</td>
<td>[MLT⁻²]</td>
</tr>
<tr>
<td>Moment of force (or torque)</td>
<td>force × distance</td>
<td>[ML²T⁻²]</td>
</tr>
<tr>
<td>Angular velocity</td>
<td>angular displacement / time</td>
<td>[T⁻¹]</td>
</tr>
<tr>
<td>Stress</td>
<td>force / area</td>
<td>[ML⁻¹T⁻²]</td>
</tr>
<tr>
<td>Heat</td>
<td>energy</td>
<td>[ML²T⁻²]</td>
</tr>
<tr>
<td>Heat capacity</td>
<td>heat energy / temperature</td>
<td>[ML²T⁻²K⁻¹]</td>
</tr>
<tr>
<td>Charge</td>
<td>current × time</td>
<td>[AT]</td>
</tr>
<tr>
<td>Faraday constant</td>
<td>Avogadro constant ×</td>
<td>[AT mol⁻¹]</td>
</tr>
<tr>
<td>Magnetic induction</td>
<td>force / (current × length)</td>
<td>[MT⁻² A⁻¹]</td>
</tr>
</tbody>
</table>
**Dimensional quantities**

Constants which possess dimensions are called dimensional constants. Planck's constant, universal gravitational constant are dimensional constants.

Dimensional variables are those physical quantities which possess dimensions but do not have a fixed value. Example – velocity, force, etc.

**Dimensionless quantities**

There are certain quantities which do not possess dimensions. They are called dimensionless quantities. Examples are strain, angle, specific gravity, etc. They are dimensionless as they are the ratio of two quantities having the same dimensional formula.

**Principle of homogeneity of dimensions**

An equation is dimensionally correct if the dimensions of the various terms on either side of the equation are the same. This is called the principle of homogeneity of dimensions. This principle is based on the fact that two quantities of the same dimension only can be added up, the resulting quantity also possessing the same dimension.

The equation $A + B = C$ is valid only if the dimensions of $A$, $B$ and $C$ are the same.

**1.9.1 Uses of dimensional analysis**

The method of dimensional analysis is used to

(i) convert a physical quantity from one system of units to another.

(ii) check the dimensional correctness of a given equation.

(iii) establish a relationship between different physical quantities in an equation.

**1.9.1(i) To convert a physical quantity from one system of units to another**

Given the value of $G$ in cgs system is $6.67 \times 10^{-8} \text{dyne cm}^2 \text{g}^{-2}$.

Calculate its value in SI units.

<table>
<thead>
<tr>
<th>In cgs system</th>
<th>In SI system</th>
</tr>
</thead>
<tbody>
<tr>
<td>$G_{\text{cgs}} = 6.67 \times 10^{-8}$</td>
<td>$G =$</td>
</tr>
<tr>
<td>$M_1 = 1 \text{g}$</td>
<td>$M_2 = 1 \text{ kg}$</td>
</tr>
<tr>
<td>$L_1 = 1 \text{ cm}$</td>
<td>$L_2 = 1 \text{ m}$</td>
</tr>
<tr>
<td>$T_1 = 1 \text{ s}$</td>
<td>$T_2 = 1 \text{ s}$</td>
</tr>
</tbody>
</table>
The dimensional formula for gravitational constant is \([M^{-1}L^3T^{-2}]\).

In cgs system, dimensional formula for \(G\) is \([M_1^x L_1^y T_1^z]\)

In SI system, dimensional formula for \(G\) is \([M_2^x L_2^y T_2^z]\)

Here \(x = -1,\ y = 3,\ z = -2\)

\[
\therefore \quad G \left[ M_2^x L_2^y T_2^z \right] = G_{\text{cgs}} \left[ M_1^x L_1^y T_1^z \right]
\]

or

\[
G = G_{\text{cgs}} \left[ \frac{M_1}{M_2} \right]^x \left[ \frac{L_1}{L_2} \right]^y \left[ \frac{T_1}{T_2} \right]^z
\]

\[
= 6.67 \times 10^{-8} \left[ \frac{1 \, g}{1 \, kg} \right]^{-1} \left[ \frac{1 \, cm}{1 \, m} \right]^3 \left[ \frac{1 \, s}{1 \, s} \right]^{-2}
\]

\[
= 6.67 \times 10^{-8} \left[ \frac{1 \, g}{1000 \, g} \right]^{-1} \left[ \frac{1 \, cm}{100 \, cm} \right]^3 \left[ 1 \right]^{-2}
\]

\[
= 6.67 \times 10^{-11}
\]

\[
\therefore \quad \text{In SI units,}\quad G = 6.67 \times 10^{-11} \, \text{N} \, \text{m}^2 \, \text{kg}^{-2}
\]

(ii) **To check the dimensional correctness of a given equation**

Let us take the equation of motion

\[
s = ut + \left( \frac{1}{2} \right)at^2
\]

Applying dimensions on both sides,

\[
[L] = [LT^{-1}] [T] + [LT^{-2}] [T^2]
\]

\(\left( \frac{1}{2} \right)\) is a constant having no dimension\)

\[
[L] = [L] + [L]
\]

As the dimensions on both sides are the same, the equation is dimensionally correct.

(iii) **To establish a relationship between the physical quantities in an equation**

Let us find an expression for the time period \(T\) of a simple pendulum. The time period \(T\) may depend upon (i) mass \(m\) of the bob (ii) length \(l\) of the pendulum and (iii) acceleration due to gravity \(g\) at the place where the pendulum is suspended.
\[ T \propto m^x l^y g^z \]
or \[ T = k m^x l^y g^z \quad \text{...(1)} \]
where \( k \) is a dimensionless constant of proportionality. Rewriting equation (1) with dimensions,
\[ [T^1] = [M^x] [L^y] [L^{T-2}] \]
\[ [T^1] = [M^x L^y + z T^{-2z}] \]
Comparing the powers of M, L and T on both sides
\[ x = 0, \ y + z = 0 \text{ and } -2z = 1 \]
Solving for \( x, y \) and \( z \), \( x = 0, \ y = \frac{1}{2} \text{ and } z = -\frac{1}{2} \)
From equation (1), \[ T = k m^{\frac{1}{2}} l^{\frac{1}{2}} g^{-\frac{1}{2}} \]
\[ T = k \left[ \frac{1}{g} \right]^{1/2} = k \frac{T}{\sqrt{g}} \]
Experimentally the value of \( k \) is determined to be \( 2\pi \).
\[ \therefore T = 2\pi \frac{T}{\sqrt{g}} \]

1.9.2 Limitations of Dimensional Analysis

(i) The value of dimensionless constants cannot be determined by this method.

(ii) This method cannot be applied to equations involving exponential and trigonometric functions.

(iii) It cannot be applied to an equation involving more than three physical quantities.

(iv) It can check only whether a physical relation is dimensionally correct or not. It cannot tell whether the relation is absolutely correct or not. For example applying this technique \( s = ut + \frac{1}{4} at^2 \) is dimensionally correct whereas the correct relation is \( s = ut + \frac{1}{2} at^2 \).
Solved Problems

1.1 A laser signal is beamed towards a distant planet from the Earth and its reflection is received after seven minutes. If the distance between the planet and the Earth is $6.3 \times 10^{10}$ m, calculate the velocity of the signal.

**Data** : $d = 6.3 \times 10^{10}$ m, $t = 7$ minutes $= 7 \times 60 = 420$ s

**Solution**: If $d$ is the distance of the planet, then total distance travelled by the signal is $2d$.

\[
\text{velocity} = \frac{2d}{t} = \frac{2 \times 6.3 \times 10^{10}}{420} = 3 \times 10^8 \text{ m s}^{-1}
\]

1.2 A goldsmith put a ruby in a box weighing 1.2 kg. Find the total mass of the box and ruby applying principle of significant figures. The mass of the ruby is 5.42 g.

**Data** : Mass of box $= 1.2$ kg

Mass of ruby $= 5.42$ g $= 5.42 \times 10^{-3}$ kg $= 0.00542$ kg

**Solution**: Total mass $= \text{mass of box} + \text{mass of ruby}$

$= 1.2 + 0.00542 = 1.20542$ kg

After rounding off, total mass $= 1.2$ kg

1.3 Check whether the equation $\lambda = \frac{h}{mv}$ is dimensionally correct ($\lambda$ - wavelength, $h$ - Planck’s constant, $m$ - mass, $v$ - velocity).

**Solution**: Dimension of Planck’s constant $h$ is $[ML^2T^{-1}]$

Dimension of $\lambda$ is $[L]$

Dimension of $m$ is $[M]$

Dimension of $v$ is $[LT^{-1}]$

Rewriting $\lambda = \frac{h}{mv}$ using dimension

$[L] = \left[\frac{[ML^2T^{-1}]}{[M][LT^{-1}]}\right]$

$[L] = [L]$

As the dimensions on both sides of the equation are same, the given equation is dimensionally correct.
1.4 Multiply 2.2 and 0.225. Give the answer correct to significant figures.

**Solution**: \(2.2 \times 0.225 = 0.495\)

Since the least number of significant figure in the given data is 2, the result should also have only two significant figures.

\[2.2 \times 0.225 = 0.50\]

1.5 Convert 76 cm of mercury pressure into N m\(^{-2}\) using the method of dimensions.

**Solution**: In cgs system, 76 cm of mercury

pressure = \(76 \times 13.6 \times 980 \text{ dyne cm}^{-2}\)

Let this be \(P_1\). Therefore \(P_1 = 76 \times 13.6 \times 980 \text{ dyne cm}^{-2}\)

In cgs system, the dimension of pressure is \([M_1^{a}L_1^{b}T_1^{c}]\)

Dimension of pressure is \([ML^{-1}T^{-2}]\). Comparing this with \([M_2^{a}L_2^{b}T_2^{c}]\)

we have \(a = 1, b = -1 \text{ and } c = -2\).

\[\therefore \text{Pressure in SI system } P_2 = P_1 \left[ \frac{M_1}{M_2} \right]^{a} \left[ \frac{L_2}{L_1} \right]^{b} \left[ \frac{T_2}{T_1} \right]^{c}\]

\[\text{ie } P_2 = 76 \times 13.6 \times 980 \left[ \frac{10^{-3}\text{ kg}}{1\text{ kg}} \right]^{1} \left[ \frac{10^{-2}\text{ m}}{1\text{ m}} \right]^{-1} \left[ \frac{1\text{ s}}{1\text{ s}} \right]^{-2}\]

\[= 76 \times 13.6 \times 980 \times 10^{-3} \times 10^{2}\]

\[= 101292.8 \text{ N m}^{-2}\]

\[P_2 = 1.01 \times 10^5 \text{ N m}^{-2}\]
Self evaluation

(The questions and problems given in this self evaluation are only samples. In the same way any question and problem could be framed from the text matter. Students must be prepared to answer any question and problem from the text matter, not only from the self evaluation.)

1.1 Which of the following are equivalent?
(a) 6400 km and $6.4 \times 10^8$ cm  
(b) $2 \times 10^4$ cm and $2 \times 10^6$ mm
(c) $800$ m and $80 \times 10^2$ m  
(d) $100\mu$m and $1$ mm

1.2 Red light has a wavelength of 7000 Å. In μm it is
(a) $0.7\mu$m  
(b) $7\mu$m
(c) $70\mu$m  
(d) $0.07\mu$m

1.3 A speck of dust weighs $1.6 \times 10^{-10}$ kg. How many such particles would weigh $1.6$ kg?
(a) $10^{-10}$  
(b) $10^{10}$
(c) $10$  
(d) $10^{-1}$

1.4 The force acting on a particle is found to be proportional to velocity. The constant of proportionality is measured in terms of
(a) kg s$^{-1}$  
(b) kg s
(c) kg m s$^{-1}$  
(d) kg m s$^{-2}$

1.5 The number of significant digits in $0.0006032$ is
(a) 8  
(b) 7
(c) 4  
(d) 2

1.6 The length of a body is measured as $3.51$ m. If the accuracy is $0.01$ m, then the percentage error in the measurement is
(a) $351\%$  
(b) $1\%$
(c) $0.28\%$  
(d) $0.035\%$

1.7 The dimensional formula for gravitational constant is
(a) $M^1L^3T^{-2}$  
(b) $M^{-1}L^3T^{-2}$
(c) $M^{-1}L^{-3}T^{-2}$  
(d) $M^1L^{-3}T^2$
1.8 The velocity of a body is expressed as \( v = \left( \frac{x}{t} \right) + yt \). The dimensional formula for \( x \) is

(a) \( ML^0T^0 \)  
(b) \( M^0LT^0 \)  
(c) \( M^0LT \)  
(d) \( ML^0T^0 \)

1.9 The dimensional formula for Planck’s constant is

(a) \( MLT \)  
(b) \( ML^3T^{-2} \)  
(c) \( ML^0T^4 \)  
(d) \( ML^2T^{-1} \)

1.10 \( \text{________}_{} \) have the same dimensional formula

(a) Force and momentum  
(b) Stress and strain  
(c) Density and linear density  
(d) Work and potential energy

1.11 What is the role of Physics in technology?

1.12 Write a note on the basic forces in nature.

1.13 Distinguish between fundamental units and derived units.

1.14 Give the SI standard for (i) length (ii) mass and (iii) time.

1.15 Why SI system is considered superior to other systems?

1.16 Give the rules and conventions followed while writing SI units.

1.17 What is the need for measurement of physical quantities?

1.18 You are given a wire and a metre scale. How will you estimate the diameter of the wire?

1.19 Name four units to measure extremely small distances.

1.20 What are random errors? How can we minimise these errors?

1.21 Show that \( \frac{1}{2} gt^2 \) has the same dimensions of distance.

1.22 What are the limitations of dimensional analysis?

1.23 What are the uses of dimensional analysis? Explain with one example.

**Problems**

1.24 How many astronomical units are there in 1 metre?
1.25 If mass of an electron is $9.11 \times 10^{-31}$ kg how many electrons would weigh 1 kg?

1.26 In a submarine fitted with a SONAR, the time delay between generation of a signal and reception of its echo after reflection from an enemy ship is observed to be 73.0 seconds. If the speed of sound in water is $1450$ m s$^{-1}$, then calculate the distance of the enemy ship.

1.27 State the number of significant figures in the following:
   (i) 600900  (ii) 5212.0  (iii) 6.320  (iv) 0.0631  (v) $2.64 \times 10^{24}$

1.28 Find the value of $\pi^2$ correct to significant figures, if $\pi = 3.14$.

1.29 5.74 g of a substance occupies a volume of 1.2 cm$^3$. Calculate its density applying the principle of significant figures.

1.30 The length, breadth and thickness of a rectangular plate are 4.234 m, 1.005 m and 2.01 cm respectively. Find the total area and volume of the plate to correct significant figures.

1.31 The length of a rod is measured as 25.0 cm using a scale having an accuracy of 0.1 cm. Determine the percentage error in length.

1.32 Obtain by dimensional analysis an expression for the surface tension of a liquid rising in a capillary tube. Assume that the surface tension $T$ depends on mass $m$ of the liquid, pressure $P$ of the liquid and radius $r$ of the capillary tube (Take the constant $k = \frac{1}{2}$).

1.33 The force $F$ acting on a body moving in a circular path depends on mass $m$ of the body, velocity $v$ and radius $r$ of the circular path. Obtain an expression for the force by dimensional analysis (Take the value of $k = 1$).

1.34 Check the correctness of the following equation by dimensional analysis

   (i) $F = \frac{mv^2}{r^2}$ where $F$ is force, $m$ is mass, $v$ is velocity and $r$ is radius

   (ii) $n = \frac{1}{2\pi} \sqrt{\frac{g}{l}}$ where $n$ is frequency, $g$ is acceleration due to gravity and $l$ is length.
\[ \frac{1}{2} m v^2 = m g h^2 \] where \( m \) is mass, \( v \) is velocity, \( g \) is acceleration due to gravity and \( h \) is height.

1.35 Convert using dimensional analysis

(i) \( \frac{18}{5} \) kmph into m s\(^{-1}\)

(ii) \( \frac{5}{18} \) m s\(^{-1}\) into kmph

(iii) 13.6 g cm\(^{-3}\) into kg m\(^{-3}\)

Answers

1.1 (a) 1.2 (a) 1.3 (b) 1.4 (a)

1.5 (c) 1.6 (c) 1.7 (b) 1.8 (b)

1.9 (d) 1.10 (d)

1.24 6.68 \times 10^{-12} \text{ AU} 1.25 1.097 \times 10^{30}

1.26 52.925 km 1.27 4, 5, 4, 3, 3

1.28 9.86 1.29 4.8 g cm\(^{-3}\)

1.30 4.255 m\(^2\), 0.0855 m\(^3\) 1.31 0.4 \%

1.32 \( T = \frac{p r}{2} \) 1.33 \( F = \frac{m v^2}{r} \)

1.34 wrong, correct, wrong

1.35 1 m s\(^{-1}\), 1 kmph, 1.36 \times 10^4 \text{ kg m}^{-3}
2. Kinematics

Mechanics is one of the oldest branches of physics. It deals with the study of particles or bodies when they are at rest or in motion. Modern research and development in the spacecraft design, its automatic control, engine performance, electrical machines are highly dependent upon the basic principles of mechanics. Mechanics can be divided into statics and dynamics.

Statics is the study of objects at rest; this requires the idea of forces in equilibrium.

Dynamics is the study of moving objects. It comes from the Greek word dynamis which means power. Dynamics is further subdivided into kinematics and kinetics.

Kinematics is the study of the relationship between displacement, velocity, acceleration and time of a given motion, without considering the forces that cause the motion.

Kinetics deals with the relationship between the motion of bodies and forces acting on them.

We shall now discuss the various fundamental definitions in kinematics.

Particle

A particle is ideally just a piece or a quantity of matter, having practically no linear dimensions but only a position.

Rest and Motion

When a body does not change its position with respect to time, then it is said to be at rest.

Motion is the change of position of an object with respect to time. To study the motion of the object, one has to study the change in position \((x,y,z)\) coordinates of the object with respect to the surroundings. It may be noted that the position of the object changes even due to the change in one, two or all the three coordinates of the position of the
objects with respect to time. Thus motion can be classified into three types:

(i) Motion in one dimension

Motion of an object is said to be one dimensional, if only one of the three coordinates specifying the position of the object changes with respect to time. Example: An ant moving in a straight line, running athlete, etc.

(ii) Motion in two dimensions

In this type, the motion is represented by any two of the three coordinates. Example: a body moving in a plane.

(iii) Motion in three dimensions

Motion of a body is said to be three dimensional, if all the three coordinates of the position of the body change with respect to time.

Examples: motion of a flying bird, motion of a kite in the sky, motion of a molecule, etc.

2.1 Motion in one dimension (rectilinear motion)

The motion along a straight line is known as rectilinear motion. The important parameters required to study the motion along a straight line are position, displacement, velocity, and acceleration.

2.1.1 Position, displacement and distance travelled by the particle

The motion of a particle can be described if its position is known continuously with respect to time.

The total length of the path is the distance travelled by the particle and the shortest distance between the initial and final position of the particle is the displacement.

The distance travelled by a particle, however, is different from its displacement from the origin. For example, if the particle moves from a point O to position P₁ and then to
position $P_2$, its displacement at the position $P_2$ is $-x_2$ from the origin but, the distance travelled by the particle is $x_1 + x_1 + x_2 = (2x_1 + x_2)$ (Fig 2.1).

The distance travelled is a scalar quantity and the displacement is a vector quantity.

2.1.2 Speed and velocity

Speed

It is the distance travelled in unit time. It is a scalar quantity.

Velocity

The velocity of a particle is defined as the rate of change of displacement of the particle. It is also defined as the speed of the particle in a given direction. The velocity is a vector quantity. It has both magnitude and direction.

$$\text{Velocity} = \frac{\text{displacement}}{\text{time taken}}$$

Its unit is m s$^{-1}$ and its dimensional formula is $LT^{-1}$.

Uniform velocity

A particle is said to move with uniform velocity if it moves along a fixed direction and covers equal displacements in equal intervals of time, however small these intervals of time may be.

In a displacement - time graph, (Fig. 2.2) the slope is constant at all the points, when the particle moves with uniform velocity.

Non uniform or variable velocity

The velocity is variable (non-uniform), if it covers unequal displacements in equal intervals of time or if the direction of motion changes or if both the rate of motion and the direction change.
**Average velocity**

Let $s_1$ be the displacement of a body in time $t_1$ and $s_2$ be its displacement in time $t_2$ (Fig. 2.3). The average velocity during the time interval $(t_2 - t_1)$ is defined as

$$v_{\text{average}} = \frac{\text{change in displacement}}{\text{change in time}} = \frac{s_2 - s_1}{t_2 - t_1} = \frac{\Delta s}{\Delta t}$$

From the graph, it is found that the slope of the curve varies.

**Instantaneous velocity**

It is the velocity at any given instant of time or at any given point of its path. The instantaneous velocity $v$ is given by

$$v = \lim_{\Delta t \to 0} \frac{s(t + \Delta t) - s(t)}{\Delta t} = \frac{ds}{dt}$$

**2.1.3 Acceleration**

*If the magnitude or the direction or both of the velocity changes with respect to time, the particle is said to be under acceleration.*

Acceleration of a particle is defined as the rate of change of velocity. Acceleration is a vector quantity.

$$\text{Acceleration} = \frac{\text{change in velocity}}{\text{time taken}}$$

If $u$ is the initial velocity and $v$, the final velocity of the particle after a time $t$, then the acceleration,

$$a = \frac{v - u}{t}$$

Its unit is m s$^{-2}$ and its dimensional formula is LT$^{-2}$.

The instantaneous acceleration is,

$$a = \frac{dv}{dt} = \frac{d}{dt} \left( \frac{ds}{dt} \right) = \frac{d^2s}{dt^2}$$

**Uniform acceleration**

If the velocity changes by an equal amount in equal intervals of time, however small these intervals of time may be, the acceleration is said to be uniform.
Retardation or deceleration

If the velocity decreases with time, the acceleration is negative. The negative acceleration is called retardation or deceleration.

Uniform motion

A particle is in uniform motion when it moves with constant velocity (i.e) zero acceleration.

2.1.4 Graphical representations

The graphs provide a convenient method to present pictorially, the basic informations about a variety of events. Line graphs are used to show the relation of one quantity say displacement or velocity with another quantity such as time.

If the displacement, velocity and acceleration of a particle are plotted with respect to time, they are known as,

(i) displacement – time graph (s - t graph)
(ii) velocity – time graph (v - t graph)
(iii) acceleration – time graph (a - t graph)

Displacement – time graph

When the displacement of the particle is plotted as a function of time, it is displacement - time graph.

As $v = \frac{ds}{dt}$, the slope of the s - t graph at any instant gives the velocity of the particle at that instant. In Fig. 2.4 the particle at time $t_1$, has a positive velocity, at time $t_2$, has zero velocity and at time $t_3$, has negative velocity.

Velocity – time graph

When the velocity of the particle is plotted as a function of time, it is velocity-time graph.

As $a = \frac{dv}{dt}$, the slope of the $v - t$ curve at any instant gives the
acceleration of the particle (Fig. 2.5).

But, \( v = \frac{ds}{dt} \) or \( ds = v \, dt \)

If the displacements are \( s_1 \) and \( s_2 \) in times \( t_1 \) and \( t_2 \), then

\[
\int_{s_1}^{s_2} ds = \int_{t_1}^{t_2} v \, dt
\]

\( s_2 - s_1 = \int_{t_1}^{t_2} v \, dt = \text{area ABCD} \)

The area under the \( v - t \) curve, between the given intervals of time, gives the change in displacement or the distance travelled by the particle during the same interval.

**Acceleration - time graph**

When the acceleration is plotted as a function of time, it is acceleration - time graph (Fig. 2.6).

\[ a = \frac{dv}{dt} \quad \text{(or)} \quad dv = a \, dt \]

If the velocities are \( v_1 \) and \( v_2 \) at times \( t_1 \) and \( t_2 \) respectively, then

\[
\int_{v_1}^{v_2} dv = \int_{t_1}^{t_2} a \, dt \quad \text{(or)} \quad v_2 - v_1 = \int_{t_1}^{t_2} a \, dt = \text{area PQRS} \]

The area under the \( a - t \) curve, between the given intervals of time, gives the change in velocity of the particle during the same interval. If the graph is parallel to the time axis, the body moves with constant acceleration.

**2.1.5 Equations of motion**

For uniformly accelerated motion, some simple equations that relate displacement \( s \), time \( t \), initial velocity \( u \), final velocity \( v \) and acceleration \( a \) are obtained.

(i) As acceleration of the body at any instant is given by the first derivative of the velocity with respect to time,
\[ a = \frac{dv}{dt} \quad \text{(or)} \quad dv = a \, dt \]

If the velocity of the body changes from \( u \) to \( v \) in time \( t \) then from the above equation,

\[
\int_{u}^{v} dv = \int_{0}^{t} a \, dt \quad \Rightarrow \quad [v]_{u}^{v} = a [t]_{0}^{t}
\]

\[
\therefore \quad v - u = at \quad \text{(or)} \quad v = u + at \quad \ldots(1)
\]

(ii) The velocity of the body is given by the first derivative of the displacement with respect to time.

\[ v = \frac{ds}{dt} \quad \text{(or)} \quad ds = v \, dt \]

Since \( v = u + at \), \( ds = (u + at) \, dt \)

The distance \( s \) covered in time \( t \) is,

\[
\int_{0}^{s} ds = \int_{0}^{t} u \, dt + \int_{0}^{t} at \, dt \quad \text{(or)} \quad s = ut + \frac{1}{2}at^2 \quad \ldots(2)
\]

(iii) The acceleration is given by the first derivative of velocity with respect to time. \((i.e)\)

\[ a = \frac{dv}{dt} = \frac{dv}{ds} \cdot \frac{ds}{dt} = \frac{dv}{ds} \cdot v \quad \text{(or)} \quad ds = \frac{1}{a} \cdot v \, dv \]

Therefore,

\[
\int_{0}^{s} ds = \int_{u}^{v} \frac{v \, dv}{a} \quad \text{(i.e)} \quad s = \frac{1}{a} \left[ \frac{v^2}{2} - \frac{u^2}{2} \right]
\]

\[ s = \frac{1}{2a} \left( v^2 - u^2 \right) \quad \text{(or)} \quad 2as = (v^2 - u^2) \]

\[ \therefore \quad v^2 = u^2 + 2as \quad \ldots(3) \]

The equations (1), (2) and (3) are called equations of motion.

**Expression for the distance travelled in \( n^{th} \) second**

Let a body move with an initial velocity \( u \) and travel along a straight line with uniform acceleration \( a \).

Distance travelled in the \( n^{th} \) second of motion is,

\[ s_{n} = \text{distance travelled during first } n \text{ seconds} - \text{distance travelled during (} n-1 \text{) seconds} \]
Distance travelled during \( n \) seconds
\[
D_n = u n + \frac{1}{2} a n^2
\]

Distance travelled during \((n-1)\) seconds
\[
D_{(n-1)} = u(n-1) + \frac{1}{2} a(n-1)^2
\]

\[
\therefore \text{the distance travelled in the } n^{th} \text{ second} = D_n - D_{(n-1)}
\]

\[
(i.e) \ s_n = \left( u n + \frac{1}{2} a n^2 \right) - \left[ u(n-1) + \frac{1}{2} a(n-1)^2 \right]
\]

\[
s_n = u + a \left( n - \frac{1}{2} \right)
\]

\[
s_n = u + \frac{1}{2} a(2n - 1)
\]

**Special Cases**

**Case (i) : For downward motion**

For a particle moving downwards, \( a = g \), since the particle moves in the direction of gravity.

**Case (ii) : For a freely falling body**

For a freely falling body, \( a = g \) and \( u = 0 \), since it starts from rest.

**Case (iii) : For upward motion**

For a particle moving upwards, \( a = -g \), since the particle moves against the gravity.

**2.2 Scalar and vector quantities**

A study of motion will involve the introduction of a variety of quantities, which are used to describe the physical world. Examples of such quantities are distance, displacement, speed, velocity, acceleration, mass, momentum, energy, work, power etc. All these quantities can be divided into two categories – *scalars* and *vectors*.

*The scalar quantities have magnitude only.* It is denoted by a number and unit. Examples : length, mass, time, speed, work, energy,
temperature etc. Scalars of the same kind can be added, subtracted, multiplied or divided by ordinary laws.

The vector quantities have both magnitude and direction. Examples: displacement, velocity, acceleration, force, weight, momentum, etc.

2.2.1 Representation of a vector

Vector quantities are often represented by a scaled vector diagrams. Vector diagrams represent a vector by the use of an arrow drawn to scale in a specific direction. An example of a scaled vector diagram is shown in Fig 2.7.

From the figure, it is clear that

(i) The scale is listed.

(ii) A line with an arrow is drawn in a specified direction.

(iii) The magnitude and direction of the vector are clearly labelled. In the above case, the diagram shows that the magnitude is 4 N and direction is $30^\circ$ to x-axis. The length of the line gives the magnitude and arrow head gives the direction. In notation, the vector is denoted in bold face letter such as $\mathbf{A}$ or with an arrow above the letter as $\vec{A}$, read as vector A or A vector while its magnitude is denoted by $|\mathbf{A}|$ or by $|\vec{A}|$.

2.2.2 Different types of vectors

(i) Equal vectors

Two vectors are said to be equal if they have the same magnitude and same direction, wherever be their initial positions. In Fig. 2.8 the vectors $\vec{A}$ and $\vec{B}$ have the same magnitude and direction. Therefore $\vec{A}$ and $\vec{B}$ are equal vectors.
(ii) **Like vectors**

Two vectors are said to be like vectors, if they have same direction but different magnitudes as shown in Fig. 2.9.

(iii) **Opposite vectors**

The vectors of same magnitude but opposite in direction, are called opposite vectors (Fig. 2.10).

(iv) **Unlike vectors**

The vectors of different magnitude acting in opposite directions are called unlike vectors. In Fig. 2.11 the vectors \( \vec{A} \) and \( \vec{B} \) are unlike vectors.

(v) **Unit vector**

A vector having unit magnitude is called a unit vector. It is also defined as a vector divided by its own magnitude. A unit vector in the direction of a vector \( \vec{A} \) is written as \( \hat{A} \) and is read as ‘A cap’ or ‘A caret’ or ‘A hat’. Therefore,

\[
\hat{A} = \frac{\vec{A}}{|\vec{A}|} \quad \text{(or)} \quad \vec{A} = \hat{A} |\vec{A}|\
\]

Thus, a vector can be written as the product of its magnitude and unit vector along its direction.

**Orthogonal unit vectors**

There are three most common unit vectors in the positive directions of X,Y and Z axes of Cartesian coordinate system, denoted by \( i, j \) and \( k \) respectively. Since they are along the mutually perpendicular directions, they are called orthogonal unit vectors.

(vi) **Null vector or zero vector**

A vector whose magnitude is zero, is called a null vector or zero vector. It is represented by \( \vec{0} \) and its starting and end points are the same. The direction of null vector is not known.
(vii) **Proper vector**

All the non-zero vectors are called proper vectors.

(viii) **Co-initial vectors**

Vectors having the same starting point are called co-initial vectors. In Fig. 2.12, \( \overrightarrow{A} \) and \( \overrightarrow{B} \) start from the same origin O. Hence, they are called as co-initial vectors.

(ix) **Coplanar vectors**

Vectors lying in the same plane are called coplanar vectors and the plane in which the vectors lie are called plane of vectors.

2.2.3 **Addition of vectors**

As vectors have both magnitude and direction they cannot be added by the method of ordinary algebra.

Vectors can be added graphically or geometrically. We shall now discuss the addition of two vectors graphically using head to tail method.

Consider two vectors \( \overrightarrow{P} \) and \( \overrightarrow{Q} \) which are acting along the same line. To add these two vectors, join the tail of \( \overrightarrow{Q} \) with the head of \( \overrightarrow{P} \) (Fig. 2.13).

The resultant of \( \overrightarrow{P} \) and \( \overrightarrow{Q} \) is \( \overrightarrow{R} = \overrightarrow{P} + \overrightarrow{Q} \). The length of the line AD gives the magnitude of \( \overrightarrow{R} \). \( \overrightarrow{R} \) acts in the same direction as that of \( \overrightarrow{P} \) and \( \overrightarrow{Q} \).

In order to find the sum of two vectors, which are inclined to each other, triangle law of vectors or parallelogram law of vectors, can be used.

(i) **Triangle law of vectors**

*If two vectors are represented in magnitude and direction by the two adjacent sides of a triangle taken in order, then their resultant is the closing side of the triangle taken in the reverse order.*
To find the resultant of two vectors $\vec{P}$ and $\vec{Q}$ which are acting at an angle $\theta$, the following procedure is adopted.

First draw $\overrightarrow{OA} = \vec{P}$ (Fig. 2.14) Then starting from the arrow head of $\vec{P}$, draw the vector $\overrightarrow{AB} = \vec{Q}$. Finally, draw a vector $\overrightarrow{OB} = \vec{R}$ from the tail of vector $\vec{P}$ to the head of vector $\vec{Q}$. Vector $\overrightarrow{OB} = \vec{R}$ is the sum of the vectors $\vec{P}$ and $\vec{Q}$. Thus $\vec{R} = \vec{P} + \vec{Q}$.

The magnitude of $\vec{P} + \vec{Q}$ is determined by measuring the length of $\vec{R}$ and direction by measuring the angle between $\vec{P}$ and $\vec{R}$.

The magnitude and direction of $\vec{R}$, can be obtained by using the sine law and cosine law of triangles. Let $\alpha$ be the angle made by the resultant $\vec{R}$ with $\vec{P}$. The magnitude of $\vec{R}$ is,

$$R^2 = P^2 + Q^2 - 2PQ \cos (180^\circ - \theta)$$

$$R = \sqrt{P^2 + Q^2 + 2PQ \cos \theta}$$

The direction of $R$ can be obtained by,

$$\frac{P}{\sin \beta} = \frac{Q}{\sin \alpha} = \frac{R}{\sin (180^\circ - \theta)}$$

(ii) Parallelogram law of vectors

If two vectors acting at a point are represented in magnitude and direction by the two adjacent sides of a parallelogram, then their resultant is represented in magnitude and direction by the diagonal passing through the common tail of the two vectors.

Let us consider two vectors $\vec{P}$ and $\vec{Q}$ which are inclined to each other at an angle $\theta$ as shown in Fig. 2.15. Let the vectors $\vec{P}$ and $\vec{Q}$ be represented in magnitude and direction by the two sides $OA$ and $OB$ of a parallelogram $OACB$. The diagonal $OC$ passing through the common tail $O$, gives the magnitude and direction of the resultant $\vec{R}$.

$CD$ is drawn perpendicular to the extended $OA$, from $C$. Let $|COD|$ made by $\vec{R}$ with $\vec{P}$ be $\alpha$. 

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From right angled triangle $OCD$,
\[OC^2 = OD^2 + CD^2\]
\[= (OA + AD)^2 + CD^2\]
\[= OA^2 + AD^2 + 2OA \cdot AD + CD^2\]  \hspace{1cm} \ldots(1)

In Fig. 2.15  \hspace{1cm} |BOA = \theta = |CAD|

From right angled $\Delta CAD$,
\[AC^2 = AD^2 + CD^2\]  \hspace{1cm} \ldots(2)

Substituting (2) in (1)
\[OC^2 = OA^2 + AC^2 + 2OA \cdot AD\]  \hspace{1cm} \ldots(3)

From $\Delta ACD$,
\[CD = AC \sin \theta\]  \hspace{1cm} \ldots(4)

\[AD = AC \cos \theta\]  \hspace{1cm} \ldots(5)

Substituting (5) in (3)
\[OC^2 = OA^2 + AC^2 + 2OA \cdot AC \cos \theta\]

Substituting $OC = R$, $OA = P$, $OB = AC = Q$ in the above equation
\[R^2 = P^2 + Q^2 + 2PQ \cos \theta\]

(or) \[R = \sqrt{P^2 + Q^2 + 2PQ \cos \theta}\]  \hspace{1cm} \ldots(6)

Equation (6) gives the magnitude of the resultant. From $\Delta OCD$,
\[
\tan \alpha = \frac{CD}{OD} = \frac{CD}{OA + AD}
\]

Substituting (4) and (5) in the above equation,
\[
\tan \alpha = \frac{AC \sin \theta}{OA + AC \cos \theta} = \frac{Q \sin \theta}{P + Q \cos \theta}
\]

(or) \[\alpha = \tan^{-1}\left[\frac{Q \sin \theta}{P + Q \cos \theta}\right]\]  \hspace{1cm} \ldots(7)

Equation (7) gives the direction of the resultant.

**Special Cases**

(i) When two vectors act in the same direction

In this case, the angle between the two vectors $\theta = 0^o$, \[\cos 0^o = 1, \sin 0^o = 0\]
Thus, the resultant vector acts in the same direction as the individual vectors and is equal to the sum of the magnitude of the two vectors.

(ii) When two vectors act in the opposite direction

In this case, the angle between the two vectors $\theta = 180^\circ$, $\cos 180^\circ = -1$, $\sin 180^\circ = 0$.

From (6) $R = \sqrt{P^2 + Q^2 - 2PQ} = (P - Q)$

From (7) $\alpha = \tan^{-1} \left( \frac{0}{P - Q} \right) = \tan^{-1}(0) = 0$

Thus, the resultant vector has a magnitude equal to the difference in magnitude of the two vectors and acts in the direction of the bigger of the two vectors

(iii) When two vectors are at right angles to each other

In this case, $\theta = 90^\circ$, $\cos 90^\circ = 0$, $\sin 90^\circ = 1$

From (6) $R = \sqrt{P^2 + Q^2}$

From (7) $\alpha = \tan^{-1} \left( \frac{Q}{P} \right)$

The resultant $\vec{R}$ vector acts at an angle $\alpha$ with vector $\vec{P}$.

2.2.4 Subtraction of vectors

The subtraction of a vector from another is equivalent to the addition of one vector to the negative of the other.

For example $\vec{Q} - \vec{P} = \vec{Q} + (-\vec{P})$.

Thus to subtract $\vec{P}$ from $\vec{Q}$, one has to add $-\vec{P}$ with $\vec{Q}$ (Fig 2.16a). Therefore, to subtract $\vec{P}$ from $\vec{Q}$, reversed $\vec{P}$ is added to the
\( \vec{Q} \). For this, first draw \( \overrightarrow{AB} = \vec{Q} \) and then starting from the arrow head of \( \vec{Q} \), draw \( \overrightarrow{BC} = (-\vec{P}) \) and finally join the head of \( -\vec{P} \). Vector \( \vec{R} \) is the sum of \( \vec{Q} \) and \( -\vec{P} \). (i.e) difference \( \vec{Q} - \vec{P} \).

\[ \text{Fig 2.16 Subtraction of vectors} \]

The resultant of two vectors which are antiparallel to each other is obtained by subtracting the smaller vector from the bigger vector as shown in Fig 2.16b. The direction of the resultant vector is in the direction of the bigger vector.

### 2.2.5 Product of a vector and a scalar

Multiplication of a scalar and a vector gives a vector quantity which acts along the direction of the vector.

**Examples**

(i) If \( \vec{a} \) is the acceleration produced by a particle of mass \( m \) under the influence of the force, then \( \vec{F} = m\vec{a} \).

(ii) momentum = mass \times velocity (i.e) \( \vec{P} = m\vec{v} \).

### 2.2.6 Resolution of vectors and rectangular components

A vector directed at an angle with the co-ordinate axis, can be resolved into its components along the axes. This process of splitting a vector into its components is known as resolution of a vector.

Consider a vector \( \vec{R} = \overrightarrow{OA} \) making an angle \( \theta \) with X - axis. The vector \( \vec{R} \) can be resolved into two components along X - axis and Y-axis respectively. Draw two perpendiculars from \( A \) to X and Y axes respectively. The intercepts on these axes are called the scalar components \( R_x \) and \( R_y \).
Then, OP is $R_x$, which is the magnitude of $x$ component of $\vec{R}$ and OQ is $R_y$, which is the magnitude of $y$ component of $\vec{R}$.

From $\triangle OPA$,

\[
\cos \theta = \frac{OP}{OA} = \frac{R_x}{R} \quad \text{(or)} \quad R_x = R \cos \theta
\]

\[
\sin \theta = \frac{OQ}{OA} = \frac{R_y}{R} \quad \text{(or)} \quad R_y = R \sin \theta
\]

and

\[
R^2 = R_x^2 + R_y^2
\]

Also, $\vec{R}$ can be expressed as $\vec{R} = R_x \hat{i} + R_y \hat{j}$ where $\hat{i}$ and $\hat{j}$ are unit vectors.

In terms of $R_x$ and $R_y$, $\theta$ can be expressed as $\theta = \tan^{-1} \left[ \frac{R_y}{R_x} \right]$

### 2.2.7 Multiplication of two vectors

Multiplication of a vector by another vector does not follow the laws of ordinary algebra. There are two types of vector multiplication:

(i) **Scalar product and (ii) Vector product.**

#### (i) Scalar product or Dot product of two vectors

If the product of two vectors is a scalar, then it is called scalar product. If $\vec{A}$ and $\vec{B}$ are two vectors, then their scalar product is written as $\vec{A} \cdot \vec{B}$ and read as $\vec{A}$ dot $\vec{B}$. Hence scalar product is also called dot product. This is also referred as inner or direct product.

The scalar product of two vectors is a scalar, which is equal to the product of magnitudes of the two vectors and the cosine of the angle between them. The scalar product of two vectors $\vec{A}$ and $\vec{B}$ may be expressed as $\vec{A} \cdot \vec{B} = |\vec{A}| |\vec{B}| \cos \theta$ where $|\vec{A}|$ and $|\vec{B}|$ are the magnitudes of $\vec{A}$ and $\vec{B}$ respectively and $\theta$ is the angle between $\vec{A}$ and $\vec{B}$ as shown in Fig 2.18.
(ii) **Vector product or Cross product of two vectors**

If the product of two vectors is a vector, then it is called vector product. If \(\vec{A}\) and \(\vec{B}\) are two vectors then their vector product is written as \(\vec{A} \times \vec{B}\) and read as \(\vec{A}\) cross \(\vec{B}\). This is also referred as outer product.

The vector product or cross product of two vectors is a vector whose magnitude is equal to the product of their magnitudes and the sine of the smaller angle between them and the direction is perpendicular to a plane containing the two vectors.

If \(\theta\) is the smaller angle through which \(\vec{A}\) should be rotated to reach \(\vec{B}\), then the cross product of \(\vec{A}\) and \(\vec{B}\) (Fig. 2.19) is expressed as,

\[
\vec{A} \times \vec{B} = |\vec{A}| \cdot |\vec{B}| \sin \theta \hat{n} = \vec{C}
\]

where \(|\vec{A}|\) and \(|\vec{B}|\) are the magnitudes of \(\vec{A}\) and \(\vec{B}\) respectively. \(\vec{C}\) is perpendicular to the plane containing \(\vec{A}\) and \(\vec{B}\). The direction of \(\vec{C}\) is along the direction in which the tip of a screw moves when it is rotated from \(\vec{A}\) to \(\vec{B}\). Hence \(\vec{C}\) acts along OC. By the same argument, \(\vec{B} \times \vec{A}\) acts along OD.

**2.3 Projectile motion**

A body thrown with some initial velocity and then allowed to move under the action of gravity alone, is known as a projectile.

If we observe the path of the projectile, we find that the projectile moves in a path, which can be considered as a part of parabola. Such a motion is known as projectile motion.

A few examples of projectiles are (i) a bomb thrown from an aeroplane (ii) a javelin or a shot-put thrown by an athlete (iii) motion of a ball hit by a cricket bat etc.

The different types of projectiles are shown in Fig. 2.20. A body can be projected in two ways:
(i) It can be projected horizontally from a certain height.

(ii) It can be thrown from the ground in a direction inclined to it.

The projectiles undergo a vertical motion as well as horizontal motion. The two components of the projectile motion are (i) vertical component and (ii) horizontal component. These two perpendicular components of motion are independent of each other.

A body projected with an initial velocity making an angle with the horizontal direction possesses uniform horizontal velocity and variable vertical velocity, due to force of gravity. The object therefore has horizontal and vertical motions simultaneously. The resultant motion would be the vector sum of these two motions and the path following would be curvilinear.

The above discussion can be summarised as in the Table 2.1

<table>
<thead>
<tr>
<th>Motion</th>
<th>Forces</th>
<th>Velocity</th>
<th>Acceleration</th>
</tr>
</thead>
<tbody>
<tr>
<td>Horizontal</td>
<td>No force acts</td>
<td>Constant</td>
<td>Zero</td>
</tr>
<tr>
<td>Vertical</td>
<td>The force of gravity acts</td>
<td>Changes (~10 m s⁻¹)</td>
<td>Downwards (~10 m s⁻²)</td>
</tr>
</tbody>
</table>

In the study of projectile motion, it is assumed that the air resistance is negligible and the acceleration due to gravity remains constant.
**Angle of projection**

The angle between the initial direction of projection and the horizontal direction through the point of projection is called the angle of projection.

**Velocity of projection**

The velocity with which the body is projected is known as velocity of projection.

**Range**

Range of a projectile is the horizontal distance between the point of projection and the point where the projectile hits the ground.

**Trajectory**

The path described by the projectile is called the trajectory.

**Time of flight**

Time of flight is the total time taken by the projectile from the instant of projection till it strikes the ground.

2.3.1 **Motion of a projectile thrown horizontally**

Let us consider an object thrown horizontally with a velocity \( u \) from a point A, which is at a height \( h \) from the horizontal plane \( OX \) (Fig 2.21). The object acquires the following motions simultaneously:

(i) Uniform velocity with which it is projected in the horizontal direction \( OX \)

(ii) Vertical velocity, which is non-uniform due to acceleration due to gravity.

The two velocities are independent of each other. The horizontal velocity of the object shall remain constant as no acceleration is acting in the horizontal direction. The velocity in the vertical direction shall go on changing because of acceleration due to gravity.
Path of a projectile

Let the time taken by the object to reach C from A = t

Vertical distance travelled by the object in time \( t = s = y \)

From equation of motion, \( s = u_1 t + \frac{1}{2} a t^2 \) ...(1)

Substituting the known values in equation (1),

\[ y = (0) t + \frac{1}{2} g t^2 = \frac{1}{2} g t^2 \] ...(2)

At A, the initial velocity in the horizontal direction is \( u \).

Horizontal distance travelled by the object in time \( t \) is \( x \).

\[ \therefore x = \text{horizontal velocity} \times \text{time} = u t \quad \text{(or) } t = \frac{x}{u} \quad \ldots(3) \]

Substituting \( t \) in equation (2),

\[ y = \frac{1}{2} g \left( \frac{x}{u} \right)^2 = \frac{1}{2} g \frac{x^2}{u^2} \quad \ldots(4) \]

(or) \[ y = kx^2 \]

where \( k = \frac{g}{2u^2} \) is a constant.

The above equation is the equation of a parabola. Thus the path taken by the projectile is a parabola.

Resultant velocity at C

At an instant of time \( t \), let the body be at C.

At A, initial vertical velocity \( (u_1) = 0 \)

At C, the horizontal velocity \( (u_2) = u \)

At C, the vertical velocity = \( u_2 \)

From equation of motion, \( u_2 = u_1 + g t \)

Substituting all the known values, \( u_2 = 0 + g t \) \ldots(5)

The resultant velocity at C is \( v = \sqrt{u_x^2 + u_2^2} = \sqrt{u^2 + g^2 t^2} \) \ldots(6)

The direction of \( v \) is given by \[ \tan \theta = \frac{u_2}{u_x} = \frac{gt}{u} \] \ldots(7)

where \( \theta \) is the angle made by \( v \) with X axis.
**Time of flight and range**

The distance $OB = R$, is called as range of the projectile.

Range = horizontal velocity $\times$ time taken to reach the ground

\[ R = u \cdot t_f \] \hspace{1cm} ...(8)

where $t_f$ is the time of flight

At A, initial vertical velocity $(u_y) = 0$

The vertical distance travelled by the object in time $t_f = h$

From the equations of motion

\[ s_y = u_t t_f + \frac{1}{2} g t_f^2 \] \hspace{1cm} ...(9)

Substituting the known values in equation (9),

\[ h = (0) \cdot t_f^2 + \frac{1}{2} g t_f^2 \] \hspace{1cm} (or) \hspace{1cm} \[ t_f = \sqrt{\frac{2h}{g}} \] \hspace{1cm} ...(10)

Substituting $t_f$ in equation (8), Range

\[ R = u \sqrt{\frac{2h}{g}} \] \hspace{1cm} ...(11)

2.3.2 **Motion of a projectile projected at an angle with the horizontal (oblique projection)**

Consider a body projected from a point O on the surface of the Earth with an initial velocity $u$ at an angle $\theta$ with the horizontal as shown in Fig. 2.23. The velocity $u$ can be resolved into two components

\[ u_x = u \cos \theta \]

\[ u_y = u \sin \theta \]

![Fig 2.23  Motion of a projectile projected at an angle with horizontal](image-url)
(i) $u_x = u \cos \theta$, along the horizontal direction OX and
(ii) $u_y = u \sin \theta$, along the vertical direction OY

The horizontal velocity $u_x$ of the object shall remain constant as no acceleration is acting in the horizontal direction. But the vertical component $u_y$ of the object continuously decreases due to the effect of the gravity and it becomes zero when the body is at the highest point of its path. After this, the vertical component $u_y$ is directed downwards and increases with time till the body strikes the ground at $B$.

**Path of the projectile**

Let $t_1$ be the time taken by the projectile to reach the point C from the instant of projection.

Horizontal distance travelled by the projectile in time $t_1$ is,

$$x = \text{horizontal velocity } \times \text{time}$$

$$x = u \cos \theta \times t_1 \quad \text{(or) } t_1 = \frac{x}{u \cos \theta} \quad \ldots \ldots \text{(1)}$$

Let the vertical distance travelled by the projectile in time $t_1 = s = y$

At O, initial vertical velocity $u_1 = u \sin \theta$

From the equation of motion $s = u_1 t_1 - \frac{1}{2} a t_1^2$

Substituting the known values,

$$y = (u \sin \theta) t_1 - \frac{1}{2} g t_1^2 \quad \ldots \ldots \text{(2)}$$

Substituting equation (1) in equation (2),

$$y = (u \sin \theta) \left( \frac{x}{u \cos \theta} \right) - \frac{1}{2} (g) \left( \frac{x}{u \cos \theta} \right)^2$$

$$y = x \tan \theta - \frac{gx^2}{2u^2 \cos^2 \theta} \quad \ldots \ldots \text{(3)}$$

The above equation is of the form $y = Ax + Bx^2$ and represents a parabola. Thus the path of a projectile is a parabola.

**Resultant velocity of the projectile at any instant $t_1$**

At C, the velocity along the horizontal direction is $u_x = u \cos \theta$ and the velocity along the vertical direction is $u_y = u_2$. 

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From the equation of motion,

\[ u_2 = u_1 - gt_1 \]
\[ u_2 = u \sin \theta - gt_1 \]

:. The resultant velocity at

C is \( v = \sqrt{u_x^2 + u_y^2} \)

\[ v = \sqrt{(u \cos \theta)^2 + (u \sin \theta - gt_1)^2} \]
\[ = \sqrt{u^2 + g^2t_1^2 - 2ut_1gsin\theta} \]

The direction of \( v \) is given by

\[ \tan \alpha = \frac{u_y}{u_x} = \frac{u \sin \theta - gt_1}{u \cos \theta} \quad \text{(or)} \quad \alpha = \tan^{-1} \left( \frac{u \sin \theta - gt_1}{u \cos \theta} \right) \]

where \( \alpha \) is the angle made by \( v \) with the horizontal line.

**Maximum height reached by the projectile**

The maximum vertical displacement produced by the projectile is known as the maximum height reached by the projectile. In Fig 2.23, EA is the maximum height attained by the projectile. It is represented as \( h_{\text{max}} \)

At O, the initial vertical velocity \( (u_y) = u \sin \theta \)

At A, the final vertical velocity \( (u_y) = 0 \)

The vertical distance travelled by the object = \( s_y = h_{\text{max}} \)

From equation of motion, \( u_y^2 = u_y^2 - 2gs_y \)

Substituting the known values, \( (0)^2 - (u \sin \theta)^2 = 2gh_{\text{max}} \)

\[ 2gh_{\text{max}} = u^2 \sin^2 \theta \quad \text{(or)} \quad h_{\text{max}} = \frac{u^2 \sin^2 \theta}{2g} \quad \ldots(4) \]

**Time taken to attain maximum height**

Let \( t \) be the time taken by the projectile to attain its maximum height.

From equation of motion \( u_y = u_i - gt \)
Substituting the known values \( 0 = u \sin \theta - gt \)

\[
t = \frac{u \sin \theta}{g}
\] ... (5)

**Time of flight**

Let \( t_f \) be the time of flight (i.e) the time taken by the projectile to reach \( B \) from \( O \) through \( A \). When the body returns to the ground, the net vertical displacement made by the projectile

\[
s_y = h_{\text{max}} - h_{\text{max}} = 0
\]

From the equation of motion

\[
s_y = u_1 t_f - \frac{1}{2} g t_f^2
\]

Substituting the known values \( 0 = (u \sin \theta) t_f - \frac{1}{2} g t_f^2 \)

\[
\frac{1}{2} g t_f^2 = (u \sin \theta) t_f \quad \text{(or)} \quad t_f = \frac{2u\sin \theta}{g} \quad ... (6)
\]

From equations (5) and (6)

\[
t_f = 2t \quad ... (7)
\]

(i.e) the time of flight is twice the time taken to attain the maximum height.

**Horizontal range**

The horizontal distance \( OB \) is called the range of the projectile.

Horizontal range = horizontal velocity \( \times \) time of flight

(i.e) \( R = u \cos \theta \times t_f \)

Substituting the value of \( t_f \), \( R = (u \cos \theta) \frac{2u\sin \theta}{g} \)

\[
R = \frac{u^2 (2 \sin \theta \cos \theta)}{g}
\]

\[
: R = \frac{u^2 \sin 2\theta}{g} \quad ... (8)
\]

**Maximum Range**

From (8), it is seen that for the given velocity of projection, the horizontal range depends on the angle of projection only. The range is maximum only if the value of \( \sin 2\theta \) is maximum.
For maximum range \( R_{\text{max}} \) \( \sin 2\theta = 1 \)

(i.e) \( \theta = 45^\circ \)

Therefore the range is maximum when the angle of projection is \( 45^\circ \).

\[
R_{\text{max}} = 2u^2 \times \frac{1}{g} \Rightarrow R_{\text{max}} = \frac{u^2}{g}
\] ...(9)

2.4 Newton’s laws of motion

Various philosophers studied the basic ideas of cause of motion. According to Aristotle, a constant external force must be applied continuously to an object in order to keep it moving with uniform velocity. Later this idea was discarded and Galileo gave another idea on the basis of the experiments on an inclined plane. According to him, no force is required to keep an object moving with constant velocity. It is the presence of frictional force that tends to stop moving object. The smaller the frictional force between the object and the surface on which it is moving, the larger the distance it will travel before coming to rest. After Galileo, it was Newton who made a systematic study of motion and extended the ideas of Galileo.

Newton formulated the laws concerning the motion of the object. There are three laws of motion. A deep analysis of these laws lead us to the conclusion that these laws completely define the force. The first law gives the fundamental definition of force; the second law gives the quantitative and dimensional definition of force while the third law explains the nature of the force.

2.4.1 Newton’s first law of motion

It states that every body continues in its state of rest or of uniform motion along a straight line unless it is compelled by an external force to change that state.

This law is based on Galileo’s law of inertia. Newton’s first law of motion deals with the basic property of matter called inertia and the definition of force.

Inertia is that property of a body by virtue of which the body is unable to change its state by itself in the absence of external force.
The inertia is of three types

(i) Inertia of rest
(ii) Inertia of motion
(iii) Inertia of direction.

(i) Inertia of rest

It is the inability of the body to change its state of rest by itself.

Examples

(i) A person standing in a bus falls backward when the bus suddenly starts moving. This is because, the person who is initially at rest continues to be at rest even after the bus has started moving.

(ii) A book lying on the table will remain at rest, until it is moved by some external agencies.

(iii) When a carpet is beaten by a stick, the dust particles fall off vertically downwards once they are released and do not move along the carpet and fall off.

(ii) Inertia of motion

Inertia of motion is the inability of the body to change its state of motion by itself.

Examples

(a) When a passenger gets down from a moving bus, he falls down in the direction of the motion of the bus.

(b) A passenger sitting in a moving car falls forward, when the car stops suddenly.

(c) An athlete running in a race will continue to run even after reaching the finishing point.

(iii) Inertia of direction

It is the inability of the body to change its direction of motion by itself.

Examples

When a bus moving along a straight line takes a turn to the right, the passengers are thrown towards left. This is due to inertia which makes the passengers travel along the same straight line, even though the bus has turned towards the right.
This inability of a body to change by itself its state of rest or of uniform motion along a straight line or direction, is known as inertia. The inertia of a body is directly proportional to the mass of the body.

From the first law, we infer that to change the state of rest or uniform motion, an external agency called, the force is required.

Force is defined as that which when acting on a body changes or tends to change the state of rest or of uniform motion of the body along a straight line.

A force is a push or pull upon an object, resulting the change of state of a body. Whenever there is an interaction between two objects, there is a force acting on each other. When the interaction ceases, the two objects no longer experience a force. Forces exist only as a result of an interaction.

There are two broad categories of forces between the objects, contact forces and non-contact forces resulting from action at a distance.

Contact forces are forces in which the two interacting objects are physically in contact with each other.

Tensile force, normal force, force due to air resistance, applied forces and frictional forces are examples of contact forces.

Action-at-a-distance forces (non-contact forces) are forces in which the two interacting objects are not in physical contact which each other, but are able to exert a push or pull despite the physical separation. Gravitational force, electrical force and magnetic force are examples of non-contact forces.

**Momentum of a body**

It is observed experimentally that the force required to stop a moving object depends on two factors: (i) mass of the body and (ii) its velocity

A body in motion has momentum. The momentum of a body is defined as the product of its mass and velocity. If \( m \) is the mass of the body and \( \vec{v} \) its velocity, the linear momentum of the body is given by \( \vec{p} = m\vec{v} \).

Momentum has both magnitude and direction and it is, therefore, a vector quantity. The momentum is measured in terms of kg m s\(^{-1}\) and its dimensional formula is MLT\(^{-1}\).
When a force acts on a body, its velocity changes, consequently, its momentum also changes. The slowly moving bodies have smaller momentum than fast moving bodies of same mass.

If two bodies of unequal masses and velocities have same momentum, then,
\[ \vec{p}_1 = \vec{p}_2 \]
(i.e) \[ m_1 \vec{v}_1 = m_2 \vec{v}_2 \quad \Rightarrow \quad \frac{m_1}{m_2} = \frac{\vec{v}_2}{\vec{v}_1} \]

Hence for bodies of same momenta, their velocities are inversely proportional to their masses.

**2.4.2 Newton’s second law of motion**

Newton’s first law of motion deals with the behaviour of objects on which all existing forces are balanced. Also, it is clear from the first law of motion that a body in motion needs a force to change the direction of motion or the magnitude of velocity or both. This implies that force is such a physical quantity that causes or tends to cause an acceleration.

Newton’s second law of motion deals with the behaviour of objects on which all existing forces are not balanced.

According to this law, *the rate of change of momentum of a body is directly proportional to the external force applied on it and the change in momentum takes place in the direction of the force.*

If \( \vec{p} \) is the momentum of a body and \( \vec{F} \) the external force acting on it, then according to Newton’s second law of motion,
\[ \vec{F} \propto \frac{d\vec{p}}{dt} \quad \text{(or)} \quad \vec{F} = k \frac{d\vec{p}}{dt} \]
where \( k \) is a proportionality constant.

If a body of mass \( m \) is moving with a velocity \( \vec{v} \) then, its momentum is given by \( \vec{p} = m \vec{v} \).

\[ \therefore \vec{F} = k \frac{d}{dt}(m \vec{v}) = k m \frac{d\vec{v}}{dt} \]

Unit of force is chosen in such a manner that the constant \( k \) is equal to unity. (i.e) \( k = 1 \).
\[ F = m \frac{d\vec{v}}{dt} = m \vec{a} \quad \text{where} \quad \vec{a} = \frac{d\vec{v}}{dt} \]

is the acceleration produced in the motion of the body.

The force acting on a body is measured by the product of mass of the body and acceleration produced by the force acting on the body. The second law of motion gives us a measure of the force.

The acceleration produced in the body depends upon the inertia of the body (i.e) greater the inertia, lesser the acceleration. One newton is defined as that force which, when acting on unit mass produces unit acceleration. Force is a vector quantity. The unit of force is kg m s\(^{-2}\) or newton. Its dimensional formula is ML\(^{-2}\).

**Impulsive force and Impulse of a force**

(i) **Impulsive Force**

An impulsive force is a very great force acting for a very short time on a body, so that the change in the position of the body during the time the force acts on it may be neglected.

(e.g.) The blow of a hammer, the collision of two billiard balls etc.

(ii) **Impulse of a force**

The impulse \( J \) of a constant force \( F \) acting for a time \( t \) is defined as the product of the force and time.

\( \text{(i.e) Impulse} = \text{Force} \times \text{time} \)

\( J = F \times t \)

The impulse of force \( F \) acting over a time interval \( t \) is defined by the integral,

\[ J = \int_{0}^{t} F \, dt \quad \text{...(1)} \]

The impulse of a force, therefore can be visualised as the area under the force versus time graph as shown in Fig. 2.25. When a variable force acting for a short interval of time, then the impulse can be measured as,

\[ J = F_{\text{average}} \times dt \quad \text{...(2)} \]
Impulse of a force is a vector quantity and its unit is N s.

**Principle of impulse and momentum**

By Newton’s second law of motion, the force acting on a body = \( m \ a \) where \( m \) = mass of the body and \( a \) = acceleration produced

The impulse of the force = \( F \times t = (m \ a) \ t \)

If \( u \) and \( v \) be the initial and final velocities of the body then,

\[
a = \frac{(v - u)}{t},
\]

Therefore, impulse of the force = \( m \times \frac{(v - u)}{t} \times t = m(v - u) = mv - mu \)

Impulse = final momentum of the body

- initial momentum of the body.

(i.e) Impulse of the force = Change in momentum

The above equation shows that the total change in the momentum of a body during a time interval is equal to the impulse of the force acting during the same interval of time. This is called principle of impulse and momentum.

**Examples**

(i) A cricket player while catching a ball lowers his hands in the direction of the ball.

If the total change in momentum is brought about in a very short interval of time, the average force is very large according to the equation, \( F = \frac{mv - mu}{t} \)

By increasing the time interval, the average force is decreased. It is for this reason that a cricket player while catching a ball, to increase the time of contact, the player should lower his hand in the direction of the ball, so that he is not hurt.

(ii) A person falling on a cemented floor gets injured more where as a person falling on a sand floor does not get hurt. For the same reason, in wrestling, high jump etc., soft ground is provided.

(iii) The vehicles are fitted with springs and shock absorbers to reduce jerks while moving on uneven or wavy roads.
2.4.3 Newton’s third Law of motion

It is a common observation that when we sit on a chair, our body exerts a downward force on the chair and the chair exerts an upward force on our body. There are two forces resulting from this interaction: a force on the chair and a force on our body. These two forces are called action and reaction forces. Newton’s third law explains the relation between these action forces. It states that for every action, there is an equal and opposite reaction.

(i.e.) whenever one body exerts a certain force on a second body, the second body exerts an equal and opposite force on the first. Newton’s third law is sometimes called as the law of action and reaction.

Let there be two bodies 1 and 2 exerting forces on each other. Let the force exerted on the body 1 by the body 2 be $\vec{F}_{12}$ and the force exerted on the body 2 by the body 1 be $\vec{F}_{21}$. Then according to third law, $\vec{F}_{12} = -\vec{F}_{21}$.

One of these forces, say $\vec{F}_{12}$ may be called as the action whereas the other force $\vec{F}_{21}$ may be called as the reaction or vice versa. This implies that we cannot say which is the cause (action) or which is the effect (reaction). It is to be noted that always the action and reaction do not act on the same body; they always act on different bodies. The action and reaction never cancel each other and the forces always exist in pair.

The effect of third law of motion can be observed in many activities in our everyday life. The examples are

(i) When a bullet is fired from a gun with a certain force (action), there is an equal and opposite force exerted on the gun in the backward direction (reaction).

(ii) When a man jumps from a boat to the shore, the boat moves away from him. The force he exerts on the boat (action) is responsible for its motion and his motion to the shore is due to the force of reaction exerted by the boat on him.

(iii) The swimmer pushes the water in the backward direction with a certain force (action) and the water pushes the swimmer in the forward direction with an equal and opposite force (reaction).
(iv) We will not be able to walk if there were no reaction force. In order to walk, we push our foot against the ground. The Earth in turn exerts an equal and opposite force. This force is inclined to the surface of the Earth. The vertical component of this force balances our weight and the horizontal component enables us to walk forward.

(v) A bird flies by with the help of its wings. The wings of a bird push air downwards (action). In turn, the air reacts by pushing the bird upwards (reaction).

(vi) When a force exerted directly on the wall by pushing the palm of our hand against it (action), the palm is distorted a little because, the wall exerts an equal force on the hand (reaction).

**Law of conservation of momentum**

From the principle of impulse and momentum, impulse of a force, \( J = mv - mu \)

If \( J = 0 \) then \( mv - mu = 0 \) (or) \( mv = mu \)

(i.e) final momentum = initial momentum

In general, *the total momentum of the system is always a constant (i.e) when the impulse due to external forces is zero, the momentum of the system remains constant. This is known as law of conservation of momentum.*

We can prove this law, in the case of a head on collision between two bodies.

**Proof**

Consider a body A of mass \( m_1 \) moving with a velocity \( u_1 \) collides head on with another body B of mass \( m_2 \) moving in the same direction as A with velocity \( u_2 \) as shown in Fig 2.26.
After collision, let the velocities of the bodies be changed to $v_1$ and $v_2$ respectively, and both moves in the same direction. During collision, each body experiences a force.

The force acting on one body is equal in magnitude and opposite in direction to the force acting on the other body. Both forces act for the same interval of time.

Let $F_1$ be force exerted by A on B (action), $F_2$ be force exerted by B on A (reaction) and $t$ be the time of contact of the two bodies during collision.

Now, $F_1$ acting on the body B for a time $t$, changes its velocity from $u_2$ to $v_2$.

∴ $F_1 = \text{mass of the body B} \times \text{acceleration of the body B}$

$$= m_2 \times \frac{(v_2 - u_2)}{t} \quad \ldots(1)$$

Similarly, $F_2$ acting on the body A for the same time $t$ changes its velocity from $u_1$ to $v_1$.

∴ $F_2 = \text{mass of the body A} \times \text{acceleration of the body A}$

$$= m_1 \times \frac{(v_1 - u_1)}{t} \quad \ldots(2)$$

Then by Newton’s third law of motion $F_1 = -F_2$

(i.e) $m_2 \times \frac{(v_2 - u_2)}{t} = -m_1 \times \frac{(v_1 - u_1)}{t}$

$$m_2 v_2 - m_2 u_2 = -m_1 v_1 + m_1 u_1$$

$$m_1 u_1 + m_2 u_2 = m_1 v_1 + m_2 v_2 \quad \ldots(3)$$

(i.e) total momentum before impact = total momentum after impact.

(i.e) total momentum of the system is a constant.

This proves the law of conservation of linear momentum.

**Applications of law of conservation of momentum**

The following examples illustrate the law of conservation of momentum.

(i) **Recoil of a gun**

Consider a gun and bullet of mass $m_g$ and $m_b$ respectively. The gun and the bullet form a single system. Before the gun is fired, both
the gun and the bullet are at rest. Therefore the velocities of the gun
and bullet are zero. Hence total momentum of the system before firing
is $m_y(0) + m_b(0) = 0$

When the gun is fired, the bullet moves forward and the gun
recoils backward. Let $v_b$ and $v_g$ are their respective velocities, the total
momentum of the bullet – gun system, after firing is $m_b v_b + m_g v_g$

According to the law of conservation of momentum, total
momentum before firing is equal to total momentum after firing.

\[(\text{i.e})\quad 0 = m_b v_b + m_g v_g \quad \text{(or)} \quad v_g = -\frac{m_b}{m_g} v_b\]

It is clear from this equation, that $v_g$ is directed opposite to $v_b$.
Knowing the values of $m_b$, $m_g$ and $v_b$, the recoil velocity of the gun $v_g$
can be calculated.

(i) Explosion of a bomb

Suppose a bomb is at rest before it explodes. Its momentum is
zero. When it explodes, it breaks up into many parts, each part having
a particular momentum. A part flying in one direction with a certain
momentum, there is another part moving in the opposite direction with
the same momentum. If the bomb explodes into two equal parts, they
will fly off in exactly opposite directions with the same speed, since
each part has the same mass.

Applications of Newton’s third law of motion

(i) Apparent loss of weight in a lift

Let us consider a man of mass $M$ standing on a weighing machine
placed inside a lift. The actual weight of the man = $Mg$. This weight
(action) is measured by the weighing machine and in turn, the machine
offers a reaction $R$. This reaction offered by the surface of contact on
the man is the apparent weight of the man.

Case (i)

When the lift is at rest:

The acceleration of the man = 0
Therefore, net force acting on the man = 0
From Fig. 2.27(i), \( R - Mg = 0 \) \( \text{or} \quad R = Mg \)
That is, the apparent weight of the man is equal to the actual weight.

Case (ii)

When the lift is moving uniformly in the upward or downward direction:

For uniform motion, the acceleration of the man is zero. Hence, in this case also the apparent weight of the man is equal to the actual weight.

Case (iii)

When the lift is accelerating upwards:

If $a$ be the upward acceleration of the man in the lift, then the net upward force on the man is $F = Ma$

From Fig 2.27(ii), the net force

$$F = R - Mg = Ma \ (or) \ R = M (g + a)$$

Therefore, apparent weight of the man is greater than actual weight.

Case (iv)

When the lift is accelerating downwards:

Let $a$ be the downward acceleration of the man in the lift, then the net downward force on the man is $F = Ma$

From Fig. 2.27 (iii), the net force

$$F = Mg - R = Ma \ (or) \ R = M (g - a)$$
Therefore, apparent weight of the man is less than the actual weight.

When the downward acceleration of the man is equal to the acceleration due to the gravity of earth, (i.e) \( a = g \)

\[
\therefore R = M (g - g) = 0
\]

Hence, the apparent weight of the man becomes zero. This is known as the weightlessness of the body.

(ii) Working of a rocket and jet plane

The propulsion of a rocket is one of the most interesting examples of Newton’s third law of motion and the law of conservation of momentum. The rocket is a system whose mass varies with time. In a rocket, the gases at high temperature and pressure, produced by the combustion of the fuel, are ejected from a nozzle. The reaction of the escaping gases provides the necessary thrust for the launching and flight of the rocket.

From the law of conservation of linear momentum, the momentum of the escaping gases must be equal to the momentum gained by the rocket. Consequently, the rocket is propelled in the forward direction opposite to the direction of the jet of escaping gases. Due to the thrust imparted to the rocket, its velocity and acceleration will keep on increasing. The mass of the rocket and the fuel system keeps on decreasing due to the escaping mass of gases.

2.5 Concurrent forces and Coplanar forces

The basic knowledge of various kinds of forces and motion is highly desirable for engineering and practical applications. The Newton’s laws of motion defines and gives the expression for the force. Force is a vector quantity and can be combined according to the rules of vector algebra. A force can be graphically represented by a straight line with an arrow, in which the length of the line is proportional to the magnitude of the force and the arrowhead indicates its direction.
A force system is said to be concurrent, if the lines of all forces intersect at a common point (Fig 2.28).

A force system is said to be coplanar, if the lines of the action of all forces lie in one plane (Fig 2.29).

2.5.1 Resultant of a system of forces acting on a rigid body

If two or more forces act simultaneously on a rigid body, it is possible to replace the forces by a single force, which will produce the same effect on the rigid body as the effect produced jointly by several forces. This single force is the resultant of the system of forces.

If \( \vec{P} \) and \( \vec{Q} \) are two forces acting on a body simultaneously in the same direction, their resultant is \( \vec{R} = \vec{P} + \vec{Q} \) and it acts in the same direction as that of the forces. If \( \vec{P} \) and \( \vec{Q} \) act in opposite directions, their resultant \( \vec{R} \) is \( \vec{R} = \vec{P} - \vec{Q} \) and the resultant is in the direction of the greater force.

If the forces \( \vec{P} \) and \( \vec{Q} \) act in directions which are inclined to each other, their resultant can be found by using parallelogram law of forces and triangle law of forces.

2.5.2 Parallelogram law of forces

If two forces acting at a point are represented in magnitude and direction by the two adjacent sides of a parallelogram, then their resultant is represented in magnitude and direction by the diagonal passing through the point.

Explanation

Consider two forces \( \vec{P} \) and \( \vec{Q} \) acting at a point \( O \) inclined at an angle \( \theta \) as shown in Fig. 2.30.

The forces \( \vec{P} \) and \( \vec{Q} \) are represented in magnitude and direction by the sides \( OA \) and \( OB \) of a parallelogram \( OACB \) as shown in Fig 2.30.
The resultant $\vec{R}$ of the forces $\vec{P}$ and $\vec{Q}$ is the diagonal $OC$ of the parallelogram. The magnitude of the resultant is

$$R = \sqrt{P^2 + Q^2 + 2PQ \cos \theta}$$

The direction of the resultant is

$$\alpha = \tan^{-1}\left(\frac{Q \sin \theta}{P + Q \cos \theta}\right)$$

### 2.5.3 Triangle law of forces

The resultant of two forces acting at a point can also be found by using triangle law of forces.

If two forces acting at a point are represented in magnitude and direction by the two adjacent sides of a triangle taken in order, then the closing side of the triangle taken in the reversed order represents the resultant of the forces in magnitude and direction.

Forces $\vec{P}$ and $\vec{Q}$ act at an angle $\theta$. In order to find the resultant of $\vec{P}$ and $\vec{Q}$, one can apply the head to tail method, to construct the triangle.

In Fig. 2.31, OA and AB represent $\vec{P}$ and $\vec{Q}$ in magnitude and direction. The closing side OB of the triangle taken in the reversed order represents the resultant $\vec{R}$ of the forces $\vec{P}$ and $\vec{Q}$. The magnitude and the direction of $\vec{R}$ can be found by using sine and cosine laws of triangles.

The triangle law of forces can also be stated as, if a body is in equilibrium under the action of three forces acting at a point, then the three forces can be completely represented by the three sides of a triangle taken in order.

If $\vec{P}$, $\vec{Q}$ and $\vec{R}$ are the three forces acting at a point and they are represented by the three sides of a triangle then $\frac{P}{OA} = \frac{Q}{AB} = \frac{R}{OB}$.  

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2.5.4 Equilibrant

According to Newton’s second law of motion, a body moves with a velocity if it is acted upon by a force. When the body is subjected to number of concurrent forces, it moves in a direction of the resultant force. However, if another force, which is equal in magnitude of the resultant but opposite in direction, is applied to a body, the body comes to rest. Hence, equilibrant of a system of forces is a single force, which acts along with the other forces to keep the body in equilibrium.

Let us consider the forces \( F_1, F_2, F_3 \) and \( F_4 \) acting on a body O as shown in Fig. 2.32a. If \( F \) is the resultant of all the forces and in order to keep the body at rest, an equal force (known as equilibrant) should act on it in the opposite direction as shown in Fig. 2.32b.

![Fig 2.32 Resultant and equilibrant](image)

From Fig. 2.32b, it is found that, resultant = − equilibrant

2.5.5 Resultant of concurrent forces

Consider a body O, which is acted upon by four forces as shown in Fig. 2.33a. Let \( \theta_1, \theta_2, \theta_3 \) and \( \theta_4 \) be the angles made by the forces with respect to X-axis.

Each force acting at O can be replaced by its rectangular components \( F_{1x} \) and \( F_{1y} \), \( F_{2x} \) and \( F_{2y} \), \( \ldots \) etc.

For example, for the force \( F_1 \) making an angle \( \theta_1 \), its components are, \( F_{1x} = F_1 \cos \theta_1 \) and \( F_{1y} = F_1 \sin \theta_1 \).

These components of forces produce the same effect on the body as the forces themselves. The algebraic sum of the horizontal components
F_1x, F_2x, F_3x, ... gives a single horizontal component \( R_x \)

\[
R_x = F_1x + F_2x + F_3x + F_4x = \sum F_x
\]

Similarly, the algebraic sum of the vertical components \( F_1y, F_2y, \ldots \) gives a single vertical component \( R_y \).

\[
R_y = F_1y + F_2y + F_3y + F_4y = \sum F_y
\]

Now, these two perpendicular components \( R_x \) and \( R_y \) can be added vectorially to give the resultant \( \vec{R} \).

\[
\therefore \text{From Fig. 2.33b, } R^2 = R_x^2 + R_y^2 \quad \text{(or) } \quad R = \sqrt{R_x^2 + R_y^2}
\]

and

\[
\tan \alpha = \frac{R_y}{R_x} \quad \text{(or) } \quad \alpha = \tan^{-1} \left( \frac{R_y}{R_x} \right)
\]

2.5.6 Lami's theorem

It gives the conditions of equilibrium for three forces acting at a point. Lami’s theorem states that if three forces acting at a point are in equilibrium, then each of the force is directly proportional to the sine of the angle between the remaining two forces.

Let us consider three forces \( \vec{P}, \vec{Q} \) and \( \vec{R} \) acting at a point O (Fig 2.34). Under the action of three forces, the point O is at rest, then by Lami’s theorem,
\[ P \propto \sin \alpha \]
\[ Q \propto \sin \beta \]
and \[ R \propto \sin \gamma, \]
then
\[ \frac{P}{\sin \alpha} = \frac{Q}{\sin \beta} = \frac{R}{\sin \gamma} = \text{constant} \]

**2.5.7 Experimental verification of triangle law, parallelogram law and Lami’s theorem**

Two smooth small pulleys are fixed, one each at the top corners of a drawing board kept vertically on a wall as shown in Fig. 2.35. The pulleys should move freely without any friction. A light string is made to pass over both the pulleys. Two slotted weights P and Q (of the order of 50 g) are taken and are tied to the two free ends of the string. Another short string is tied to the centre of the first string at O. A third slotted weight R is attached to the free end of the short string. The weights P, Q and R are adjusted such that the system is at rest.

The point O is in equilibrium under the action of the three forces P, Q and R acting along the strings. Now, a sheet of white paper is held just behind the string without touching them. The common knot O and the directions of OA, OB and OD are marked to represent in magnitude, the three forces P, Q and R on any convenient scale (like 50 g = 1 cm).
The experiment is repeated for different values of $P$, $Q$ and $R$ and the values are tabulated.

**To verify parallelogram law**

To determine the resultant of two forces $P$ and $Q$, a parallelogram $OACB$ is completed, taking $OA$ representing $P$, $OB$ representing $Q$ and the diagonal $OC$ gives the resultant. The length of the diagonal $OC$ and the angle $\angle COD$ are measured and tabulated (Table 2.2).

$OC$ is the resultant $R'$ of $P$ and $Q$. Since $O$ is at rest, this resultant $R'$ must be equal to the third force $R$ (equilibrant) which acts in the opposite direction. $OC = OD$. Also, both $OC$ and $OD$ are acting in the opposite direction. $\angle COD$ must be equal to $180^\circ$.

If $OC = OD$ and $\angle COD = 180^\circ$, one can say that parallelogram law of force is verified experimentally.

**Table 2.2 Verification of parallelogram law**

<table>
<thead>
<tr>
<th>S.No.</th>
<th>$P$</th>
<th>$Q$</th>
<th>$R$</th>
<th>$OA$</th>
<th>$OB$</th>
<th>$OD$ (R)</th>
<th>$OC$ (R')</th>
<th>$\angle COD$</th>
</tr>
</thead>
<tbody>
<tr>
<td>1.</td>
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</table>

**To verify Triangle Law**

According to triangle law of forces, the resultant of $P$ ($= OA = BC$) and $Q$ ($OB$) is represented in magnitude and direction by $OC$ which is taken in the reverse direction.

Alternatively, to verify the triangle law of forces, the ratios $\frac{P}{OA}$, $\frac{Q}{OB}$ and $\frac{R'}{OC}$ are calculated and are tabulated (Table 2.3). It will be found out that, all the three ratios are equal, which proves the triangle law of forces experimentally.

**Table 2.3 Verification of triangle law**

<table>
<thead>
<tr>
<th>S.No.</th>
<th>$P$</th>
<th>$Q$</th>
<th>$R'$</th>
<th>$OA$</th>
<th>$OB$</th>
<th>$OC$</th>
<th>$\frac{P}{OA}$</th>
<th>$\frac{Q}{OB}$</th>
<th>$\frac{R'}{OC}$</th>
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<tbody>
<tr>
<td>1.</td>
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</tbody>
</table>
To verify Lami’s theorem

To verify Lami’s theorem, the angles between the three forces, P, Q and R (i.e.) \( \angle BOD = \alpha, \angle AOD = \beta \) and \( \angle AOB = \gamma \) are measured using protractor and tabulated (Table 2.4). The ratios \( \frac{P}{\sin \alpha}, \frac{Q}{\sin \beta} \) and \( \frac{R}{\sin \gamma} \) are calculated and it is found that all the three ratios are equal and this verifies the Lami’s theorem.

**Table 2.4 Verification of Lami’s theorem**

<table>
<thead>
<tr>
<th>S.No.</th>
<th>P</th>
<th>Q</th>
<th>R</th>
<th>( \alpha )</th>
<th>( \beta )</th>
<th>( \gamma )</th>
<th>( \frac{P}{\sin \alpha} )</th>
<th>( \frac{Q}{\sin \beta} )</th>
<th>( \frac{R}{\sin \gamma} )</th>
</tr>
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<tbody>
<tr>
<td>1.</td>
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</table>

2.5.8 Conditions of equilibrium of a rigid body acted upon by a system of concurrent forces in plane

(i) If an object is in equilibrium under the action of three forces, the resultant of two forces must be equal and opposite to the third force. Thus, the line of action of the third force must pass through the point of intersection of the lines of action of the other two forces. In other words, the system of three coplanar forces in equilibrium, must obey parallelogram law, triangle law of forces and Lami’s theorem. This condition ensures the absence of translational motion in the system.

(ii) The algebraic sum of the moments about any point must be equal to zero, \( \Sigma M = 0 \) (i.e) the sum of clockwise moments about any point must be equal to the sum of anticlockwise moments about the same point. This condition ensures, the absence of rotational motion.

2.6 Uniform circular motion

The revolution of the Earth around the Sun, rotating fly wheel, electrons revolving around the nucleus, spinning top, the motion of a fan blade, revolution of the moon around the Earth etc. are some examples of circular motion. In all the above cases, the bodies or particles travel in a circular path. So, it is necessary to understand the motion of such bodies.
When a particle moves on a circular path with a constant speed, then its motion is known as uniform circular motion in a plane. The magnitude of velocity in circular motion remains constant but the direction changes continuously.

Let us consider a particle of mass $m$ moving with a velocity $v$ along the circle of radius $r$ with centre $O$ as shown in Fig 2.36. $P$ is the position of the particle at a given instant of time such that the radial line $OP$ makes an angle $\theta$ with the reference line $DA$. The magnitude of the velocity remains constant, but its direction changes continuously. The linear velocity always acts tangentially to the position of the particle (i.e) in each position, the linear velocity $\vec{v}$ is perpendicular to the radius vector $\vec{r}$.

**2.6.1 Angular displacement**

Let us consider a particle of mass $m$ moving along the circular path of radius $r$ as shown in Fig. 2.37. Let the initial position of the particle be $A$. $P$ and $Q$ are the positions of the particle at any instants of time $t$ and $t + dt$ respectively. Suppose the particle traverses a distance $ds$ along the circular path in time interval $dt$. During this interval, it moves through an angle $d\theta = \theta_2 - \theta_1$. The angle swept by the radius vector at a given time is called the angular displacement of the particle.

If $r$ be the radius of the circle, then the angular displacement is given by $d\theta = \frac{ds}{r}$. The angular displacement is measured in terms of radian.

**2.6.2 Angular velocity**

The rate of change of angular displacement is called the angular velocity of the particle.
Let $d\theta$ be the angular displacement made by the particle in time $dt$, then the angular velocity of the particle is $\omega = \frac{d\theta}{dt}$. Its unit is rad s$^{-1}$ and dimensional formula is $T^{-1}$.

For one complete revolution, the angle swept by the radius vector is $360^\circ$ or $2\pi$ radians. If $T$ is the time taken for one complete revolution, known as period, then the angular velocity of the particle is $\omega = \frac{\theta}{t} = \frac{2\pi}{T}$.

If the particle makes $n$ revolutions per second, then

$$\omega = 2\pi \left(\frac{1}{T}\right) = 2\pi n \text{ where } n = \frac{1}{T} \text{ is the frequency of revolution.}$$

### 2.6.3 Relation between linear velocity and angular velocity

Let us consider a body P moving along the circumference of a circle of radius $r$ with linear velocity $v$ and angular velocity $\omega$ as shown in Fig. 2.38. Let it move from P to Q in time $dt$ and $d\theta$ be the angle swept by the radius vector.

Let $\text{PQ} = ds$, be the arc length covered by the particle moving along the circle, then the angular displacement $d\theta$ is expressed as

$$d\theta = \frac{ds}{r} \quad \text{(or)} \quad \frac{d\theta}{dt} = \frac{v}{r}$$

(i.e) Angular velocity $\omega = \frac{v}{r}$ or $v = \omega r$

In vector notation, $\vec{v} = \vec{\omega} \times \vec{r}$

Thus, for a given angular velocity $\omega$, the linear velocity $v$ of the particle is directly proportional to the distance of the particle from the centre of the circular path (i.e. for a body in a uniform circular motion, the angular velocity is the same for all points in the body but linear velocity is different for different points of the body.

### 2.6.4 Angular acceleration

If the angular velocity of the body performing rotatory motion is non-uniform, then the body is said to possess angular acceleration.
The rate of change of angular velocity is called angular acceleration.

If the angular velocity of a body moving in a circular path changes from \(\omega_1\) to \(\omega_2\) in time \(t\) then its angular acceleration is

\[
\alpha = \frac{d\omega}{dt} = \frac{d}{dt} \left(\frac{d\theta}{dt}\right) = \frac{d^2\theta}{dt^2} = \frac{\omega_2 - \omega_1}{t}.
\]

The angular acceleration is measured in terms of rad s\(^{-2}\) and its dimensional formula is T\(^{-2}\).

2.6.5 Relation between linear acceleration and angular acceleration

If \(dv\) is the small change in linear velocity in a time interval \(dt\) then linear acceleration is

\[
a = \frac{dv}{dt} = \frac{d}{dt} (r\omega) = r \frac{d\omega}{dt} = ra.
\]

2.6.6 Centripetal acceleration

The speed of a particle performing uniform circular motion remains constant throughout the motion but its velocity changes continuously due to the change in direction (i.e., the particle executing uniform circular motion is said to possess an acceleration).

Consider a particle executing circular motion of radius \(r\) with linear velocity \(v\) and angular velocity \(\omega\). The linear velocity of the particle acts along the tangential line. Let \(d\theta\) be the angle described by the particle at the centre when it moves from A to B in time \(dt\).

At A and B, linear velocity \(v\) acts along AH and BT respectively. In Fig. 2.39 \(\angle AOB = d\theta = \angle HET \therefore\) angle subtended by the two radii of a circle = angle subtended by the two tangents).

The velocity \(v\) at B of the particle makes an angle \(d\theta\) with the line BC and hence it is resolved horizontally as \(v \cos d\theta\) along BC and vertically as \(v \sin d\theta\) along BD.

:. The change in velocity along the horizontal direction = \(v \cos d\theta - v\)

If \(d\theta\) is very small, \(\cos d\theta = 1\)
∴ Change in velocity along the horizontal direction = \(v - v = 0\)

(i.e) there is no change in velocity in the horizontal direction.

The change in velocity in the vertical direction (i.e along AO) is
\[dv = v \sin d\theta - 0 = v \sin d\theta\]
If \(d\theta\) is very small, \(\sin d\theta = d\theta\)

∴ The change in velocity in the vertical direction (i.e) along radius of the circle
\[dv = v.d\theta \quad \ldots(1)\]
But, acceleration \(a = \frac{dv}{dt} = \frac{v.d\theta}{dt} = v.\omega \quad \ldots(2)\)
where \(\omega = \frac{d\theta}{dt}\) is the angular velocity of the particle.

We know that \(v = r.\omega \quad \ldots(3)\)

From equations (2) and (3),
\[a = r.\omega.\omega = r.\omega^2 = \frac{v^2}{r} \quad \ldots(4)\]

Hence, the acceleration of the particle producing uniform circular motion is equal to \(\frac{v^2}{r}\) and is along AO (i.e) directed towards the centre of the circle. This acceleration is directed towards the centre of the circle along the radius and perpendicular to the velocity of the particle. This acceleration is known as centripetal or radial or normal acceleration.

2.6.7 Centripetal force

According to Newton’s first law of motion, a body possesses the property called directional inertia (i.e) the inability of the body to change its direction. This means that without the application of an external force, the direction of motion can not be changed. Thus when a body is moving along a circular path, some force must be acting upon it, which continuously changes the body from its straight-line path (Fig 2.40). It makes clear that the applied force should have no component in the direction of the motion of the body or the force must act at every
point perpendicular to the direction of motion of the body. This force, therefore, must act along the radius and should be directed towards the centre.

Hence for circular motion, a constant force should act on the body, along the radius towards the centre and perpendicular to the velocity of the body. This force is known as centripetal force.

If \( m \) is the mass of the body, then the magnitude of the centripetal force is given by

\[
F = \text{mass} \times \text{centripetal acceleration} = m \left( \frac{v^2}{r} \right) = \frac{mv^2}{r} = m (r \omega^2)
\]

**Examples**

Any force like gravitational force, frictional force, electric force, magnetic force etc. may act as a centripetal force. Some of the examples of centripetal force are:

(i) In the case of a stone tied to the end of a string whirled in a circular path, the centripetal force is provided by the tension in the string.

(ii) When a car takes a turn on the road, the frictional force between the tyres and the road provides the centripetal force.

(iii) In the case of planets revolving round the Sun or the moon revolving round the earth, the centripetal force is provided by the gravitational force of attraction between them.

(iv) For an electron revolving round the nucleus in a circular path, the electrostatic force of attraction between the electron and the nucleus provides the necessary centripetal force.

### 2.6.8 Centrifugal reaction

According to Newton’s third law of motion, for every action there is an equal and opposite reaction. *The equal and opposite reaction to the centripetal force is called centrifugal reaction, because it tends to take the body away from the centre.* In fact, the centrifugal reaction is a pseudo or apparent force, acts or assumed to act because of the acceleration of the rotating body.

In the case of a stone tied to the end of the string is whirled in a circular path, not only the stone is acted upon by a force (centripetal force) along the string towards the centre, but the stone also exerts an equal and opposite force on the hand (centrifugal force) away from the
centre, along the string. On releasing the string, the tension disappears and the stone flies off tangentially to the circular path along a straight line as enunciated by Newton’s first law of motion.

When a car is turning round a corner, the person sitting inside the car experiences an outward force. It is because of the fact that no centripetal force is supplied by the person. Therefore, to avoid the outward force, the person should exert an inward force.

2.6.9 Applications of centripetal forces

(i) Motion in a vertical circle

Let us consider a body of mass \( m \) tied to one end of the string which is fixed at O and it is moving in a vertical circle of radius \( r \) about the point O as shown in Fig. 2.41. The motion is circular but is not uniform, since the body speeds up while coming down and slows down while going up.

Suppose the body is at \( P \) at any instant of time \( t \), the tension \( T \) in the string always acts towards \( O \).

The weight \( mg \) of the body at \( P \) is resolved along the string as \( mg \cos \theta \) which acts outwards and \( mg \sin \theta \), perpendicular to the string.

When the body is at \( P \), the following forces acts on it along the string.

(i) \( mg \cos \theta \) acts along \( OP \) (outwards)
(ii) tension \( T \) acts along \( PO \) (inwards)

Net force on the body at \( P \) acting along \( PO = T - mg \cos \theta \)

This must provide the necessary centripetal force \( \frac{mv^2}{r} \).

Therefore, \( T - mg \cos \theta = \frac{mv^2}{r} \)

\[ T = mg \cos \theta + \frac{mv^2}{r} \]  \( \ldots(1) \)
At the lowest point A of the path, $\theta = 0^\circ$, $\cos 0^\circ = 1$ then from equation (1), \[ T_A = mg + \frac{mv_A^2}{r} \] \hspace{1cm} \text{...(2)}

At the highest point of the path, i.e. at B, $\theta = 180^\circ$. Hence $\cos 180^\circ = -1$

\[ T_B = -mg + \frac{mv_B^2}{r} = \frac{mv_B^2}{r} - mg \]

\[ T_B = m \left( \frac{v_B^2}{r} - g \right) \] \hspace{1cm} \text{...(3)}

If $T_B > 0$, then the string remains taut while if $T_B < 0$, the string slackens and it becomes impossible to complete the motion in a vertical circle.

If the velocity $v_B$ is decreased, the tension $T_B$ in the string also decreases, and becomes zero at a certain minimum value of the speed called critical velocity. Let $v_C$ be the minimum value of the velocity, then at $v_B = v_C$, $T_B = 0$. Therefore from equation (3),

\[ \frac{mv_C^2}{r} - mg = 0 \] (or) \[ v_C^2 = rg \]

(i.e) \[ v_C = \sqrt{rg} \] \hspace{1cm} \text{...(4)}

If the velocity of the body at the highest point B is below this critical velocity, the string becomes slack and the body falls downwards instead of moving along the circular path. In order to ensure that the velocity $v_B$ at the top is not lesser than the critical velocity $\sqrt{rg}$, the minimum velocity $v_A$ at the lowest point should be in such a way that $v_B$ should be $\sqrt{rg}$. (i.e) the motion in a vertical circle is possible only if $v_B \geq \sqrt{rg}$.

The velocity $v_A$ of the body at the bottom point A can be obtained by using law of conservation of energy. When the stone rises from A to B, i.e through a height $2r$, its potential energy increases by an amount equal to the decrease in kinetic energy. Thus,

\[ \text{(Potential energy at A + Kinetic energy at A ) = (Potential energy at B + Kinetic energy at B)} \]

(i.e.) \[ \theta + \frac{1}{2} \ m \ v_A^2 = mg \ (2r) + \frac{1}{2} \ m \ v_B^2 \]

Dividing by $\frac{m}{2}$, \[ v_A^2 = v_B^2 + 4gr \] \hspace{1cm} \text{...(5)}
But from equation (4), \( v_B^2 = gr \) \( \because v_B = v_C \)

\[ \therefore \text{Equation (5) becomes, } v_A^2 = gr + 4gr \text{ (or) } v_A = \sqrt{5gr} \quad \ldots(6) \]

Substituting \( v_A \) from equation (6) in (2),

\[ T_A = mg + \frac{m(5gr)}{r} = mg + 5mg = 6mg \quad \ldots(7) \]

While rotating in a vertical circle, the stone must have a velocity greater than \( \sqrt{5gr} \) or tension greater than \( 6mg \) at the lowest point, so that its velocity at the top is greater than \( \sqrt{gr} \) or tension \( > 0 \).

An aeroplane while looping a vertical circle must have a velocity greater than \( \sqrt{5gr} \) at the lowest point, so that its velocity at the top is greater than \( \sqrt{gr} \). In that case, pilot sitting in the aeroplane will not fall.

**(ii) Motion on a level circular road**

When a vehicle goes round a level curved path, it should be acted upon by a centripetal force. While negotiating the curved path, the wheels of the car have a tendency to leave the curved path and regain the straight-line path. Frictional force between the tyres and the road opposes this tendency of the wheels. This frictional force, therefore, acts towards the centre of the circular path and provides the necessary centripetal force.

In Fig. 2.42, weight of the vehicle \( mg \) acts vertically downwards. \( R_1, \ R_2 \) are the forces of normal reaction of the road on the wheels. As the road is level (horizontal), \( R_1, \ R_2 \) act vertically upwards. Obviously,

\[ R_1 + R_2 = mg \quad \ldots(1) \]

Let \( \mu \) be the coefficient of friction between the tyres and the road.

Friction: Whenever a body slides over another body, a force comes into play between the two surfaces in contact and this force is known as frictional force. The frictional force always acts in the opposite direction to that of the motion of the body. The frictional force depends on the normal reaction. (Normal reaction is a perpendicular reactional force that acts on the body at the point of contact due to its own weight) (i.e Frictional force a normal reaction \( F = \mu R \) or \( F = \mu R \) where \( \mu \) is a proportionality constant and is known as the coefficient of friction. The coefficient of friction depends on the nature of the surface.
road, $F_1$ and $F_2$ be the forces of friction between the tyres and the road, directed towards the centre of the curved path.

\[ F_1 = \mu R_1 \text{ and } F_2 = \mu R_2 \quad \ldots (2) \]

If $v$ is velocity of the vehicle while negotiating the curve, the centripetal force required is $\frac{mv^2}{r}$.

As this force is provided only by the force of friction,

\[ \frac{mv^2}{r} \leq (F_1 + F_2) \leq (\mu R_1 + \mu R_2) \leq \mu (R_1 + R_2) \]

\[ \frac{mv^2}{r} \leq \mu mg \quad (\because R_1 + R_2 = mg) \]

\[ v^2 \leq \mu rg \]

\[ v \leq \sqrt{\mu rg} \]

Hence the maximum velocity with which a car can go round a level curve without skidding is $v = \sqrt{\mu rg}$. The value of $v$ depends on radius $r$ of the curve and coefficient of friction $\mu$ between the tyres and the road.

(iii) Banking of curved roads and tracks

When a car goes round a level curve, the force of friction between the tyres and the road provides the necessary centripetal force. If the frictional force, which acts as centripetal force and keeps the body moving along the circular road is not enough to provide the necessary centripetal force, the car will skid. In order to avoid skidding, while going round a curved path the outer edge of the road is raised above the level of the inner edge. This is known as banking of curved roads or tracks.

**Bending of a cyclist round a curve**

A cyclist has to bend slightly towards the centre of the circular track in order to take a safe turn without slipping.

Fig. 2.43 shows a cyclist taking a turn towards his right on a circular path of radius $r$. Let $m$ be the mass of the cyclist along with the bicycle and $v$, the velocity. When the cyclist negotiates the curve, he bends inwards from the vertical, by an angle $\theta$. Let $R$ be the reaction
of the ground on the cyclist. The reaction \( R \) may be resolved into two components: (i) the component \( R \sin \theta \), acting towards the centre of the curve, providing necessary centripetal force for circular motion and (ii) the component \( R \cos \theta \), balancing the weight of the cyclist along with the bicycle.

\[
\begin{align*}
\text{(i.e)} \quad R \sin \theta &= \frac{mv^2}{r} \quad \text{...(1)} \\
\text{and} \quad R \cos \theta &= mg \quad \text{...(2)}
\end{align*}
\]

Dividing equation (1) by (2),

\[
\frac{R \sin \theta}{R \cos \theta} = \frac{mv^2}{rg} \quad \text{...(3)}
\]

Thus for less bending of cyclist (i.e for \( \theta \) to be small), the velocity \( v \) should be smaller and radius \( r \) should be larger.

For a banked road (Fig. 2.44), let \( h \) be the elevation of the outer edge of the road above the inner edge and \( l \) be the width of the road then,

\[
\sin \theta = \frac{h}{l} \quad \text{...(4)}
\]
For small values of $\theta$, $\sin \theta = \tan \theta$

Therefore from equations (3) and (4)

$$\tan \theta = \frac{h}{L} = \frac{v^2}{rg}$$ ...

Obviously, a road or track can be banked correctly only for a particular speed of the vehicle. Therefore, the driver must drive with a particular speed at the circular turn. If the speed is higher than the desired value, the vehicle tends to slip outward at the turn but then the frictional force acts inwards and provides the additional centripetal force. Similarly, if the speed of the vehicle is lower than the desired speed it tends to slip inward at the turn but now the frictional force acts outwards and reduces the centripetal force.

**Condition for skidding**

When the centripetal force is greater than the frictional force, skidding occurs. If $\mu$ is the coefficient of friction between the road and tyre, then the limiting friction (frictional force) is $f = \mu R$ where normal reaction $R = mg$

$$\therefore f = \mu (mg)$$

Thus for skidding,

Centripetal force $>$ Frictional force

$$\frac{mv^2}{r} > \mu (mg)$$

$$\frac{v^2}{rg} > \mu$$

But $$\frac{v^2}{rg} = \tan \theta$$

$$\therefore \tan \theta > \mu$$

(i.e) when the tangent of the angle of banking is greater than the coefficient of friction, skidding occurs.

**2.7 Work**

The terms work and energy are quite familiar to us and we use them in various contexts. In everyday life, the term work is used to refer to any form of activity that requires the exertion of mental or muscular efforts. In physics, *work is said to be done by a force or*
against the direction of the force, when the point of application of the force moves towards or against the direction of the force. If no displacement takes place, no work is said to be done. Therefore for work to be done, two essential conditions should be satisfied:

(i) a force must be exerted
(ii) the force must cause a motion or displacement

If a particle is subjected to a force $F$ and if the particle is displaced by an infinitesimal displacement $ds$, the work done $dw$ by the force is $dw = \vec{F} \cdot ds$.

The magnitude of the above dot product is $F \cos \theta ds$.

(i.e) $dw = F \cos \theta ds$ where $\theta = \text{angle between } \vec{F} \text{ and } ds$. (Fig. 2.45)

Thus, the work done by a force during an infinitesimal displacement is equal to the product of the displacement $ds$ and the component of the force $F \cos \theta$ in the direction of the displacement.

Work is a scalar quantity and has magnitude but no direction.

The work done by a force when the body is displaced from position $P$ to $P_1$ can be obtained by integrating the above equation,

$W = \int dw = \int (F \cos \theta) \, ds$.

**Work done by a constant force**

When the force $F$ acting on a body has a constant magnitude and acts at a constant angle $\theta$ from the straight line path of the particle as shown as Fig. 2.46, then,

$W = F \cos \theta \int_{s_1}^{s_2} ds = F \cos \theta (s_2 - s_1)$

The graphical representation of work done by a constant force is shown in Fig 2.47.

$W = F \cos \theta (s_2 - s_1) = \text{area ABCD}$
Work done by a variable force

If the body is subjected to a varying force $F$ and displaced along $X$ axis as shown in Fig 2.48, work done

$$dw = F \cos \theta \ ds = \text{area of the small element abcd.}$$

∴ The total work done when the body moves from $s_1$ to $s_2$ is

$$\Sigma dw = W = \text{area under the curve } P_1P_2 = \text{area } S_1P_1P_2S_2$$

The unit of work is joule. One joule is defined as the work done by a force of one newton when its point of application moves by one metre along the line of action of the force.

Special cases

(i) When $\theta = 0$, the force $F$ is in the same direction as the displacement $s$.

∴ Work done, $W = F \ s \cos 0 = F \ s$

(ii) When $\theta = 90^\circ$, the force under consideration is normal to the direction of motion.

∴ Work done, $W = F \ s \cos 90^\circ = 0$

For example, if a body moves along a frictionless horizontal surface, its weight and the reaction of the surface, both normal to the surface, do no work. Similarly, when a stone tied to a string is whirled around in a circle with uniform speed, the centripetal force continuously changes the direction of motion. Since this force is always normal to the direction of motion of the object, it does no work.

(iii) When $\theta = 180^\circ$, the force $F$ is in the opposite direction to the displacement.
Work done \( (W) = F \cdot s \cos 180^\circ = -F \cdot s \)

(eg.) The frictional force that slows the sliding of an object over a surface does a negative work.

A positive work can be defined as the work done by a force and a negative work as the work done against a force.

2.8 Energy

Energy can be defined as the capacity to do work. Energy can manifest itself in many forms like mechanical energy, thermal energy, electric energy, chemical energy, light energy, nuclear energy, etc.

The energy possessed by a body due to its position or due to its motion is called mechanical energy.

The mechanical energy of a body consists of potential energy and kinetic energy.

2.8.1 Potential energy

The potential energy of a body is the energy stored in the body by virtue of its position or the state of strain. Hence water stored in a reservoir, a wound spring, compressed air, stretched rubber chord, etc, possess potential energy.

Potential energy is given by the amount of work done by the force acting on the body, when the body moves from its given position to some other position.

Expression for the potential energy

Let us consider a body of mass \( m \), which is at rest at a height \( h \) above the ground as shown in Fig 2.49. The work done in raising the body from the ground to the height \( h \) is stored in the body as its potential energy and when the body falls to the ground, the same amount of work can be got back from it. Now, in order to lift the body vertically up, a force \( mg \) equal to the weight of the body should be applied.

When the body is taken vertically up through a height \( h \), then work done, \( W = \text{Force} \times \text{displacement} \)

\[ \therefore W = mg \times h \]

This work done is stored as potential energy in the body

\[ \therefore E_p = mgh \]
2.8.2 Kinetic energy

The kinetic energy of a body is the energy possessed by the body by virtue of its motion. It is measured by the amount of work that the body can perform against the impressed forces before it comes to rest. A falling body, a bullet fired from a rifle, a swinging pendulum, etc. possess kinetic energy.

A body is capable of doing work if it moves, but in the process of doing work its velocity gradually decreases. The amount of work that can be done depends both on the magnitude of the velocity and the mass of the body. A heavy bullet will penetrate a wooden plank deeper than a light bullet of equal size moving with equal velocity.

Expression for Kinetic energy

Let us consider a body of mass \(m\) moving with a velocity \(v\) in a straightline as shown in Fig. 2.50. Suppose that it is acted upon by a constant force \(F\) resisting its motion, which produces retardation \(a\) (decrease in acceleration is known as retardation). Then

\[
F = \text{mass} \times \text{retardation} = -ma \quad \ldots(1)
\]

Let \(dx\) be the displacement of the body before it comes to rest.

But the retardation is

\[
\alpha = \frac{dv}{dt} = \frac{dv}{dx} \times \frac{dx}{dt} = \frac{dv}{dx} \times v \quad \ldots(2)
\]

Substituting equation (2) in (1),

\[
F = -mv \frac{dv}{dx} \quad \ldots(3)
\]

Hence the work done in bringing the body to rest is given by,

\[
W = \int F \, dx = \left[ \int_{0}^{v} \left( -\frac{mv}{v} \, \frac{dv}{dx} \right) \, dx = \left[ -m \right]_{v}^{0} \, v \, dx \right] = \frac{1}{2} mv^{2} \quad \ldots(4)
\]

This work done is equal to kinetic energy of the body.
2.8.3 Principle of work and energy (work – energy theorem)

Statement

The work done by a force acting on the body during its displacement is equal to the change in the kinetic energy of the body during that displacement.

Proof

Let us consider a body of mass \( m \) acted upon by a force \( F \) and moving with a velocity \( v \) along a path as shown in Fig. 2.51. At any instant, let \( P \) be the position of the body from the origin \( O \). Let \( \theta \) be the angle made by the direction of the force with the tangential line drawn at \( P \).

The force \( F \) can be resolved into two rectangular components:

(i) \( F_t = F \cos \theta \), tangentially and

(ii) \( F_n = F \sin \theta \), normally at \( P \).

But \( F_t = ma_t \) 

\[
\therefore \quad F \cos \theta = ma_t \quad \ldots(1)
\]

where \( a_t \) is the acceleration of the body in the tangential direction

\[
\therefore \quad F \cos \theta = ma_t \quad \ldots(2)
\]

But \( a_t = \frac{dv}{dt} \quad \ldots(3) \]

\[
\therefore \quad \text{substituting equation (3) in (2),}
\]

\[
F \cos \theta = m \left( \frac{dv}{dt} \right) = m \frac{dv}{ds} \frac{ds}{dt} \quad \ldots(4)
\]

\[
F \cos \theta \, ds = mv \, dv \quad \ldots(5)
\]

where \( ds \) is the small displacement.

Let \( v_1 \) and \( v_2 \) be the velocities of the body at the positions 1 and 2 and the corresponding distances be \( s_1 \) and \( s_2 \).

Integrating the equation (5),

\[
\int_{s_1}^{s_2} (F \cos \theta) \, ds = \int_{v_1}^{v_2} mv \, dv \quad \ldots(6)
\]
But \[ \int_{s_1}^{s_2} (F \cos \theta) \, ds = W_{1\to2} \]  
...(7)

where \( W_{1\to2} \) is the work done by the force.

From equation (6) and (7),

\[
W_{1\to2} = \int_{v_1}^{v_2} mv \, dv
\]

\[ = m \left[ \frac{v^2}{2} \right]_{v_1}^{v_2} = \frac{mv_2^2}{2} - \frac{mv_1^2}{2} \]  
...(8)

Therefore work done

\[ = \text{final kinetic energy} - \text{initial kinetic energy} \]

\[ = \text{change in kinetic energy} \]

This is known as Work–energy theorem.

2.8.4 Conservative forces and non-conservative forces

**Conservative forces**

If the work done by a force in moving a body between two positions is independent of the path followed by the body, then such a force is called as a conservative force.

Examples: force due to gravity, spring force and elastic force.

The work done by the conservative forces depends only upon the initial and final position of the body.

\[ (i.e.) \quad \int \vec{F} \cdot d\vec{r} = 0 \]

The work done by a conservative force around a closed path is zero.

**Non conservative forces**

Non-conservative force is the force, which can perform some resultant work along an arbitrary closed path of its point of application.

The work done by the non-conservative force depends upon the path of the displacement of the body.
\[ (\text{i.e.}) \oint \vec{F} \cdot d\vec{r} \neq 0 \]

(e.g) Frictional force, viscous force, etc.

### 2.8.5 Law of conservation of energy

The law states that, if a body or system of bodies is in motion under a conservative system of forces, the sum of its kinetic energy and potential energy is constant.

**Explanation**

From the principle of work and energy,

\[
\text{Work done} = \text{change in the kinetic energy}
\]

(i.e) \( W_{1 \to 2} = E_{k2} - E_{k1} \) \...(1)

If a body moves under the action of a conservative force, work done is stored as potential energy.

\[ W_{1 \to 2} = -(E_{p2} - E_{p1}) \] \...(2)

Work done is equal to negative change of potential energy. Combining the equation (1) and (2),

\[ E_{k2} - E_{k1} = -(E_{p2} - E_{p1}) \text{ (or) } E_{p1} + E_{k1} = E_{p2} + E_{k2} \] \...(3)

which means that the sum of the potential energy and kinetic energy of a system of particles remains constant during the motion under the action of the conservative forces.

### 2.8.6 Power

It is defined as the rate at which work is done.

\[
\text{power} = \frac{\text{work done}}{\text{time}}
\]

Its unit is watt and dimensional formula is \( \text{ML}^2\text{T}^{-3} \).

*Power is said to be one watt, when one joule of work is said to be done in one second.*

If \( dw \) is the work done during an interval of time \( dt \) then,

\[
\text{power} = \frac{dw}{dt} \] \...(1)

But \( dw = (F \cos \theta) \, ds \) \...(2)
where $\theta$ is the angle between the direction of the force and displacement. $F \cos \theta$ is component of the force in the direction of the small displacement $ds$.

Substituting equation (2) in (1) power = \frac{(F \cos \theta) \, ds}{dt} = (F \cos \theta) \, \frac{ds}{dt} = (F \cos \theta) \, v \quad \left( \because \frac{ds}{dt} = v \right) \\
\therefore \text{power} = (F \cos \theta) \, v$

If $F$ and $v$ are in the same direction, then

$\text{power} = F \, v \cos 0 = F \, v = \text{Force} \times \text{velocity}$

It is also represented by the dot product of $F$ and $v$.

(i.e) $P = \vec{F} \cdot \vec{v}$

### 2.9 Collisions

A collision between two particles is said to occur if they physically strike against each other or if the path of the motion of one is influenced by the other. In physics, the term collision does not necessarily mean that a particle actually strikes. In fact, two particles may not even touch each other and yet they are said to collide if one particle influences the motion of the other.

When two bodies collide, each body exerts a force on the other. The two forces are exerted simultaneously for an equal but short interval of time. According to Newton’s third law of motion, each body exerts an equal and opposite force on the other at each instant of collision. During a collision, the two fundamental conservation laws namely, the law of conservation of momentum and that of energy are obeyed and these laws can be used to determine the velocities of the bodies after collision.

Collisions are divided into two types: (i) elastic collision and (ii) inelastic collision

#### 2.9.1 Elastic collision

If the kinetic energy of the system is conserved during a collision, it is called an elastic collision. (i.e) The total kinetic energy before collision and after collision remains unchanged. The collision between subatomic
particles is generally elastic. The collision between two steel or glass balls is nearly elastic. In elastic collision, the linear momentum and kinetic energy of the system are conserved.

**Elastic collision in one dimension**

If the two bodies after collision move in a straight line, the collision is said to be of one dimension.

Consider two bodies A and B of masses $m_1$ and $m_2$ moving along the same straight line in the same direction with velocities $u_1$ and $u_2$ respectively as shown in Fig. 2.54. Let us assume that $u_1$ is greater than $u_2$. The bodies A and B suffer a head on collision when they strike and continue to move along the same straight line with velocities $v_1$ and $v_2$ respectively.

![Fig 2.54 Elastic collision in one dimension](image)

From the law of conservation of linear momentum,

Total momentum before collision =

Total momentum after collision

$m_1u_1 + m_2u_2 = m_1v_1 + m_2v_2$ ... (1)

Since the kinetic energy of the bodies is also conserved during the collision

Total kinetic energy before collision =

Total kinetic energy after collision

$\frac{1}{2}m_1u_1^2 + \frac{1}{2}m_2u_2^2 = \frac{1}{2}m_1v_1^2 + \frac{1}{2}m_2v_2^2$ ... (2)

$m_1u_1^2 - m_1v_1^2 = m_2v_2^2 - m_2u_2^2$ ... (3)

From equation (1) $m_1(u_1 - v_1) = m_2(v_2 - u_2)$ ... (4)

Dividing equation (3) by (4),

$\frac{u_1^2 - v_1^2}{u_1 - v_1} = \frac{v_2^2 - u_2^2}{v_2 - u_2}$  (or)  $u_1 + v_1 = u_2 + v_2$

$(u_1 - u_2) = (v_2 - v_1)$ ... (5)
Equation (5) shows that in an elastic one-dimensional collision, the relative velocity with which the two bodies approach each other before collision is equal to the relative velocity with which they recede from each other after collision.

From equation (5), \[ v_2 = u_1 - u_2 + v_1 \] ... (6)

Substituting \( v_2 \) in equation (4),
\[ m_1 (u_1 - v_1) = m_2 (v_1 - u_2 + u_1 - u_2) \]
\[ m_1 u_1 - m_1 v_1 = m_2 u_1 - 2m_2 u_2 + m_2 v_1 \]
\[ (m_1 + m_2) v_1 = m_1 u_1 - m_2 u_1 + 2m_2 u_2 \]
\[ (m_1 + m_2) v_1 = u_1 (m_1 - m_2) + 2m_2 u_2 \]
\[ v_1 = u_1 \left( \frac{m_1 - m_2}{m_1 + m_2} \right) + \frac{2m_2 u_2}{m_1 + m_2} \] ... (7)

Similarly, \[ v_2 = \frac{2m_1 u_1}{m_1 + m_2} + \frac{u_2 (m_2 - m_1)}{m_1 + m_2} \] ... (8)

**Special cases**

**Case (i) :** If the masses of colliding bodies are equal, i.e. \( m_1 = m_2 \)
\[ v_1 = u_2 \quad \text{and} \quad v_2 = u_1 \] ... (9)

After head on elastic collision, the velocities of the colliding bodies are mutually interchanged.

**Case (ii) :** If the particle \( B \) is initially at rest, (i.e) \( u_2 = 0 \) then
\[ v_1 = \frac{(m_1 - m_2)}{(m_1 + m_2)} u_1 \] ... (10)

and \[ v_2 = \frac{2m_1}{(m_1 + m_2)} u_1 \] ... (11)

**2.9.2 Inelastic collision**

During a collision between two bodies if there is a loss of kinetic energy, then the collision is said to be an inelastic collision. Since there is always some loss of kinetic energy in any collision, collisions are generally inelastic. In inelastic collision, the linear momentum is conserved but the energy is not conserved. If two bodies stick together, after colliding, the collision is perfectly inelastic but it is a special case of inelastic collision called plastic collision. (eg) a bullet striking a block
of wood and being embedded in it. The loss of kinetic energy usually results in the form of heat or sound energy.

Let us consider a simple situation in which the inelastic head on collision between two bodies of masses \( m_A \) and \( m_B \) takes place. Let the colliding bodies be initially move with velocities \( u_1 \) and \( u_2 \). After collision both bodies stick together and moves with common velocity \( v \).

Total momentum of the system before collision = \( m_A u_1 + m_B u_2 \)

Total momentum of the system after collision =

mass of the composite body \( \times \) common velocity = \( (m_A + m_B) v \)

By law of conservation of momentum

\[
m_A u_1 + m_B u_2 = (m_A + m_B) v \quad \text{(or)} \quad v = \frac{m_A u_1 + m_B u_2}{m_A + m_B}
\]

Thus, knowing the masses of the two bodies and their velocities before collision, the common velocity of the system after collision can be calculated.

If the second particle is initially at rest i.e. \( u_2 = 0 \) then

\[
v = \frac{m_A u_1}{(m_A + m_B)}
\]

kinetic energy of the system before collision

\[
E_{K1} = \frac{1}{2} m_A u_1^2 \quad [\because u_2 = 0]
\]

and kinetic energy of the system after collision

\[
E_{K2} = \frac{1}{2} (m_A + m_B) v^2
\]

Hence,

\[
\frac{E_{K2}}{E_{K1}} = \frac{\text{kinetic energy after collision}}{\text{kinetic energy before collision}} = \frac{(m_A + m_B) v^2}{m_A u_1^2}
\]

Substituting the value of \( v \) in the above equation,

\[
\frac{E_{K2}}{E_{K1}} = \frac{m_A}{m_A + m_B} \quad \text{(or)} \quad \frac{E_{K2}}{E_{K1}} < 1
\]

It is clear from the above equation that in a perfectly inelastic collision, the kinetic energy after impact is less than the kinetic energy before impact. The loss in kinetic energy may appear as heat energy.
Solved Problems

2.1. The driver of a car travelling at 72 kmph observes the light 300 m ahead of him turning red. The traffic light is timed to remain red for 20 s before it turns green. If the motorist wishes to pass the light without stopping to wait for it to turn green, determine (i) the required uniform acceleration of the car (ii) the speed with which the motorist crosses the traffic light.

Data: \( u = 72 \text{ kmph} = 72 \times \frac{5}{18} \text{ m s}^{-1} = 20 \text{ m s}^{-1} \); \( S = 300 \text{ m} \); \( t = 20 \text{ s} \); \( a = ? \); \( v = ? \)

Solution:  
(i) \[ s = ut + \frac{1}{2} at^2 \]
\[ 300 = (20 \times 20) + \frac{1}{2} a (20)^2 \]
\[ a = -0.5 \text{ m s}^{-2} \]
(ii) \[ v = u + at = 20 - 0.5 \times 20 = 10 \text{ m s}^{-1} \]

2.2. A stone is dropped from the top of the tower 50 m high. At the same time another stone is thrown up from the foot of the tower with a velocity of 25 m s\(^{-1}\). At what distance from the top and after how much time the stones cross each other?

Data: Height of the tower = 50 m \( u_1 = 0 \); \( u_2 = 25 \text{ m s}^{-1} \)

Let \( s_1 \) and \( s_2 \) be the distances travelled by the two stones at the time of crossing \( t \). Therefore \( s_1 + s_2 = 50 \text{m} \)

\[ s_1 = ? ; t = ? \]

Solution:  
For I stone: \[ s_1 = \frac{1}{2} gt^2 \]

For II stone: \[ s_2 = u_2 t - \frac{1}{2} gt^2 \]
\[ s_2 = 25 t - \frac{1}{2} gt^2 \]

Therefore, \[ s_1 + s_2 = 50 = \frac{1}{2} gt^2 + 25 t - \frac{1}{2} gt^2 \]
\[ t = 2 \text{ seconds} \]
\[ s_1 = \frac{1}{2} gt^2 = \frac{1}{2} (9.8) (2)^2 = 19.6 \text{ m} \]
2.3 A boy throws a ball so that it may just clear a wall 3.6m high. The boy is at a distance of 4.8 m from the wall. The ball was found to hit the ground at a distance of 3.6m on the other side of the wall. Find the least velocity with which the ball can be thrown.

**Data:**  
Range of the ball = $4.8 + 3.6 = 8.4m$  
Height of the wall = 3.6m  
$u = ?$; $\theta = ?$

**Solution:** The top of the wall AC must lie on the path of the projectile.

The equation of the projectile is  
$$y = x\tan \theta - \frac{gx^2}{2u^2 \cos^2 \theta}$$  
...(1)

The point C ($x = 4.8m$, $y = 3.6m$) lies on the trajectory. Substituting the known values in (1),

$$3.6 = 4.8\tan \theta - \frac{g \times (4.8)^2}{2u^2 \cos^2 \theta}$$  
...(2)

The range of the projectile is  
$$R = \frac{u^2 \sin 2\theta}{g} = 8.4$$  
...(3)

From (3),  
$$\frac{u^2}{g} = \frac{8.4}{\sin 2\theta}$$  
...(4)

Substituting (4) in (2),

$$3.6 = (4.8)\tan \theta - \frac{(4.8)^2}{2 \cos^2 \theta} \times \frac{\sin 2\theta}{(8.4)}$$
\[ 3.6 = (4.8) \tan \theta - \frac{(4.8)^2}{2 \cos^2 \theta} \times \frac{2 \sin \theta \cos \theta}{(8.4)} \]

\[ 3.6 = (4.8) \tan \theta - (2.7429) \tan \theta \]

Substituting the value of \( \theta \) in (4),

\[ u^2 = \frac{8.4 \times g}{\sin 2\theta} = \frac{8.4 \times 9.8}{\sin 2(60^\circ 15')} = 95.5399 \]

\[ u = 9.7745 \text{ m s}^{-1} \]

2.4 Prove that for a given velocity of projection, the horizontal range is same for two angles of projection \( \alpha \) and \((90^\circ - \alpha)\).

The horizontal range is given by, \[ R = \frac{u^2 \sin 2\theta}{g} \] \[ R_1 = \frac{u^2 \sin 2\alpha}{g} \]

When \( \theta = \alpha \),

\[ R_1 = \frac{u^2 \sin 2\alpha}{g} \]

When \( \theta = (90^\circ - \alpha) \),

\[ \theta = \tan^{-1}(1.75) = 60^\circ 15' \]

\[ R_2 = \frac{u^2 \sin 2(90^\circ - \alpha)}{g} = \frac{u^2 \sin(90^\circ - \alpha) \cos(90^\circ - \alpha)}{g} \]

But \( \sin(90^\circ - \alpha) = \cos \alpha \); \( \cos(90^\circ - \alpha) = \sin \alpha \)

\[ \ldots(4) \]

From (2) and (4), it is seen that at both angles \( \alpha \) and \((90 - \alpha)\), the horizontal range remains the same.

2.5 The pilot of an aeroplane flying horizontally at a height of 2000 m with a constant speed of 540 kmph wishes to hit a target on the ground. At what distance from the target should release the bomb to hit the target?
Data: Initial velocity of the bomb in the horizontal is the same as that of the airplane.

Initial velocity of the bomb in the horizontal direction = 540 kmph = \( \frac{540 \times 5}{18} \) m s\(^{-1} \) = 150 m s\(^{-1} \)

Initial velocity in the vertical direction (u) = 0 ; vertical distance (s) = 2000 m ; time of flight t = ?

Solution: From equation of motion,

\[ s = ut + \frac{1}{2} at^2 \]

Substituting the known values,

\[ 2000 = 0 \times t + \frac{1}{2} \times 9.8 \times t^2 \]

\[ 2000 = 4.9t^2 \] (or)

\[ t = \frac{2000}{4.9} = 20.20 \text{ s} \]

∴ horizontal range = horizontal velocity \( \times \) time of flight

= 150 \( \times \) 20.20 = 3030 m

2.6 Two equal forces are acting at a point with an angle of 60° between them. If the resultant force is equal to 20\( \sqrt{3} \) N, find the magnitude of each force.

Data: Angle between the forces, \( \theta = 60^\circ \); Resultant R = 20\( \sqrt{3} \) N

\( P = Q = P \) (say) = ?

Solution: \[ R = \sqrt{P^2 + Q^2 + 2PQ \cos \theta} \]

\[ = \sqrt{P^2 + P^2 + 2P \cdot P \cos 60^\circ} \]

\[ = \sqrt{2P^2 + P^2 \cdot \frac{1}{2}} \]

\[ = P \sqrt{3} \]

\[ 20 \sqrt{3} = P \sqrt{3} \]

\[ P = 20 \text{ N} \]
2.7 If two forces \( F_1 = 20 \text{ kN} \) and \( F_2 = 15 \text{ kN} \) act on a particle as shown in figure, find their resultant by triangle law.

**Data:** \( F_1 = 20 \text{ kN} \); \( F_2 = 15 \text{ kN} \); \( R = ? \)

**Solution:** Using law of cosines,

\[
R^2 = P^2 + Q^2 - 2PQ \cos (180 - \theta)
\]

\[
R^2 = 20^2 + 15^2 - 2(20)(15) \cos 110^\circ
\]

\[
\therefore \quad R = 28.813 \text{ kN}.
\]

Using law of sines,

\[
\frac{R}{\sin 110^\circ} = \frac{15}{\sin \alpha}
\]

\[
\therefore \quad \alpha = 29.3^\circ
\]

2.8 Two forces act at a point in directions inclined to each other at 120°. If the bigger force is 5 kg wt and their resultant is at right angles to the smaller force, find the resultant and the smaller force.

**Data:** Bigger force = 5 kg wt

Angle made by the resultant with the smaller force = 90°

Resultant = ? Smaller force = ?

**Solution:** Let the forces \( P \) and \( Q \) are acting along \( OA \) and \( OD \) where \( \angle AOD = 120^\circ \)

Complete the parallelogram \( OACD \) and join \( OC \). \( OC \) therefore which represents the resultant which is perpendicular to \( OA \).

In \( \triangle OAC \)

\[
\angle OCA = \angle COD = 30^\circ
\]

\[
\angle AOC = 90^\circ
\]

Therefore \( \angle OAC = 60^\circ \)

\[
\text{(i.e)} \quad \frac{P}{\sin 30^\circ} = \frac{Q}{\sin 90^\circ} = \frac{R}{\sin 60^\circ}
\]

Since \( Q = 5 \text{ kg. wt.} \)

\[
P = \frac{5 \sin 30^\circ}{\sin 90^\circ} = 2.5 \text{ kg wt}
\]

\[
R = \frac{5 \sin 60^\circ}{\sin 90^\circ} = \frac{5 \sqrt{3}}{2} \text{ kg wt}
\]
2.9 Determine analytically the magnitude and direction of the resultant of the following four forces acting at a point.

(i) 10 kN pull N 30° E;  (ii) 20 kN push S 45° W;
(iii) 5 kN push N 60° W;  (iv) 15 kN push S 60° E.

Data:  \( F_1 = 10 \) kN;
\( F_2 = 20 \) kN;
\( F_3 = 5 \) kN;
\( F_4 = 15 \) kN;
\( R = \) ?;  \( \alpha = \) ?

Solution: The various forces acting at a point are shown in figure.

Resolving the forces horizontally, we get

\[
\Sigma F_x = 10 \sin 30° + 5 \sin 60° + 20 \sin 45° - 15 \sin 60° \\
= 10.48 \text{ kN}
\]

Similarly, resolving forces vertically, we get

\[
\Sigma F_y = 10 \cos 30° - 5 \cos 60° + 20 \cos 45° + 15 \cos 60° \\
= 27.8 \text{ kN}
\]

Resultant \( R = \sqrt{(\Sigma F_x)^2 + (\Sigma F_y)^2} \)

\[
= \sqrt{(10.48)^2 + (27.8)^2} \\
= 29.7 \text{ kN}
\]

\[
\tan \alpha = \frac{\Sigma F_y}{\Sigma F_x} = \frac{27.8}{10.48} = 2.65
\]

\( \alpha = 69.34° \)

2.10 A machine weighing 1500 N is supported by two chains attached to some point on the machine. One of these ropes goes to a nail in the wall and is inclined at 30° to the horizontal and
other goes to the hook in ceiling and is inclined at 45° to the horizontal. Find the tensions in the two chains.

**Data :** \( W = 1500 \text{ N} \). Tensions in the strings = ?

**Solution :** The machine is in equilibrium under the following forces:

(i) \( W \) (weight of the machine) acting vertically down;

(ii) Tension \( T_1 \) in the chain OA;

(iii) Tension \( T_2 \) in the chain OB.

Now applying Lami’s theorem at O, we get

\[
\frac{T_1}{\sin (90^\circ + 45^\circ)} = \frac{T_2}{\sin (90^\circ + 30^\circ)} = \frac{T_3}{\sin 105^\circ}
\]

\[
\frac{T_1}{\sin 135^\circ} = \frac{T_2}{\sin 120^\circ} = \frac{1500}{\sin 105^\circ}
\]

\[
T_1 = \frac{1500 \times \sin 135^\circ}{\sin 105^\circ} = 1098.96 \text{ N}
\]

\[
T_2 = \frac{1500 \times \sin 120^\circ}{\sin 105^\circ} = 1346.11 \text{ N}
\]

2.11 The radius of curvature of a railway line at a place when a train is moving with a speed of 72 kmph is 1500 m. If the distance between the rails is 1.54 m, find the elevation of the outer rail above the inner rail so that there is no side pressure on the rails.

**Data :** \( r = 1500 \text{ m} ; \ v = 72 \text{ kmph} = 20 \text{ m s}^{-1} ; \ l = 1.54 \text{ m} ; \ h = ? \)

**Solution :** \( \tan \theta = \frac{h}{l} = \frac{v^2}{rg} \)

Therefore \( h = \frac{lv^2}{rg} = \frac{1.54 \times (20)^2}{1500 \times 9.8} = 0.0419 \text{ m} \)
2.12 A truck of weight 2 tonnes is slipped from a train travelling at 9 kmph and comes to rest in 2 minutes. Find the retarding force on the truck.

**Data:**

- \( m = 2 \) tonne = 2 \( \times 1000 \) kg = 2000 kg
- \( v_1 = 9 \) kmph = \( 9 \times \frac{5}{18} \) = \( \frac{5}{2} \) m s\(^{-1} \);
- \( v_2 = 0 \)

**Solution:** Let \( R \) newton be the retarding force.

By the momentum - impulse theorem,

\[ (m v_1 - m v_2) = Rt \quad \text{(or)} \quad m v_1 - Rt = m v_2 \]

\[ 2000 \times \frac{5}{2} - R \times 120 = 2000 \times 0 \quad \text{(or)} \quad 5000 - 120 R = 0 \]

\[ R = 41.67 \text{ N} \]

2.13 A body of mass 2 kg initially at rest is moved by a horizontal force of 0.5N on a smooth frictionless table. Obtain the work done by the force in 8 s and show that this is equal to change in kinetic energy of the body.

**Data:**

- \( M = 2 \) kg
- \( F = 0.5 \) N
- \( t = 8 \) s
- \( W = ? \)

**Solution:**

\[ \text{Acceleration produced (a)} = \frac{F}{m} = \frac{0.5}{2} = 0.25 \text{ m s}\(^{-2} \) \]

The velocity of the body after 8s = \( a \times t = 0.25 \times 8 = 2 \) m \( \text{s}^{-1} \)

The distance covered by the body in 8 s = \( S = ut + \frac{1}{2} at^2 \)

\[ S = (0 \times 8) + \frac{1}{2} (0.25) (8)^2 = 8 \text{ m} \]

\[ \therefore \text{ Work done by the force in 8 s = } Force \times Distance = 0.5 \times 8 = 4 \text{ J} \]

**Initial kinetic energy** = \( \frac{1}{2} m (0)^2 = 0 \)

**Final kinetic energy** = \( \frac{1}{2} \times m v^2 = \frac{1}{2} \times 2 \times (2)^2 = 4 \text{ J} \)

\[ \therefore \text{ Change in kinetic energy = Final K.E} - \text{ Initial K.E} = 4 - 0 = 4 \text{ J} \]

The work done is equal to the change in kinetic energy of the body.
2.14 A body is thrown vertically up from the ground with a velocity of 39.2 m s\(^{-1}\). At what height will its kinetic energy be reduced to one-fourth of its original kinetic energy.

**Data:** \(v = 39.2 \text{ m s}^{-1} ; h = ?\)

**Solution:** When the body is thrown up, its velocity decreases and hence potential energy increases.

Let \(h\) be the height at which the potential energy is reduced to one-fourth of its initial value.

(i.e) loss in kinetic energy = gain in potential energy

\[
\frac{3}{4} \times \frac{1}{2} m v^2 = mg h
\]

\[
\frac{3}{4} \times \frac{1}{2} (39.2)^2 = 9.8 \times h
\]

\(h = 58.8\ \text{m}\)

2.15 A 10 g bullet is fired from a rifle horizontally into a 5 kg block of wood suspended by a string and the bullet gets embedded in the block. The impact causes the block to swing to a height of 5 cm above its initial level. Calculate the initial velocity of the bullet.

**Data:** Mass of the bullet = \(m_A = 10\ g = 0.01\ \text{kg}\)

Mass of the wooden block = \(m_B = 5\ \text{kg}\)

Initial velocity of the bullet before impact = \(u_A = ?\)

Initial velocity of the block before impact = \(u_B = 0\)

Final velocity of the bullet and block = \(v\)
**Solution :** By law of conservation of linear momentum,

\[ m_Au_A + m_Bu_B = (m_A + m_B) \cdot v \]

\[ (0.01)u_A + (5 \times 0) = (0.01 + 5) \cdot v \]

\[ (or) \quad v = \left( \frac{0.01}{5.01} \right) u_A = \frac{u_A}{501} \quad \ldots (1) \]

Applying the law of conservation of mechanical energy.

KE of the combined mass = PE at the highest point

\[ (or) \quad \frac{1}{2} (m_A + m_B) \cdot v^2 = (m_A + m_B) \cdot gh \quad \ldots (2) \]

From equation (1) and (2),

\[ \frac{u_A^2}{(501)^2} = 2gh \quad (or) \quad u_A = \sqrt{2.46 \times 10^3} = 496.0 \text{ m s}^{-1} \]
Self evaluation

(The questions and problems given in this self evaluation are only samples. In the same way any question and problem could be framed from the text matter. Students must be prepared to answer any question and problem from the text matter, not only from the self evaluation.)

2.1 A particle at rest starts moving in a horizontal straight line with uniform acceleration. The ratio of the distance covered during the fourth and the third second is

(a) \( \frac{4}{3} \)  
(b) \( \frac{26}{9} \)
(c) \( \frac{7}{5} \)  
(d) 2

2.2 The distance travelled by a body, falling freely from rest in one, two and three seconds are in the ratio

(a) 1 : 2 : 3  
(b) 1 : 3 : 5  
(c) 1 : 4 : 9  
(d) 9 : 4 : 1

2.3 The displacement of the particle along a straight line at time \( t \) is given by, \( x = a_0 + a_1 t + a_2 t^2 \) where \( a_0, a_1 \) and \( a_2 \) are constants. The acceleration of the particle is

(a) \( a_0 \)  
(b) \( a_1 \)  
(c) \( a_2 \)  
(d) 2\( a_2 \)

2.4 The acceleration of a moving body can be found from:

(a) area under velocity-time graph  
(b) area under distance-time graph  
(c) slope of the velocity-time graph  
(d) slope of the distance-time graph

2.5 Which of the following is a vector quantity?

(a) Distance  
(b) Temperature  
(c) Mass  
(d) Momentum

2.6 An object is thrown along a direction inclined at an angle 45° with the horizontal. The horizontal range of the object is

(a) vertical height  
(b) twice the vertical height  
(c) thrice the vertical height  
(d) four times the vertical height

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2.7. Two bullets are fired at angle \( \theta \) and \((90 - \theta)\) to the horizontal with some speed. The ratio of their times of flight is
(a) 1:1  
(b) \( \tan \theta :1 \)
(c) 1: \( \tan \theta \)  
(d) \( \tan^2 \theta :1 \)

2.8 A stone is dropped from the window of a train moving along a horizontal straight track, the path of the stone as observed by an observer on ground is
(a) Straight line  
(b) Parabola
(c) Circular  
(d) Hyperbola

2.9 A gun fires two bullets with same velocity at 60° and 30° with horizontal. The bullets strike at the same horizontal distance. The ratio of maximum height for the two bullets is in the ratio
(a) 2 : 1  
(b) 3 : 1
(c) 4 : 1  
(d) 1 : 1

2.10 Newton’s first law of motion gives the concept of
(a) energy  
(b) work
(c) momentum  
(d) Inertia

2.11 Inertia of a body has direct dependence on
(a) Velocity  
(b) Mass
(c) Area  
(d) Volume

2.12 The working of a rocket is based on
(a) Newton’s first law of motion  
(b) Newton’s second law of motion
(c) Newton’s third law of motion  
(d) Newton’s first and second law

2.13 When three forces acting at a point are in equilibrium
(a) each force is equal to the vector sum of the other two forces.
(b) each force is greater than the sum of the other two forces.
(c) each force is greater than the difference of the other two force.
(d) each force is to product of the other two forces.

2.14 For a particle revolving in a circular path, the acceleration of the particle is
(a) along the tangent
(b) along the radius
(c) along the circumference of the circle
(d) Zero

2.15 If a particle travels in a circle, covering equal angles in equal times, its velocity vector
(a) changes in magnitude only
(b) remains constant
(c) changes in direction only
(d) changes both in magnitude and direction

2.16 A particle moves along a circular path under the action of a force. The work done by the force is
(a) positive and nonzero (b) Zero
(c) Negative and nonzero (d) None of the above

2.17 A cyclist of mass \( m \) is taking a circular turn of radius \( R \) on a frictional level road with a velocity \( v \). In order that the cyclist does not skid,
(a) \( (mv^2/2) > \mu mg \) (b) \( (mu^2/r) > \mu mg \)
(c) \( (mu^2/r) < \mu mg \) (d) \( (v/r) = \mu g \)

2.18 If a force \( F \) is applied on a body and the body moves with velocity \( v \), the power will be
(a) \( Fv \) (b) \( F/v \)
(c) \( Fv^2 \) (d) \( F/v^2 \)

2.19 For an elastic collision
(a) the kinetic energy first increases and then decreases
(b) final kinetic energy never remains constant
(c) final kinetic energy is less than the initial kinetic energy
(d) initial kinetic energy is equal to the final kinetic energy

2.20 A bullet hits and gets embedded in a solid block resting on a horizontal frictionless table. Which of the following is conserved?
(a) momentum and kinetic energy
(b) Kinetic energy alone
(c) Momentum alone
(d) Potential energy alone

2.21 Compute the (i) distance travelled and (ii) displacement made by the student when he travels a distance of 4 km eastwards and then a further distance of 3 km northwards.

2.22 What is the (i) distance travelled and (ii) displacement produced by a cyclist when he completes one revolution?

2.23 Differentiate between speed and velocity of a body.

2.24 What is meant by retardation?

2.25 What is the significance of velocity-time graph?

2.26 Derive the equations of motion for an uniformly accelerated body.

2.27 What are scalar and vector quantities?

2.28 How will you represent a vector quantity?

2.29 What is the magnitude and direction of the resultant of two vectors acting along the same line in the same direction?

2.30 State: Parallelogram law of vectors and triangle law of vectors.

2.31 Obtain the expression for magnitude and direction of the resultant of two vectors when they are inclined at an angle ‘θ’ with each other.

2.32 State Newton’s laws of motion.

2.33 Explain the different types of inertia with examples.

2.34 State and prove law of conservation of linear momentum.

2.35 Define impulse of a force

2.36 Obtain an expression for centripetal acceleration.

2.37 What is centrifugal reaction?
2.38 Obtain an expression for the critical velocity of a body revolving in a vertical circle.

2.39 What is meant by banking of tracks?

2.40 Obtain an expression for the angle of lean when a cyclist takes a curved path.

2.41 What are the two types of collision? Explain them.

2.42 Obtain the expressions for the velocities of the two bodies after collision in the case of one dimensional motion.

2.43 Prove that in the case of one dimensional elastic collision between two bodies of equal masses, they interchange their velocities after collision.

Problems

2.44 Determine the initial velocity and acceleration of particle travelling with uniform acceleration in a straight line if it travels 55 m in the 8th second and 85 m in the 13th second of its motion.

2.45 An aeroplane takes off at an angle of 45° to the horizontal. If the vertical component of its velocity is 300 kmph, calculate its actual velocity. What is the horizontal component of velocity?

2.46 A force is inclined at 60° to the horizontal. If the horizontal component of force is 40 kg wt, calculate the vertical component.

2.47 A body is projected upwards with a velocity of 30 m s⁻¹ at an angle of 30° with the horizontal. Determine (a) the time of flight (b) the range of the body and (c) the maximum height attained by the body.

2.48 The horizontal range of a projectile is \(4\sqrt{3}\) times its maximum height. Find the angle of projection.

2.49 A body is projected at such an angle that the horizontal range is 3 times the greatest height. Find the angle of projection.

2.50 An elevator is required to lift a body of mass 65 kg. Find the acceleration of the elevator, which could cause a reaction of 800 N on the floor.

2.51 A body whose mass is 6 kg is acted on by a force which changes its velocity from 3 m s⁻¹ to 5 m s⁻¹. Find the impulse of the
force. If the force is acted for 2 seconds, find the force in newton.

2.52 A cricket ball of mass 150 g moving at 36 m s\(^{-1}\) strikes a bat and returns back along the same line at 21 m s\(^{-1}\). What is the change in momentum produced? If the bat remains in contact with the ball for 1/20 s, what is the average force exerted in newton.

2.53 Two forces of magnitude 12 N and 8 N are acting at a point. If the angle between the two forces is 60°, determine the magnitude of the resultant force?

2.54 The sum of two forces inclined to each other at an angle is 18 kg wt and their resultant which is perpendicular to the smaller force is 12 kg wt. Find the forces and the angle between them.

2.55 A weight of 20 kN supported by two cords, one 3 m long and the other 4 m long with points of support 5 m apart. Find the tensions \(T_1\) and \(T_2\) in the cords.

2.56 The following forces act at a point
(i) 20 N inclined at 30° towards North of East
(ii) 25 N towards North
(iii) 30 N inclined at 45° towards North of West
(iv) 35 N inclined at 40° towards South of West.
Find the magnitude and direction of the resultant force.

2.57 Find the magnitude of the two forces such that if they are at right angles, their resultant is \(\sqrt{10}\) N. But if they act at 60°, their resultant is \(\sqrt{13}\) N.

2.58 At what angle must a railway track with a bend of radius 880 m be banked for the safe running of a train at a velocity of 44 m s\(^{-1}\)?

2.59 A railway engine of mass 60 tonnes, is moving in an arc of radius 200 m with a velocity of 36 kmph. Find the force exerted on the rails towards the centre of the circle.

2.60 A horse pulling a cart exerts a steady horizontal pull of 300 N
and walks at the rate of 4.5 kmph. How much work is done by
the horse in 5 minutes?

2.61 A ball is thrown downward from a height of 30 m with a velocity
of 10 m s\(^{-1}\). Determine the velocity with which the ball strikes
the ground by using law of conservation of energy.

2.62 What is the work done by a man in carrying a suitcase weighing
30 kg over his head, when he travels a distance of 10 m in
(i) vertical and (ii) horizontal directions?

2.63 Two masses of 2 kg and 5 kg are moving with equal kinetic
energies. Find the ratio of magnitudes of respective linear
momenta.

2.64 A man weighing 60 kg runs up a flight of stairs 3 m high in 4 s.
Calculate the power developed by him.

2.65 A motor boat moves at a steady speed of 8 m s\(^{-1}\). If the water
resistance to the motion of the boat is 2000 N, calculate the
power of the engine.

2.66 Two blocks of mass 300 kg and 200 kg are moving toward each
other along a horizontal frictionless surface with velocities of 50
m s\(^{-1}\) and 100 m s\(^{-1}\) respectively. Find the final velocity of each
block if the collision is completely elastic.
Answers

2.1 (c) 2.2 (c) 2.3 (d) 2.4 (c)
2.5 (d) 2.6 (d) 2.7 (b) 2.8 (b)
2.9 (b) 2.10 (d) 2.11 (b) 2.12 (c)
2.13 (a) 2.14 (b) 2.15 (c) 2.16 (b)
2.17 (c) 2.18 (a) 2.19 (d) 2.20 (c)

2.44 10 m s\(^{-1}\); 6 m s\(^{-2}\)
2.45 424.26 kmph; 300 kmph

2.46 69.28 kg wt
2.47 3.06s; 79.53 m; 11.48 m

2.48 30°
2.49 53°7’

2.50 2.5 m s\(^{-2}\)
2.51 12 N s; 6 N

2.52 8.55 kg m s\(^{-1}\); 171 N
2.53 17.43 N

2.54 5 kg wt; 13 kg wt; 112°37’
2.55 16 k N, 12 k N

2.56 45.6 N; 132° 18’
2.57 3 N; 1 N

2.58 12°39’
2.59 30 kN

2.60 \(1.125 \times 10^5\) J
2.61 26.23 m s\(^{-1}\)

2.62 2940 J; 0
2.63 0.6324

2.64 441 W
2.65 16000 W

2.66 – 70 m s\(^{-1}\); 80 m s\(^{-1}\)
3. Dynamics of Rotational Motion

3.1 Centre of mass

Every body is a collection of large number of tiny particles. In translatory motion of a body, every particle experiences equal displacement with time; therefore the motion of the whole body may be represented by a particle. But when the body rotates or vibrates during translatory motion, then its motion can be represented by a point on the body that moves in the same way as that of a single particle subjected to the same external forces would move. A point in the system at which whole mass of the body is supposed to be concentrated is called centre of mass of the body. Therefore, if a system contains two or more particles, its translatory motion can be described by the motion of the centre of mass of the system.

3.1.1 Centre of mass of a two-particle system

Let us consider a system consisting of two particles of masses \(m_1\) and \(m_2\). \(P_1\) and \(P_2\) are their positions at time \(t\) and \(r_1\) and \(r_2\) are the corresponding distances from the origin O as shown in Fig. 3.1. Then the velocity and acceleration of the particles are,

\[v_1 = \frac{dr_1}{dt}\]  
\[a_1 = \frac{dv_1}{dt}\]  
\[v_2 = \frac{dr_2}{dt}\]  
\[a_2 = \frac{dv_2}{dt}\]

The particle at \(P_1\) experiences two forces:

(i) a force \(F_{12}\) due to the particle at \(P_2\) and

(ii) force \(F_{1e}\), the external force due to some particles external to the system.

If \(F_1\) is the resultant of these two forces,
\[ F_1 = F_{12} + F_{1e} \]...(5)

Similarly, the net force \( F_2 \) acting on the particle \( P_2 \) is,
\[ F_2 = F_{21} + F_{2e} \]...(6)

where \( F_{21} \) is the force exerted by the particle at \( P_1 \) on \( P_2 \).

By using Newton’s second law of motion,
\[ F_1 = m_1a_1 \]...(7)

and \[ F_2 = m_2a_2 \]...(8)

Adding equations (7) and (8), \( m_1a_1 + m_2a_2 = F_1 + F_2 \)

Substituting \( F_1 \) and \( F_2 \) from (5) and (6)
\[ m_1a_1 + m_2a_2 = F_{12} + F_{1e} + F_{21} + F_{2e} \]

By Newton’s third law, the internal force \( F_{12} \) exerted by particle at \( P_2 \) on the particle at \( P_1 \) is equal and opposite to \( F_{21} \), the force exerted by particle at \( P_1 \) on \( P_2 \).

\[ \therefore F_{12} = -F_{21} \]...(9)

\[ \therefore F = F_{1e} + F_{2e} \]...(10)

where \( F \) is the net external force acting on the system.

The total mass of the system is given by,
\[ M = m_1 + m_2 \]...(11)

Let the net external force \( F \) acting on the system produces an acceleration \( a_{CM} \) called the acceleration of the centre of mass of the system.

By Newton’s second law, for the system of two particles,
\[ F = Ma_{CM} \]...(12)

From (10) and (12), \( Ma_{CM} = m_1a_1 + m_2a_2 \) ...(13)

Let \( R_{CM} \) be the position vector of the centre of mass.

\[ a_{CM} = \frac{d^2(R_{CM})}{dt^2} \]...(14)

From (13) and (14),
\[ \frac{d^2R_{xy}}{dt^2} = \left( \frac{1}{M} \right) \left( m_1 \frac{d^2r_1}{dt^2} + m_2 \frac{d^2r_2}{dt^2} \right) \]
\[
\frac{d^2 R_{CM}}{dt^2} = \frac{1}{M} \left( \frac{d^2}{dt^2} (m_1 r_1 + m_2 r_2) \right)
\]
\[
\therefore R_{CM} = \frac{1}{M} \left( m_1 r_1 + m_2 r_2 \right)
\]
\[
R_{CM} = \frac{m_1 r_1 + m_2 r_2}{m_1 + m_2} \quad \ldots \ldots \ldots (15)
\]

This equation gives the position of the centre of mass of a system comprising two particles of masses \(m_1\) and \(m_2\).

If the masses are equal \((m_1 = m_2)\), then the position vector of the centre of mass is,
\[
R_{CM} = \frac{r_1 + r_2}{2} \quad \ldots \ldots \ldots (16)
\]
which means that the centre of mass lies exactly in the middle of the line joining the two masses.

### 3.1.2 Centre of mass of a body consisting of \(n\) particles

For a system consisting of \(n\) particles with masses \(m_1, m_2, m_3 \ldots m_n\) with position vectors \(r_1, r_2, r_3 \ldots r_n\), the total mass of the system is,
\[
M = m_1 + m_2 + m_3 + \ldots \ldots + m_n
\]

The position vector \(R_{CM}\) of the centre of mass with respect to origin \(O\) is given by
\[
R_{CM} = \frac{\sum_{i=1}^{n} m_i r_i}{m_1 + m_2 + \ldots + m_n} = \frac{\sum_{i=1}^{n} m_i r_i}{M}
\]

The \(x\) coordinate and \(y\) coordinate of the centre of mass of the system are
\[
x = \frac{m_1 x_1 + m_2 x_2 + \ldots + m_n x_n}{m_1 + m_2 + \ldots + m_n} \quad \text{and} \quad y = \frac{m_1 y_1 + m_2 y_2 + \ldots + m_n y_n}{m_1 + m_2 + \ldots + m_n}
\]

### Example for motion of centre of mass

Let us consider the motion of the centre of mass of the Earth and moon system (Fig 3.2). The moon moves round the Earth in a circular
orbit and the Earth moves round the Sun in an elliptical orbit. It is more correct to say that the Earth and the moon both move in circular orbits about their common centre of mass in an elliptical orbit round the Sun.

For the system consisting of the Earth and the moon, their mutual gravitational attractions are the internal forces in the system and Sun’s attraction on both the Earth and moon are the external forces acting on the centre of mass of the system.

### 3.1.3 Centre of Gravity

A body may be considered to be made up of an indefinitely large number of particles, each of which is attracted towards the centre of the Earth by the force of gravity. These forces constitute a system of like parallel forces. The resultant of these parallel forces known as the weight of the body always acts through a point, which is fixed relative to the body, whatever be the position of the body. This fixed point is called the centre of gravity of the body.

The centre of gravity of a body is the point at which the resultant of the weights of all the particles of the body acts, whatever may be the orientation or position of the body provided that its size and shape remain unaltered.

In the Fig. 3.3, \( W_1, W_2, W_3, \ldots \) are the weights of the first, second, third, ... particles in the body respectively. If \( W \) is the resultant weight of all the particles then the point at which \( W \) acts is known as the centre of gravity. The total weight of the body may be supposed to act at its centre of gravity. Since the weights of the particles constituting a body are practically proportional to their masses when the body is outside the Earth and near its surface, the centre of mass of a body practically coincides with its centre of gravity.
3.1.4 *Equilibrium of bodies and types of equilibrium*

If a marble $M$ is placed on a curved surface of a bowl $S$, it rolls down and settles in equilibrium at the lowest point $A$ (Fig. 3.4 a). This equilibrium position corresponds to minimum potential energy. If the marble is disturbed and displaced to a point $B$, its energy increases. When it is released, the marble rolls back to $A$. Thus the marble at the position $A$ is said to be in *stable equilibrium*.

Suppose now that the bowl $S$ is inverted and the marble is placed at its top point, at $A$ (Fig. 3.4b). If the marble is displaced slightly to the point $C$, its potential energy is lowered and tends to move further away from the equilibrium position to one of lowest energy. Thus the marble is said to be in *unstable equilibrium*.

Suppose now that the marble is placed on a plane surface (Fig. 3.4c). If it is displaced slightly, its potential energy does not change. Here the marble is said to be in *neutral equilibrium*.

Equilibrium is thus *stable*, *unstable* or neutral according to whether the potential energy is *minimum*, *maximum* or *constant*.

We may also characterize the stability of a mechanical system by noting that when the system is disturbed from its position of equilibrium, the forces acting on the system may

(i) tend to bring back to its original position if potential energy is a *minimum*, corresponding to *stable equilibrium*.

(ii) tend to move it farther away if potential energy is *maximum*, corresponding *unstable equilibrium*.

(iii) tend to move either way if potential energy is a *constant* corresponding to *neutral equilibrium*.
Consider three uniform bars shown in Fig. 3.5 a,b,c. Suppose each bar is slightly displaced from its position of equilibrium and then released. For bar A, fixed at its top end, its centre of gravity \( G \) rises to \( G_1 \) on being displaced, then the bar returns back to its original position on being released, so that the equilibrium is stable.

For bar B, whose fixed end is at its bottom, its centre of gravity \( G \) is lowered to \( G_2 \) on being displaced, then the bar B will keep moving away from its original position on being released, and the equilibrium is said to be unstable.

For bar C, whose fixed point is about its centre of gravity, the centre of gravity remains at the same height on being displaced, the bar will remain in its new position, on being released, and the equilibrium is said to be neutral.

### 3.2 Rotational motion of rigid bodies

#### 3.2.1 Rigid body

A rigid body is defined as that body which does not undergo any change in shape or volume when external forces are applied on it. When forces are applied on a rigid body, the distance between any two particles of the body will remain unchanged, however, large the forces may be.

Actually, no body is perfectly rigid. Every body can be deformed more or less by the application of the external force. The solids, in which the changes produced by external forces are negligibly small, are usually considered as rigid body.
3.2.2 Rotational motion

When a body rotates about a fixed axis, its motion is known as rotatory motion. A rigid body is said to have pure rotational motion, if every particle of the body moves in a circle, the centre of which lies on a straight line called the axis of rotation (Fig. 3.6). The axis of rotation may lie inside the body or even outside the body. The particles lying on the axis of rotation remains stationary.

The position of particles moving in a circular path is conveniently described in terms of a radius vector \( r \) and its angular displacement \( \theta \). Let us consider a rigid body that rotates about a fixed axis XOX’ passing through O and perpendicular to the plane of the paper as shown in Fig 3.7. Let the body rotate from the position A to the position B. The different particles at \( P_1, P_2, P_3, \ldots \) in the rigid body covers unequal distances \( P_1P_1', P_2P_2', P_3P_3' \ldots \) in the same interval of time. Thus their linear velocities are different. But in the same time interval, they all rotate through the same angle \( \theta \) and hence the angular velocity is the same for all the particles of the rigid body. Thus, in the case of rotational motion, different constituent particles have different linear velocities but all of them have the same angular velocity.

3.2.3 Equations of rotational motion

As in linear motion, for a body having uniform angular acceleration, we shall derive the equations of motion.

Let us consider a particle start rotating with angular velocity \( \omega_0 \) and angular acceleration \( \alpha \). At any instant \( t \), let \( \omega \) be the angular velocity of the particle and \( \theta \) be the angular displacement produced by the particle.

Therefore change in angular velocity in time \( t = \omega - \omega_0 \)

But, angular acceleration \( = \frac{\text{change in angular velocity}}{\text{time taken}} \)
(i.e) \[ \alpha = \frac{\omega - \omega_0}{t} \] ...{(1)}

\[ \omega = \omega_0 + \alpha t \] ...{(2)}

The average angular velocity = \( \frac{\omega + \omega_0}{2} \)

The total angular displacement = average angular velocity \times time taken

(i.e) \[ \theta = \left( \frac{\omega + \omega_0}{2} \right) t \] ...{(3)}

Substituting \( \omega \) from equation (2),

\[ \theta = \omega_0 t + \frac{1}{2} \alpha t^2 \] ...{(4)}

From equation (1), \( t = \left( \frac{\omega - \omega_0}{\alpha} \right) \) ...{(5)}

using equation (5) in (3),

\[ \theta = \left( \frac{\omega + \omega_0}{2} \right) \left( \frac{\omega - \omega_0}{\alpha} \right) = \frac{\omega^2 - \omega_0^2}{2\alpha} \]

\[ 2\alpha \theta = \omega^2 - \omega_0^2 \quad \text{or} \quad \omega^2 = \omega_0^2 + 2\alpha \theta \] ...{(6)}

Equations (2), (4) and (6) are the equations of rotational motion.

3.3 Moment of inertia and its physical significance

According to Newton’s first law of motion, a body must continue in its state of rest or of uniform motion unless it is compelled by some external agency called force. The inability of a material body to change its state of rest or of uniform motion by itself is called inertia. Inertia is the fundamental property of the matter. For a given force, the greater the mass, the higher will be the opposition for motion, or larger the inertia. Thus, in translatory motion, the mass of the body measures the coefficient of inertia.

Similarly, in rotational motion also, a body, which is free to rotate about a given axis, opposes any change desired to be produced in its state. The measure of opposition will depend on the mass of the body
and the distribution of mass about the axis of rotation. The coefficient of inertia in rotational motion is called the moment of inertia of the body about the given axis.

Moment of inertia plays the same role in rotational motion as that of mass in translatory motion. Also, to bring about a change in the state of rotation, torque has to be applied.

### 3.3.1 Rotational kinetic energy and moment of inertia of a rigid body

Consider a rigid body rotating with angular velocity \( \omega \) about an axis \( \text{XOX}' \). Consider the particles of masses \( m_1, m_2, m_3 \ldots \) situated at distances \( r_1, r_2, r_3 \ldots \) respectively from the axis of rotation. The angular velocity of all the particles is same but the particles rotate with different linear velocities. Let the linear velocities of the particles be \( v_1, v_2, v_3 \ldots \) respectively.

Kinetic energy of the first particle = \( \frac{1}{2} m_1 v_1^2 \)

But \( v_1 = r_1 \omega \)

\[ \therefore \text{Kinetic energy of the first particle} = \frac{1}{2} m_1 (r_1 \omega)^2 = \frac{1}{2} m_1 r_1^2 \omega^2 \]

Similarly,

Kinetic energy of second particle

\[ = \frac{1}{2} m_2 r_2^2 \omega^2 \]

Kinetic energy of third particle

\[ = \frac{1}{2} m_3 r_3^2 \omega^2 \text{ and so on.} \]

The kinetic energy of the rotating rigid body is equal to the sum of the kinetic energies of all the particles.

\[ \therefore \text{Rotational kinetic energy} \]

\[ = \frac{1}{2} (m_1 r_1^2 \omega^2 + m_2 r_2^2 \omega^2 + m_3 r_3^2 \omega^2 + \ldots + m_n r_n^2 \omega^2) \]

\[ = \frac{1}{2} \omega^2 (m_1 r_1^2 + m_2 r_2^2 + m_3 r_3^2 + \ldots + m_n r_n^2) \]
In translatory motion, kinetic energy $= \frac{1}{2} m v^2$

Comparing with the above equation, the inertial role is played by the term $\sum m_i r_i^2$. This is known as moment of inertia of the rotating rigid body about the axis of rotation. Therefore the moment of inertia is $I = \text{mass} \times (\text{distance})^2$

Kinetic energy of rotation $= \frac{1}{2} \omega^2 I$

When $\omega = 1 \text{ rad s}^{-1}$, rotational kinetic energy

$= E_R = \frac{1}{2} (1)^2 I$ \hspace{1cm} (or) \hspace{1cm} I = 2E_R

It shows that moment of inertia of a body is equal to twice the kinetic energy of a rotating body whose angular velocity is one radian per second.

The unit for moment of inertia is $\text{kg m}^2$ and the dimensional formula is $\text{ML}^2$.

**3.3.2 Radius of gyration**

The moment of inertia of the rotating rigid body is

$I = \sum m_i r_i^2 = m_1 r_1^2 + m_2 r_2^2 + \ldots m_n r_n^2$

If the particles of the rigid body are having same mass, then $m_1 = m_2 = m_3 = \ldots = m$ (say)

\[ \therefore \] The above equation becomes,

$I = m (r_1^2 + r_2^2 + r_3^2 + \ldots + r_n^2)$

$I = nm \left[ \frac{r_1^2 + r_2^2 + r_3^2 + \ldots + r_n^2}{n} \right]$

where $n$ is the number of particles in the rigid body.
\[ I = MK^2 \quad \ldots \quad (2) \]

where \( M = nm \), total mass of the body and \( K^2 = \frac{r_1^2 + r_2^2 + r_3^2 + \ldots + r_n^2}{n} \)

Here \( K = \sqrt{\frac{r_1^2 + r_2^2 + r_3^2 + \ldots + r_n^2}{n}} \) is called as the radius of gyration of the rigid body about the axis of rotation.

The radius of gyration is equal to the root mean square distances of the particles from the axis of rotation of the body.

The radius of gyration can also be defined as the perpendicular distance between the axis of rotation and the point where the whole weight of the body is to be concentrated.

Also from the equation (2) \( K^2 = \frac{I}{M} \) (or) \( K = \frac{\sqrt{I}}{M} \)

### 3.3.3 Theorems of moment of inertia

\( \text{(i) Parallel axes theorem} \)

**Statement**

The moment of inertia of a body about any axis is equal to the sum of its moment of inertia about a parallel axis through its centre of gravity and the product of the mass of the body and the square of the distance between the two axes.

**Proof**

Let us consider a body having its centre of gravity at G as shown in Fig. 3.9. The axis \( XX' \) passes through the centre of gravity and is perpendicular to the plane of the body. The axis \( X_1X_1' \) passes through the point O and is parallel to the axis \( XX' \). The distance between the two parallel axes is \( x \).

Let the body be divided into large number of particles each of mass \( m \). For a particle \( P \) at a distance \( r \) from O, its moment of inertia about the axis \( X_1OX_1' \) is equal to \( m r^2 \).

The moment of inertia of the whole body about the axis \( X_1X_1' \) is given by,

\[ I_o = \Sigma m r^2 \quad \ldots \quad (1) \]
From the point $P$, drop a perpendicular $PA$ to the extended $OG$ and join $PG$.

In the $\triangle OPA$,

\[ OP^2 = OA^2 + AP^2 \]
\[ r^2 = (x + h)^2 + AP^2 \]
\[ r^2 = x^2 + 2xh + h^2 + AP^2 \quad \ldots(2) \]

But from $\triangle GPA$,

\[ GP^2 = GA^2 + AP^2 \]
\[ y^2 = h^2 + AP^2 \quad \ldots(3) \]

Substituting equation (3) in (2),

\[ r^2 = x^2 + 2xh + y^2 \quad \ldots(4) \]

Substituting equation (4) in (1),

\[ I_o = \Sigma m \ (x^2 + 2xh + y^2) \]
\[ = \Sigma mx^2 + \Sigma 2m xh + \Sigma my^2 \]
\[ = Mx^2 + My^2 + 2x\Sigma mh \quad \ldots(5) \]

Here $My^2 = I_c$ is the moment of inertia of the body about the line passing through the centre of gravity. The sum of the turning moments of
all the particles about the centre of gravity is zero, since the body is balanced about the centre of gravity G.

\[ \Sigma (mg) (h) = 0 \quad \text{(or) } \Sigma mh = 0 \quad \text{[since } g \text{ is a constant]} \quad \ldots (6) \]

\[ \therefore \text{equation (5) becomes, } I_0 = Mx^2 + I_G \quad \ldots (7) \]

Thus the parallel axes theorem is proved.

(ii) **Perpendicular axes theorem**

**Statement**

The moment of inertia of a plane laminar body about an axis perpendicular to the plane is equal to the sum of the moments of inertia about two mutually perpendicular axes in the plane of the lamina such that the three mutually perpendicular axes have a common point of intersection.

**Proof**

Consider a plane lamina having the axes \(OX\) and \(OY\) in the plane of the lamina as shown Fig. 3.10. The axis \(OZ\) passes through \(O\) and is perpendicular to the plane of the lamina. Let the lamina be divided into a large number of particles, each of mass \(m\). A particle at \(P\) at a distance \(r\) from \(O\) has coordinates \((x,y)\).

\[ \therefore r^2 = x^2 + y^2 \quad \ldots (1) \]

The moment of inertia of the particle \(P\) about the axis \(OZ\) is

\[ I_z = \Sigma mr^2 \]

The moment of inertia of the whole lamina about the axis \(OZ\) is

\[ I_z = \Sigma mx^2 + \Sigma my^2 = I_y + I_x \]

\[ \therefore I_z = I_x + I_y \]

which proves the perpendicular axes theorem.
<table>
<thead>
<tr>
<th>Body</th>
<th>Axis of Rotation</th>
<th>Moment of Inertia</th>
</tr>
</thead>
<tbody>
<tr>
<td>Thin Uniform Rod</td>
<td>Axis passing through its centre of gravity and perpendicular to its length</td>
<td>$\frac{ML^2}{12}$ M - mass 1 - length</td>
</tr>
<tr>
<td></td>
<td>Axis passing through the end and perpendicular to its length.</td>
<td>$\frac{ML^2}{3}$ M - mass 1 - length</td>
</tr>
<tr>
<td>Thin Circular Ring</td>
<td>Axis passing through its centre and perpendicular to its plane.</td>
<td>$MR^2$ M - mass R - radius</td>
</tr>
<tr>
<td></td>
<td>Axis passing through its diameter</td>
<td>$\frac{1}{2}MR^2$ M - mass R - radius</td>
</tr>
<tr>
<td></td>
<td>Axis passing through a tangent</td>
<td>$\frac{3}{2}MR^2$ M - mass R - radius</td>
</tr>
<tr>
<td>Circular Disc</td>
<td>Axis passing through its centre and perpendicular to its plane.</td>
<td>$\frac{1}{2}MR^2$ M - mass R - radius</td>
</tr>
<tr>
<td></td>
<td>Axis passing through its diameter</td>
<td>$\frac{1}{4}MR^2$ M - mass R - radius</td>
</tr>
<tr>
<td></td>
<td>Axis passing through a tangent</td>
<td>$\frac{5}{4}MR^2$ M - mass R - radius</td>
</tr>
<tr>
<td>Solid Sphere</td>
<td>Axis passing through its diameter</td>
<td>$\frac{2}{5}MR^2$ M - mass R - radius</td>
</tr>
<tr>
<td></td>
<td>Axis passing through a tangent</td>
<td>$\frac{7}{5}MR^2$ M - mass R - radius</td>
</tr>
<tr>
<td>Solid Cylinder</td>
<td>Its own axis</td>
<td>$\frac{1}{2}MR^2$ M - mass R - radius</td>
</tr>
<tr>
<td></td>
<td>Axis passing through its centre and perpendicular to its length</td>
<td>$M \left( \frac{R^2}{4} + \frac{r^2}{12} \right)$ M - mass R - radius 1 - length</td>
</tr>
</tbody>
</table>

Table 3.1 Moment of Inertia of different bodies

(Proof is given in the annexure)
3.4 Moment of a force

A force can rotate a nut when applied by a wrench or it can open a door while the door rotates on its hinges (i.e.) in addition to the tendency to move a body in the direction of the application of a force, a force also tends to rotate the body about any axis which does not intersect the line of action of the force and also not parallel to it. This tendency of rotation is called turning effect of a force or moment of the force about the given axis. The magnitude of the moment of force \(F\) about a point is defined as the product of the magnitude of force and the perpendicular distance of the point from the line of action of the force.

Let us consider a force \(F\) acting at the point \(P\) on the body as shown in Fig. 3.11. Then, the moment of the force \(F\) about the point \(O\) = Magnitude of the force \(\times\) perpendicular distance between the direction of the force and the point about which moment is to be determined = \(F \times OA\).

If the force acting on a body rotates the body in anticlockwise direction with respect to \(O\) then the moment is called anticlockwise moment. On the other hand, if the force rotates the body in clockwise direction then the moment is said to be clockwise moment. The unit of moment of the force is N m and its dimensional formula is \(M L^2 T^{-2}\).

As a matter of convention, an anticlockwise moment is taken as positive and a clockwise moment as negative. While adding moments, the direction of each moment should be taken into account.

In terms of vector product, the moment of a force is expressed as,
\[
\mathbf{m} = \mathbf{r} \times \mathbf{F}
\]
where \(\mathbf{r}\) is the position vector with respect to \(O\). The direction of \(\mathbf{m}\) is perpendicular to the plane containing \(\mathbf{r}\) and \(\mathbf{F}\).
3.5 Couple and moment of the couple (Torque)

There are many examples in practice where two forces, acting together, exert a moment, or turning effect on some object. As a very simple case, suppose two strings are tied to a wheel at the points X and Y, and two equal and opposite forces, \( F \), are exerted tangentially to the wheels (Fig. 3.13). If the wheel is pivoted at its centre \( O \) it begins to rotate about \( O \) in an anticlockwise direction.

Two equal and opposite forces whose lines of action do not coincide are said to constitute a couple in mechanics. The two forces always have a turning effect, or moment, called a torque. The perpendicular distance between the lines of action of two forces, which constitute the couple, is called the arm of the couple.

The product of the forces forming the couple and the arm of the couple is called the moment of the couple or torque.

Torque = one of the forces × perpendicular distance between the forces

The torque in rotational motion plays the same role as the force in translational motion. A quantity that is a measure of this rotational effect produced by the force is called torque.

In vector notation, \( \tau = r \times F \)

The torque is maximum when \( \theta = 90^\circ \) (i.e) when the applied force is at right angles to \( r \).

Examples of couple are
1. Forces applied to the handle of a screw press,
2. Opening or closing a water tap,
3. Turning the cap of a pen,
4. Steering a car.

Work done by a couple

Suppose two equal and opposite forces \( F \) act tangentially to a wheel \( W \), and rotate it through an angle \( \theta \) (Fig. 3.14).
Then the work done by each force = Force $\times$ distance = $F \times r \theta$

(since $r \theta$ is the distance moved by a point on the rim)

Total work done $W = F r \theta + F r \theta = 2F r \theta$

but torque $\tau = F \times 2r = 2Fr$

$\therefore$ work done by the couple, $W = \tau \theta$

### 3.6 Angular momentum of a particle

The angular momentum in a rotational motion is similar to the linear momentum in translatory motion. The linear momentum of a particle moving along a straight line is the product of its mass and linear velocity (i.e) $p = m v$. The angular momentum of a particle is defined as the moment of linear momentum of the particle.

Let us consider a particle of mass $m$ moving in the $XY$ plane with a velocity $v$ and linear momentum $\vec{p} = m \vec{v}$ at a distance $r$ from the origin (Fig. 3.15).

The angular momentum $L$ of the particle about an axis passing through O perpendicular to XY plane is defined as the cross product of $\vec{r}$ and $\vec{p}$.

$$L = r \times \vec{p}$$

Its magnitude is given by $L = r p \sin \theta$

where $\theta$ is the angle between $\vec{r}$ and $\vec{p}$ and $L$ is along a direction perpendicular to the plane containing $\vec{r}$ and $\vec{p}$.

The unit of angular momentum is kg m$^2$ s$^{-1}$ and its dimensional formula is, M L$^2$ T$^{-1}$.

#### 3.6.1 Angular momentum of a rigid body

Let us consider a system of $n$ particles of masses $m_1, m_2 \ldots, m_n$ situated at distances $r_1, r_2, \ldots, r_n$ respectively from the axis of rotation (Fig. 3.16). Let $v_1, v_2, v_3 \ldots$ be the linear velocities of the particles respectively, then linear momentum of first particle $= m_1 v_1$. 

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Since $v_1 = r_1 \omega$, the linear momentum of first particle = $m_1 (r_1 \omega)$

The moment of linear momentum of first particle

= linear momentum $\times$ perpendicular distance

= $(m_1 r_1 \omega) \times r_1$

angular momentum of first particle = $m_1 r_1^2 \omega$

Similarly,

angular momentum of second particle = $m_2 r_2^2 \omega$
angular momentum of third particle = $m_3 r_3^2 \omega$ and so on.

The sum of the moment of the linear momenta of all the particles of a rotating rigid body taken together about the axis of rotation is known as angular momentum of the rigid body.

∴ Angular momentum of the rotating rigid body = sum of the angular momenta of all the particles.

(i.e) 

$$L = m_1 r_1^2 \omega + m_2 r_2^2 \omega + m_3 r_3^2 \omega + \ldots + m_n r_n^2 \omega$$

$$L = \omega \left[ \sum_{i=1}^{n} m_i r_i^2 \right]$$

∴ $L = \omega \ I$

where $I = \sum_{i=1}^{n} m_i r_i^2$ = moment of inertia of the rotating rigid body about the axis of rotation.

3.7 *Relation between torque and angular acceleration*

Let us consider a rigid body rotating about a fixed axis $XOX'$ with angular velocity $\omega$ (Fig. 3.17).

The force acting on a particle of mass $m_1$ situated at A, at a distance $r_1$, from the axis of rotation = mass $\times$ acceleration

$$= m_1 \times \frac{d}{dt} (r_1 \omega)$$
The moment of this force about the axis of rotation
= Force \times \text{perpendicular distance}
= m_1 r_1 \frac{d^2\theta}{dt^2} \times r_1

Therefore, the total moment of all the forces acting on all the particles
= m_1 r_1 \frac{d^2\theta}{dt^2} + m_2 r_2 \frac{d^2\theta}{dt^2} + ...

(i.e) torque = \sum_{i=1}^{n} m_i r_i^2 \frac{d^2\theta}{dt^2}

or \quad \tau = I\alpha

where \sum_{i=1}^{n} m_i r_i^2 = \text{moment of inertia} I of the rigid body and \alpha = \frac{d^2\theta}{dt^2} \text{angular acceleration.}

### 3.7.1 Relation between torque and angular momentum

The angular momentum of a rotating rigid body is, \( L = I\omega \)

Differentiating the above equation with respect to time,

\[
\frac{dL}{dt} = I \left( \frac{d\omega}{dt} \right) = I\alpha
\]

where \( \alpha = \frac{d\omega}{dt} \) angular acceleration of the body.

But torque \( \tau = I\alpha \)

Therefore, torque \( \tau = \frac{dL}{dt} \)

Thus the rate of change of angular momentum of a body is equal to the external torque acting upon the body.

### 3.8 Conservation of angular momentum

The angular momentum of a rotating rigid body is, \( L = I\omega \)

The torque acting on a rigid body is, \( \tau = \frac{dL}{dt} \)
When no external torque acts on the system, \( \tau = \frac{dL}{dt} = 0 \)

(i.e) \( L = I \omega = \text{constant} \)

Total angular momentum of the body = constant

(i.e.) when no external torque acts on the body, the net angular momentum of a rotating rigid body remains constant. This is known as law of conservation of angular momentum.

**Illustration of conservation of angular momentum**

From the law of conservation of angular momentum, \( I \omega = \text{constant} \)

(i.e) \( \omega \propto \frac{1}{I} \), the angular velocity of rotation is inversely proportional to the moment of inertia of the system.

Following are the examples for law of conservation of angular momentum.

1. A diver jumping from springboard sometimes exhibits somersaults in air before reaching the water surface, because the diver curls his body to decrease the moment of inertia and increase angular velocity. When he

---

*Fig. 3.18 A diver jumping from a spring board*
is about to reach the water surface, he again outstretches his limbs. This
again increases moment of inertia and decreases the angular velocity.
Hence, the diver enters the water surface with a gentle speed.

2. A ballet dancer can increase her angular velocity by folding her
arms, as this decreases the moment of inertia.

3. Fig. 3.19a shows a person sitting on a turntable holding a pair of
heavy dumbbells one in each hand with arms outstretched. The table is
rotating with a certain angular velocity. The person suddenly pushes the
weight towards his chest as shown Fig. 3.19b, the speed of rotation is
found to increase considerably.

4. The angular velocity of a planet in its orbit round the sun increases
when it is nearer to the Sun, as the moment of inertia of the planet about
the Sun decreases.
Solved Problems

3.1 A system consisting of two masses connected by a massless rod lies along the X-axis. A 0.4 kg mass is at a distance \( x = 2 \) m while a 0.6 kg mass is at \( x = 7 \) m. Find the \( x \) coordinate of the centre of mass.

**Data** : \( m_1 = 0.4 \) kg ; \( m_2 = 0.6 \) kg ; \( x_1 = 2 \) m ; \( x_2 = 7 \) m ; \( x = ? \)

**Solution** :

\[
x = \frac{m_1x_1 + m_2x_2}{m_1 + m_2} = \frac{(0.4 \times 2) + (0.6 \times 7)}{(0.4 + 0.6)} = 5 \text{ m}
\]

3.2 Locate the centre of mass of a system of bodies of masses \( m_1 = 1 \) kg, \( m_2 = 2 \) kg and \( m_3 = 3 \) kg situated at the corners of an equilateral triangle of side 1 m.

**Data** : \( m_1 = 1 \) kg ; \( m_2 = 2 \) kg ; \( m_3 = 3 \) kg ;

- The coordinates of \( A = (0,0) \)
- The coordinates of \( B = (1,0) \)
- Centre of mass of the system =?

**Solution** : Consider an equilateral triangle of side 1m as shown in Fig. Take X and Y axes as shown in figure.

To find the coordinate of C:

For an equilateral triangle , \( \angle CAB = 60^\circ \)

Consider the triangle ADC,

\[
\sin \theta = \frac{CD}{CA} \quad \text{(or) } CD =
\]

\[
(CA) \sin \theta = 1 \times \sin 60^\circ = \frac{\sqrt{3}}{2}
\]

Therefore from the figure, the coordinate of C are, \( (0.5, \frac{\sqrt{3}}{2}) \)

\[
x = \frac{m_1x_1 + m_2x_2 + m_3x_3}{m_1 + m_2 + m_3}
\]
A circular disc of mass \( m \) and radius \( r \) is set rolling on a table. If \( \omega \) is its angular velocity, show that its total energy \( E = \frac{3}{4} mr^2 \omega^2 \).

**Solution** : The total energy of the disc = Rotational KE + linear KE

\[
E = \frac{1}{2} I \omega^2 + \frac{1}{2} mv^2 \quad \text{...(1)}
\]

But \( I = \frac{1}{2} mr^2 \) and \( v = r \omega \) \quad \text{...(2)}

Substituting eqn. (2) in eqn. (1),

\[
E = \frac{1}{2} \times \frac{1}{2} (mr^2)(\omega^2) + \frac{1}{2} m (r \omega)^2 = \frac{1}{4} mr^2 \omega^2 + \frac{1}{2} mr^2 \omega^2
\]

\[= \frac{3}{4} mr^2 \omega^2\]

A thin metal ring of diameter 0.6m and mass 1kg starts from rest and rolls down on an inclined plane. Its linear velocity on reaching the foot of the plane is 5 m s\(^{-1}\), calculate (i) the moment of inertia of the ring and (ii) the kinetic energy of rotation at that instant.

**Data** : \( R = 0.3 \) m ; \( M = 1 \) kg ; \( v = 5 \) m s\(^{-1}\) ; \( I = ? \) K.E. = ?

**Solution** : \( I = MR^2 = 1 \times (0.3)^2 = 0.09 \) kg m\(^2\)

\[
K.E. = \frac{1}{2} I \omega^2
\]

\[
v = r \omega \quad \therefore \quad \omega = \frac{v}{r} \quad \text{K.E.} = \frac{1}{2} \times 0.09 \times \left( \frac{5}{0.3} \right)^2 = 12.5 \text{ J}\]
3.5 A solid cylinder of mass 200 kg rotates about its axis with angular speed 100 s\(^{-1}\). The radius of the cylinder is 0.25 m. What is the kinetic energy associated with the rotation of the cylinder? What is the magnitude of the angular momentum of the cylinder about its axis?

**Data:** \( M = 200 \text{ kg} \); \( \omega = 100 \text{ s}^{-1} \); \( R = 0.25 \text{ metre} \); \( E_R = ? \); \( L = ? \)

**Solution:**
\[
I = \frac{MR^2}{2} = \frac{200 \times (0.25)^2}{2} = 6.25 \text{ kg m}^2
\]
\[
K.E. = \frac{1}{2} I \omega^2
\]
\[
= \frac{1}{2} \times 6.25 \times (100)^2
\]
\[
E_R = 3.125 \times 10^4 \text{ J}
\]
\[
L = I \omega = 6.25 \times 100 = 625 \text{ kg m}^2 \text{ s}^{-1}
\]

3.6 Calculate the radius of gyration of a rod of mass 100 g and length 100 cm about an axis passing through its centre of gravity and perpendicular to its length.

**Data:** \( M = 100 \text{ g} = 0.1 \text{ kg} \); \( l = 100 \text{ cm} = 1 \text{ m} \)

**Solution:**
The moment of inertia of the rod about an axis passing through its centre of gravity and perpendicular to the length = \( I = MK^2 \)
\[
MK^2 = \frac{ML^2}{12} \quad (or) \quad K^2 = \frac{l^2}{12} \quad (or) \quad K = \frac{L}{\sqrt{12}} = \frac{1}{\sqrt{12}} = 0.2886 \text{ m}.
\]

3.7 A circular disc of mass 100 g and radius 10 cm is making 2 revolutions per second about an axis passing through its centre and perpendicular to its plane. Calculate its kinetic energy.

**Data:** \( M = 100 \text{ g} = 0.1 \text{ kg} \); \( R = 10 \text{ cm} = 0.1 \text{ m} \); \( n = 2 \)

**Solution:** \( \omega = \text{angular velocity} = 2\pi n = 2\pi \times 2 = 4\pi \text{ rad} / \text{s} \)

Kinetic energy of rotation = \( \frac{1}{2} I \omega^2 \)
\[
= \frac{1}{2} \times \frac{1}{2} \times MR^2 \omega^2 = \frac{1}{2} \times \frac{1}{2} \times (0.1) \times (0.1)^2 \times (4\pi)^2
\]
\[
= 3.947 \times 10^{-2} \text{ J}
\]
3.8 Starting from rest, the flywheel of a motor attains an angular velocity 100 rad/s from rest in 10 s. Calculate (i) angular acceleration and (ii) angular displacement in 10 seconds.

**Data :** \( \omega_0 = 0 \); \( \omega = 100 \text{ rad s}^{-1} \); \( t = 10 \text{ s} \); \( \alpha = ? \)

**Solution :** From equations of rotational dynamics,

\[
\omega = \omega_0 + at
\]

(or) \( \alpha = \frac{\omega - \omega_0}{t} = \frac{100 - 0}{10} = 10 \text{ rad s}^{-2} \)

Angular displacement \( \theta = \omega_0 t + \frac{1}{2} \alpha t^2 \)

\[
= 0 + \frac{1}{2} \times 10 \times 10^2 = 500 \text{ rad}
\]

3.9 A disc of radius 5 cm has moment of inertia of 0.02 kg m\(^2\). A force of 20 N is applied tangentially to the surface of the disc. Find the angular acceleration produced.

**Data :** \( I = 0.02 \text{ kg m}^2 \); \( r = 5 \text{ cm} = 5 \times 10^{-2} \text{ m} \); \( F = 20 \text{ N} \); \( \tau = ? \)

**Solution :** Torque \( \tau = F \times 2r = 20 \times 2 \times 5 \times 10^{-2} = 2 \text{ N m} \)

angular acceleration \( \alpha = \frac{r}{I} = \frac{2}{0.02} = 100 \text{ rad} / \text{s}^2 \)

3.10 From the figure, find the moment of the force 45 N about A?

**Data :** Force \( F = 45 \text{ N} \); Moment of the force about A = ?

**Solution :** Moment of the force about A

\[
= \text{Force} \times \text{perpendicular distance} = F \times AO
\]

\[
= 45 \times 6 \sin 30 = 135 \text{ N m}
\]
Self evaluation

(The questions and problems given in this self evaluation are only samples. In the same way any question and problem could be framed from the text matter. Students must be prepared to answer any question and problem from the text matter, not only from the self evaluation.)

3.1 The angular speed of minute arm in a watch is:
(a) $\pi/21600$ rad s$^{-1}$ (b) $\pi/12$ rad s$^{-1}$
(c) $\pi/3600$ rad s$^{-1}$ (d) $\pi/1800$ rad s$^{-1}$

3.2 The moment of inertia of a body comes into play
(a) in linear motion (b) in rotational motion
(c) in projectile motion (d) in periodic motion

3.3 Rotational analogue of mass in linear motion is
(a) Weight (b) Moment of inertia
(c) Torque (d) Angular momentum

3.4 The moment of inertia of a body does not depend on
(a) the angular velocity of the body
(b) the mass of the body
(c) the axis of rotation of the body
(d) the distribution of mass in the body

3.5 A ring of radius $r$ and mass $m$ rotates about an axis passing through its centre and perpendicular to its plane with angular velocity $\omega$. Its kinetic energy is
(a) $mr\omega^2$ (b) $\frac{1}{2}mr\omega^2$ (c) $Io^2$ (d) $\frac{1}{2}Io^2$

3.6 The moment of inertia of a disc having mass $M$ and radius $R$, about an axis passing through its centre and perpendicular to its plane is
(a) $\frac{1}{2}MR^2$ (b) $MR^2$ (c) $\frac{1}{4}MR^2$ (d) $\frac{5}{4}MR^2$

3.7 Angular momentum is the vector product of
(a) linear momentum and radius vector
(b) moment of inertia and angular velocity
(c) linear momentum and angular velocity
(d) linear velocity and radius vector
3.8 The rate of change of angular momentum is equal to
(a) Force (b) Angular acceleration
(c) Torque (d) Moment of Inertia

3.9 Angular momentum of the body is conserved
(a) always (b) never
(c) in the absence of external torque (d) in the presence of external torque

3.10 A man is sitting on a rotating stool with his arms outstretched. Suddenly he folds his arm. The angular velocity
(a) decreases (b) increases
(c) becomes zero (d) remains constant

3.11 An athlete diving off a high springboard can perform a variety of exercises in the air before entering the water below. Which one of the following parameters will remain constant during the fall. The athlete’s
(a) linear momentum (b) moment of inertia
(c) kinetic energy (d) angular momentum

3.12 Obtain an expression for position of centre of mass of two particle system.

3.13 Explain the motion of centre of mass of a system with an example.

3.14 What are the different types of equilibrium?

3.15 Derive the equations of rotational motion.

3.16 Compare linear motion with rotational motion.

3.17 Explain the physical significance of moment of inertia.

3.18 Show that the moment of inertia of a rigid body is twice the kinetic energy of rotation.

3.19 State and prove parallel axes theorem and perpendicular axes theorem.

3.20 Obtain the expressions for moment of inertia of a ring (i) about an axis passing through its centre and perpendicular to its plane, (ii) about its diameter and (iii) about a tangent.
3.21 Obtain the expressions for the moment of inertia of a circular disc (i) about an axis passing through its centre and perpendicular to its plane, (ii) about a diameter, (iii) about a tangent in its plane and (iv) about a tangent perpendicular to its plane.

3.22 Obtain an expression for the angular momentum of a rotating rigid body.

3.23 State the law of conservation of angular momentum.

3.24 A cat is able to land on its feet after a fall. Which principle of physics is being used? Explain.

Problems

3.25 A person weighing 45 kg sits on one end of a seesaw while a boy of 15 kg sits on the other end. If they are separated by 4 m, how far from the boy is the centre of mass situated. Neglect weight of the seesaw.

3.26 Three bodies of masses 2 kg, 4 kg and 6 kg are located at the vertices of an equilateral triangle of side 0.5 m. Find the centre of mass of this collection, giving its coordinates in terms of a system with its origin at the 2 kg body and with the 4 kg body located along the positive X axis.

3.27 Four bodies of masses 1 kg, 2 kg, 3 kg and 4 kg are at the vertices of a rectangle of sides a and b. If $a = 1 \text{ m}$ and $b = 2 \text{ m}$, find the location of the centre of mass. (Assume that, 1 kg mass is at the origin of the system, 2 kg body is situated along the positive x axis and 4 kg along the y axis.)

3.28 Assuming a dumbbell shape for the carbon monoxide (CO) molecule, find the distance of the centre of mass of the molecule from the carbon atom in terms of the distance d between the carbon and the oxygen atom. The atomic mass of carbon is 12 amu and for oxygen is 16 amu. (1 amu $= 1.67 \times 10^{-27} \text{ kg}$)

3.29 A solid sphere of mass 50 g and diameter 2 cm rolls without sliding with a uniform velocity of 5 m s$^{-1}$ along a straight line on a smooth horizontal table. Calculate its total kinetic energy.

( Note : Total $E_K = \frac{1}{2} m v^2 + \frac{1}{2} I_0 \omega^2$ ).
3.30 Compute the rotational kinetic energy of a 2 kg wheel rotating at 6 revolutions per second if the radius of gyration of the wheel is 0.22 m.

3.31 The cover of a jar has a diameter of 8 cm. Two equal, but oppositely directed, forces of 20 N act parallel to the rim of the lid to turn it. What is the magnitude of the applied torque?

Answers

3.1 (d) 3.2 (b) 3.3 (b) 3.4 (a)
3.5 (d) 3.6 (a) 3.7 (b) 3.8 (c)
3.9 (c) 3.10 (b) 3.11 (c)
3.25 3 m from the boy 3.26 0.2916 m, 0.2165 m
3.27 0.5 m, 1.4 m 3.28 \( \frac{16}{28} \)
3.29 0.875 J 3.30 68.71 J
3.31 1.6 N m
4. Gravitation and Space Science

We have briefly discussed the kinematics of a freely falling body under the gravity of the Earth in earlier units. The fundamental forces of nature are gravitational, electromagnetic and nuclear forces. The gravitational force is the weakest among them. But this force plays an important role in the birth of a star, controlling the orbits of planets and evolution of the whole universe.

Before the seventeenth century, scientists believed that objects fell on the Earth due to their inherent property of matter. Galileo made a systematic study of freely falling bodies.

4.1 Newton’s law of gravitation

The motion of the planets, the moon and the Sun was the interesting subject among the students of Trinity college at Cambridge in England. Isaac Newton was also one among these students. In 1665, the college was closed for an indefinite period due to plague. Newton, who was then 23 years old, went home to Lincolnshire. He continued to think about the motion of planets and the moon. One day Newton sat under an apple tree and had tea with his friends. He saw an apple falling to ground. This incident made him to think about falling bodies. He concluded that the same force of gravitation which attracts the apple to the Earth might also be responsible for attracting the moon and keeping it in its orbit. The centripetal acceleration of the moon in its orbit and the downward acceleration of a body falling on the Earth might have the same origin. Newton calculated the centripetal acceleration by assuming moon’s orbit (Fig. 4.1) to be circular.

Acceleration due to gravity on the Earth’s surface, \( g = 9.8 \text{ m s}^{-2} \)

Centripetal acceleration on the moon, \( a_c = \frac{v^2}{r} \)
where \( r \) is the radius of the orbit of the moon \((3.84 \times 10^8 \text{ m})\) and \( v \) is the speed of the moon.

Time period of revolution of the moon around the Earth, \( T = 27.3 \text{ days} \).

The speed of the moon in its orbit, \( v = \frac{2\pi r}{T} \)

\[
v = \frac{2\pi \times 3.84 \times 10^8}{27.3 \times 24 \times 60 \times 60} = 1.02 \times 10^3 \text{ m s}^{-1}
\]

\[
\therefore \text{Centripetal acceleration, } a_c = \frac{v^2}{r} = \frac{0.02 \times 10^3 \times (3.84 \times 10^8)}{2.7 \times 10^3 \times 60^2} = 1.02 \times 10^3 \text{ m s}^{-2}
\]

Newton assumed that both the moon and the apple are accelerated towards the centre of the Earth. But their motions differ, because, the moon has a tangential velocity whereas the apple does not have.

Newton found that \( a_c \) was less than \( g \) and hence concluded that force produced due to gravitational attraction of the Earth decreases with increase in distance from the centre of the Earth. He assumed that this acceleration and therefore force was inversely proportional to the square of the distance from the centre of the Earth. He had found that the value of \( a_c \) was about \( 1/3600 \) of the value of \( g \), since the radius of the lunar orbit \( r \) is nearly 60 times the radius of the Earth \( R \).

The value of \( a_c \) was calculated as follows:

\[
a_c = \frac{g}{3600} = \frac{9.8}{3600} = 2.7 \times 10^{-3} \text{ m s}^{-2}
\]

Newton suggested that gravitational force might vary inversely as the square of the distance between the bodies. He realised that this force of attraction was a case of universal attraction between any two bodies present anywhere in the universe and proposed universal gravitational law.

The law states that every particle of matter in the universe attracts every other particle with a force which is directly proportional to the product
of their masses and inversely proportional to the square of the distance between them.

Consider two bodies of masses $m_1$ and $m_2$ with their centres separated by a distance $r$. The gravitational force between them is

\[ F \alpha \frac{m_1 m_2}{r^2} \]

\[ F = G \frac{m_1 m_2}{r^2} \]

where $G$ is the universal gravitational constant.

If $m_1 = m_2 = 1$ kg and $r = 1$ m, then $F = G$.

Hence, the Gravitational constant ‘G’ is numerically equal to the gravitational force of attraction between two bodies of mass 1 kg each separated by a distance of 1 m. The value of $G$ is $6.67 \times 10^{-11}$ N m$^2$ kg$^{-2}$ and its dimensional formula is $M^{-1} L^3 T^{-2}$.

4.1.1 Special features of the law

(i) The gravitational force between two bodies is an action and reaction pair.

(ii) The gravitational force is very small in the case of lighter bodies. It is appreciable in the case of massive bodies. The gravitational force between the Sun and the Earth is of the order of $10^{27}$ N.

4.2 Acceleration due to gravity

Galileo was the first to make a systematic study of the motion of a body under the gravity of the Earth. He dropped various objects from the leaning tower of Pisa and made analysis of their motion under gravity. He came to the conclusion that “in the absence of air, all bodies will fall at the same rate”. It is the air resistance that slows down a piece of paper or a parachute falling under gravity. If a heavy stone and a parachute are dropped where there is no air, both will fall together at the same rate.

Experiments showed that the velocity of a freely falling body under
Gravity increases at a constant rate. (i.e) with a constant acceleration. The acceleration produced in a body on account of the force of gravity is called *acceleration due to gravity*. It is denoted by $g$. At a given place, the value of $g$ is the same for all bodies irrespective of their masses. It differs from place to place on the surface of the Earth. It also varies with altitude and depth.

The value of $g$ at sea-level and at a latitude of $45^\circ$ is taken as the standard (i.e) $g = 9.8 \text{ m s}^{-2}$

### 4.3 Acceleration due to gravity at the surface of the Earth

Consider a body of mass $m$ on the surface of the Earth as shown in the Fig. 4.3. Its distance from the centre of the Earth is $R$ (radius of the Earth).

The gravitational force experienced by the body is $F = \frac{GMm}{R^2}$ where $M$ is the mass of the Earth.

From Newton’s second law of motion,

Force $F = mg$.

Equating the above two forces, $\frac{GMm}{R^2} = mg$.

$\therefore g = \frac{GM}{R^2}$

This equation shows that $g$ is independent of the mass of the body $m$. But, it varies with the distance from the centre of the Earth. If the Earth is assumed to be a sphere of radius $R$, the value of $g$ on the surface of the Earth is given by $g = \frac{GM}{R^2}$

#### 4.3.1 Mass of the Earth

From the expression $g = \frac{GM}{R^2}$, the mass of the Earth can be calculated as follows:

$$M = \frac{gR^2}{G} = \frac{9.8 \times (6.38 \times 10^6)^2}{6.67 \times 10^{-11}} = 5.98 \times 10^{24} \text{ kg}$$
4.4 Variation of acceleration due to gravity

(i) Variation of $g$ with altitude

Let $P$ be a point on the surface of the Earth and $Q$ be a point at an altitude $h$. Let the mass of the Earth be $M$ and radius of the Earth be $R$. Consider the Earth as a spherical shaped body.

The acceleration due to gravity at $P$ on the surface is

$$g = \frac{GM}{R^2} \quad \cdots (1)$$

Let the body be placed at $Q$ at a height $h$ from the surface of the Earth. The acceleration due to gravity at $Q$ is

$$g_h = \frac{GM}{(R+h)^2} \quad \cdots (2)$$

By simplifying and expanding using binomial theorem, $g_h = g\left(1 - \frac{2h}{R}\right)$

The value of acceleration due to gravity decreases with increase in height above the surface of the Earth.

(ii) Variation of $g$ with depth

Consider the Earth to be a homogeneous sphere with uniform density of radius $R$ and mass $M$.

Let $P$ be a point on the surface of the Earth and $Q$ be a point at a depth $d$ from the surface.

The acceleration due to gravity at $P$ on the surface is $g = \frac{GM}{R^2}$.

If $\rho$ be the density, then, the mass of the Earth is $M = \frac{4}{3} \pi R^3 \rho$. 
\[ g = \frac{4}{3} G \pi R \rho \]  

... (1)

The acceleration due to gravity at \( Q \) at a depth \( d \) from the surface of the Earth is

\[ g_d = \frac{G M_d}{(R-d)^2} \]

where \( M_d \) is the mass of the inner sphere of the Earth of radius \((R-d)\).

\[ M_d = \frac{4}{3} \pi (R-d)^3 \rho \]

\[ \therefore g_d = \frac{4}{3} G \pi (R-d) \rho \]  

... (2)

dividing (2) by (1),

\[ \frac{g_d}{g} = \frac{R-d}{R} \]

\[ g_d = g \left( 1 - \frac{d}{R} \right) \]

The value of acceleration due to gravity decreases with increase of depth.

**iii) Variation of \( g \) with latitude** (Non-sphericity of the Earth)

The Earth is not a perfect sphere. It is an ellipsoid as shown in the Fig. 4.6. It is flattened at the poles where the latitude is 90° and bulged at the equator where the latitude is 0°.

The radius of the Earth at equatorial plane \( R_e \) is greater than the radius along the poles \( R_p \) by about 21 km.

We know that \( g = \frac{G M}{R^2} \)

\[ \therefore g \propto \frac{1}{R^2} \]

The value of \( g \) varies inversely as the square of radius of the Earth. The radius at the equator is the greatest. Hence the value of \( g \)
is minimum at the equator. The radius at poles is the least. Hence, the
value of $g$ is maximum at the poles. The value of $g$ increases from
the equator to the poles.

(iv) Variation of $g$ with latitude (Rotation of the Earth)

Let us consider the Earth as a homogeneous sphere of mass $M$
and radius $R$. The Earth rotates about an axis passing through its north
and south poles. The Earth rotates from west to east in 24 hours. Its angular
velocity is $7.3 \times 10^{-5} \text{ rad s}^{-1}$.

Consider a body of mass $m$ on the surface of the Earth at $P$ at a
latitude $\theta$. Let $\omega$ be the angular velocity. The force (weight) $F = mg$
acts along $PO$. It could be resolved into two rectangular components
(i) $mg \cos \theta$ along $PB$ and (ii) $mg \sin \theta$ along $PA$ (Fig. 4.7).

From the $\triangle OPB$, it is found that $BP = R \cos \theta$. The particle describes a
circle with $B$ as centre and radius $BP = R \cos \theta$.

The body at $P$ experiences a centrifugal force (outward force) $F_c$
due to the rotation of the Earth.

\[(\text{i.e}) \quad F_c = mR\omega^2 \cos \theta \text{.}\]

The net force along $PC = mg \cos \theta - mR\omega^2 \cos \theta$
\therefore The body is acted upon by two forces along $PA$ and $PC$.

The resultant of these two forces is

\[F = mg \sqrt{1 - \frac{2R\omega^2 \cos^2 \theta}{g} \cdot \frac{R^2 \omega^4 \cos^2 \theta}{g^2}}\]

since $\frac{R^2 \omega^4}{g^2}$ is very small, the term $\frac{R^2 \omega^4 \cos^2 \theta}{g^2}$ can be neglected.

The force, $F = mg \sqrt{1 - \frac{2R\omega^2 \cos^2 \theta}{g}}$ \quad \ldots \ (1)
If \( g' \) is the acceleration of the body at \( P \) due to this force \( F \), we have, \( F = mg' \) ... (2)

by equating (2) and (1)

\[
mg' = mg \sqrt{1 - \frac{2 R \omega^2 \cos^2 \theta}{g}}
\]

\[
g' = g \left(1 - \frac{R \omega^2 \cos^2 \theta}{g}\right)
\]

Case (i) At the poles, \( \theta = 90^\circ \); \( \cos \theta = 0 \)

\[\therefore g' = g\]

Case (ii) At the equator, \( \theta = 0 \); \( \cos \theta = 1 \)

\[\therefore g' = g \left(1 - \frac{R \omega^2}{g}\right)\]

So, the value of acceleration due to gravity is maximum at the poles.

4.5 Gravitational field

Two masses separated by a distance exert gravitational forces on one another. This is called action at–a–distance. They interact even though they are not in contact. This interaction can also be explained with the field concept. A particle or a body placed at a point modifies a space around it which is called gravitational field. When another particle is brought in this field, it experiences gravitational force of attraction. The gravitational field is defined as the space around a mass in which it can exert gravitational force on other mass.

4.5.1 Gravitational field intensity

Gravitational field intensity or strength at a point is defined as the force experienced by a unit mass placed at that point. It is denoted by \( \mathbf{E} \). It is a vector quantity. Its unit is N kg\(^{-1}\).

Consider a body of mass \( M \) placed at a point \( Q \) and another body of mass \( m \) placed at \( P \) at a distance \( r \) from \( Q \).
The mass \( M \) develops a field \( E \) at \( P \) and this field exerts a force \( F = mE \).

The gravitational force of attraction between the masses \( m \) and \( M \) is

\[
F = \frac{GMm}{r^2}
\]

The gravitational field intensity at \( P \) is

\[
E = \frac{F}{m} = \frac{GM}{r^2}
\]

\( \therefore E = \frac{GM}{r^2} \)

Gravitational field intensity is the measure of gravitational field.

### 4.5.2 Gravitational potential difference

Gravitational potential difference between two points is defined as the amount of work done in moving unit mass from one point to another point against the gravitational force of attraction.

Consider two points \( A \) and \( B \) separated by a distance \( dr \) in the gravitational field.

The work done in moving unit mass from \( A \) to \( B \) is

\[
dv = W_{A \rightarrow B}
\]

Gravitational potential difference \( dv = -E \, dr \)

Here negative sign indicates that work is done against the gravitational field.

### 4.5.3 Gravitational potential

Gravitational potential at a point is defined as the amount of work done in moving unit mass from the point to infinity against the gravitational field. It is a scalar quantity. Its unit is \( \text{N m kg}^{-1} \).

### 4.5.4 Expression for gravitational potential at a point

Consider a body of mass \( M \) at the point \( C \). Let \( P \) be a point at a distance \( r \) from \( C \). To calculate the gravitational potential at \( P \) consider two points \( A \) and \( B \). The point \( A \), where the unit mass is placed is at a distance \( x \) from \( C \).

\[
\text{Fig. 4.10 Gravitational potential}
\]
The gravitational field at \( A \) is \( E = \frac{GM}{x^2} \).

The work done in moving the unit mass from \( A \) to \( B \) through a small distance \( dx \) is \( dw = dv = -E.dx \).

Negative sign indicates that work is done against the gravitational field.

\[
dv = -\frac{GM}{x^2} \, dx
\]

The work done in moving the unit mass from the point \( P \) to infinity is

\[
\int_\infty^r dv = -\int r \frac{GM}{x^2} \, dx
\]

The gravitational potential is negative, since the work is done against the field. (i.e) the gravitational force is always attractive.

### 4.5.5 Gravitational potential energy

Consider a body of mass \( m \) placed at \( P \) at a distance \( r \) from the centre of the Earth. Let the mass of the Earth be \( M \).

When the mass \( m \) is at \( A \) at a distance \( x \) from \( Q \), the gravitational force of attraction on it due to mass \( M \) is given by

\[
F = \frac{GMm}{x^2}
\]

The work done in moving the mass \( m \) through a small distance \( dx \) from \( A \) to \( B \) along the line joining the two centres of masses \( m \) and \( M \) is \( dw = -F.dx \).

Negative sign indicates that work is done against the gravitational field.

\[
\therefore \ dw = -\frac{GMm}{x^2} \cdot dx
\]

The gravitational potential energy of a mass \( m \) at a distance \( r \) from another mass \( M \) is defined as the amount of work done in moving the mass \( m \) from a distance \( r \) to infinity.

The total work done in moving the mass \( m \) from a distance \( r \) to
infinity is 

\[ \int_0^\infty \frac{G M m}{x^2} \, dx \]

\[ W = -G M m \int_0^\infty \frac{1}{x^2} \, dx \]

\[ U = -G M m \int_0^r \frac{1}{x^2} \, dx \]

Gravitational potential energy is zero at infinity and decreases as the distance decreases. This is due to the fact that the gravitational force exerted on the body by the Earth is attractive. Hence the gravitational potential energy \( U \) is negative.

4.5.6 Gravitational potential energy near the surface of the Earth

Let the mass of the Earth be \( M \) and its radius be \( R \). Consider a point \( A \) on the surface of the Earth and another point \( B \) at a height \( h \) above the surface of the Earth. The work done in moving the mass \( m \) from \( A \) to \( B \) is 

\[ U = U_B - U_A \]

\[ U = -G M m \left[ \frac{1}{(R + h)} - \frac{1}{R} \right] \]

\[ U = G M m \left[ \frac{1}{R} - \frac{1}{(R + h)} \right] \]

\[ U = \frac{G M m h}{R(R + h)} \]

If the body is near the surface of the Earth, \( h \) is very small when compared with \( R \). Hence \( (R+h) \) could be taken as \( R \).

\[ U = \frac{G M m h}{R^2} \]

\[ U = mgh \quad \therefore \frac{G M}{R^2} = g \]

4.6 Inertial mass

According to Newton’s second law of motion \( (F = ma) \), the mass of a body can be determined by measuring the acceleration produced in it

* Potential energy is represented by \( U \) (Upsilon).
by a constant force. (i.e) \( m = F/a \). **Intertial mass of a body is a measure of the ability of a body to oppose the production of acceleration in it by an external force.**

If a constant force acts on two masses \( m_A \) and \( m_B \) and produces accelerations \( a_A \) and \( a_B \) respectively, then, \( F = m_Aa_A = m_Ba_B \)

\[
\therefore \frac{m_A}{m_B} = \frac{a_B}{a_A}
\]

The ratio of two masses is independent of the constant force. If the same force is applied on two different bodies, the inertial mass of the body is more in which the acceleration produced is less.

If one of the two masses is a standard kilogram, the unknown mass can be determined by comparing their accelerations.

### 4.7 Gravitational mass

According to Newton’s law of gravitation, the gravitational force on a body is proportional to its mass. We can measure the mass of a body by measuring the gravitational force exerted on it by a massive body like Earth. **Gravitational mass is the mass of a body which determines the magnitude of gravitational pull between the body and the Earth.** This is determined with the help of a beam balance.

If \( F_A \) and \( F_B \) are the gravitational forces of attraction on the two bodies of masses \( m_A \) and \( m_B \) due to the Earth, then

\[
F_A = \frac{G m_A M}{R^2} \quad \text{and} \quad F_B = \frac{G m_B M}{R^2}
\]

where \( M \) is mass of the Earth, \( R \) is the radius of the Earth and \( G \) is the gravitational constant.

\[
\therefore \frac{m_A}{m_B} = \frac{F_A}{F_B}
\]

If one of the two masses is a standard kilogram, the unknown mass can be determined by comparing the gravitational forces.

### 4.8 Escape speed

If we throw a body upwards, it reaches a certain height and then falls back. This is due to the gravitational attraction of the Earth. If we throw the body with a greater speed, it rises to a greater height. If the
body is projected with a speed of 11.2 km/s, it escapes from the Earth and never comes back. The escape speed is the minimum speed with which a body must be projected in order that it may escape from the gravitational pull of the planet.

Consider a body of mass $m$ placed on the Earth’s surface. The gravitational potential energy is $E_P = -\frac{GMm}{R}$ where $M$ is the mass of the Earth and $R$ is its radius.

If the body is projected up with a speed $v_e$, the kinetic energy is $E_K = \frac{1}{2}mv_e^2$.

∴ the initial total energy of the body is $E_i = \frac{1}{2}mv_e^2 - \frac{GMm}{R}$ ... (1)

If the body reaches a height $h$ above the Earth’s surface, the gravitational potential energy is $E_P = -\frac{GMm}{R+h} + \frac{GMm}{R}$.

Let the speed of the body at the height is $v$, then its kinetic energy is, $E_K = \frac{1}{2}mv^2$.

Hence, the final total energy of the body at the height is $E_f = \frac{1}{2}mv^2 - \frac{GMm}{R+h}$ ... (2)

We know that the gravitational force is a conservative force and hence the total mechanical energy must be conserved.

∴ $E_i = E_f$

(i.e) $\frac{mv_e^2}{2} - \frac{GMm}{R} = \frac{mv^2}{2} - \frac{GMm}{R+h}$

The body will escape from the Earth’s gravity at a height where the gravitational field ceases out. (i.e) $h = \infty$. At the height $h = \infty$, the speed $v$ of the body is zero.
Thus \[ \frac{mv_e^2}{2} - \frac{GMm}{R} = 0 \]

\[ v_e = \sqrt{\frac{2GM}{R}} \]

From the relation \( g = \frac{GM}{R^2} \), we get \( GM = gR^2 \)

Thus, the escape speed is \( v_e = \sqrt{2gR} \)

The escape speed for Earth is 11.2 km/s, for the planet Mercury it is 4 km/s and for Jupiter it is 60 km/s. The escape speed for the moon is about 2.5 km/s.

4.8.1 An interesting consequence of escape speed with the atmosphere of a planet

We know that the escape speed is independent of the mass of the body. Thus, molecules of a gas and very massive rockets will require the same initial speed to escape from the Earth or any other planet or moon.

The molecules of a gas move with certain average velocity, which depends on the nature and temperature of the gas. At moderate temperatures, the average velocity of oxygen, nitrogen and carbon–di–oxide is in the order of 0.5 km/s to 1 km/s and for lighter gases hydrogen and helium it is in the order of 2 to 3 km/s. It is clear that the lighter gases whose average velocities are in the order of the escape speed, will escape from the moon. The gravitational pull of the moon is too weak to hold these gases. The presence of lighter gases in the atmosphere of the Sun should not surprise us, since the gravitational attraction of the sun is very much stronger and the escape speed is very high about 620 km/s.

4.9 Satellites

A body moving in an orbit around a planet is called satellite. The moon is the natural satellite of the Earth. It moves around the Earth once in 27.3 days in an approximate circular orbit of radius 3.85 \( \times \) 10\(^5\) km. The first artificial satellite Sputnik was launched in 1956. India launched its first satellite Aryabhatta on April 19, 1975.
4.9.1 Orbital velocity

Artificial satellites are made to revolve in an orbit at a height of few hundred kilometres. At this altitude, the friction due to air is negligible. The satellite is carried by a rocket to the desired height and released horizontally with a high velocity, so that it remains moving in a nearly circular orbit.

The horizontal velocity that has to be imparted to a satellite at the determined height so that it makes a circular orbit around the planet is called orbital velocity.

Let us assume that a satellite of mass $m$ moves around the Earth in a circular orbit of radius $r$ with uniform speed $v_o$. Let the satellite be at a height $h$ from the surface of the Earth. Hence, $r = R+h$, where $R$ is the radius of the Earth.

The centripetal force required to keep the satellite in circular orbit is $F = \frac{m v_o^2}{r} = \frac{m v_o^2}{R+h}$

The gravitational force between the Earth and the satellite is $F = \frac{GMm}{r^2} = \frac{GMm}{(R+h)^2}$

For the stable orbital motion, $\frac{m v_o^2}{R+h} = \frac{GMm}{(R+h)^2}$

$v_o = \sqrt{\frac{GM}{R+h}}$

Since the acceleration due to gravity on Earth’s surface is $g = \frac{GM}{R^2}$,

$v_o = \sqrt{\frac{gR^2}{R+h}}$

If the satellite is at a height of few hundred kilometres (say 200 km), $(R+h)$ could be replaced by $R$.

$\therefore$ orbital velocity, $v_o = \sqrt{gR}$

If the horizontal velocity (injection velocity) is not equal to the calculated value, then the orbit of the satellite will not be circular. If the
injection velocity is greater than the calculated value but not greater than the escape speed \( v_e = \sqrt{2} v_o \), the satellite will move along an elliptical orbit. If the injection velocity exceeds the escape speed, the satellite will not revolve around the Earth and will escape into the space. If the injection velocity is less than the calculated value, the satellite will fall back to the Earth.

**4.9.2 Time period of a satellite**

Time taken by the satellite to complete one revolution round the Earth is called time period.

Time period, \( T = \frac{\text{circumference of the orbit}}{\text{orbital velocity}} \)

\[
T = \frac{2\pi r}{v_o} = \frac{2\pi (R+h)}{v_o}
\]

where \( r \) is the radius of the orbit which is equal to \( (R+h) \).

\[
T = \frac{2\pi (R+h)}{\sqrt{\frac{GM}{R+h}}}
\]

\[
\therefore v_o = \sqrt{\frac{GM}{R+h}}
\]

As \( GM = gR^2 \), \( T = 2\pi \sqrt{\frac{(R+h)^3}{gR^2}} \)

If the satellite orbits very close to the Earth, then \( h \ll R \)

\[
\therefore T = 2\pi \sqrt{\frac{R}{g}}
\]

**4.9.3 Energy of an orbiting satellite**

A satellite revolving in a circular orbit round the Earth possesses both potential energy and kinetic energy. If \( h \) is the height of the satellite above the Earth’s surface and \( R \) is the radius of the Earth, then the radius of the orbit of satellite is \( r = R+h \).

If \( m \) is the mass of the satellite, its potential energy is,

\[
E_p = -\frac{GMm}{r} = -\frac{GMm}{(R+h)}
\]

where \( M \) is the mass of the Earth. The satellite moves with an orbital velocity of \( v_o = \sqrt{\frac{GM}{R+h}} \).

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Hence, its kinetic energy is, \[ E_K = \frac{1}{2} m v^2 \]

The total energy of the satellite is, 
\[ E = E_p + E_K = \frac{-GMm}{2(R+h)} \]

The negative value of the total energy indicates that the satellite is bound to the Earth.

4.9.4 Geo–stationary satellites

A geo-stationary satellite is a particular type used in television and telephone communications. *A number of communication satellites which appear to remain in fixed positions at a specified height above the equator are called synchronous satellites or geo-stationary satellites.* Some television programmes or events occurring in other countries are often transmitted ‘live’ with the help of these satellites.

For a satellite to appear fixed at a position above a certain place on the Earth, its orbital period around the Earth must be exactly equal to the rotational period of the Earth about its axis.

Consider a satellite of mass \( m \) moving in a circular orbit around the Earth at a distance \( r \) from the centre of the Earth. For synchronisation, its period of revolution around the Earth must be equal to the period of rotation of the Earth (ie) 1 day = 24 hr = 86400 seconds.

The speed of the satellite in its orbit is
\[ v = \frac{\text{Circumference of orbit}}{\text{Time period}} \]
\[ v = \frac{2\pi r}{T} \]

The centripetal force is \( F = \frac{mv^2}{r} \)
\[ \therefore F = \frac{4mu^2r}{T^2} \]

The gravitational force on the satellite due to the Earth is
\[ F = \frac{GMm}{r^2} \]

For the stable orbital motion \( \frac{4mu^2r}{T^2} = \frac{GMm}{r^2} \) (or) \( r^3 = \frac{GMT^2}{4\pi^2} \)
We know that, \( g = \frac{GM}{R^2} \)

\[ \therefore r^3 = \frac{gR^2T^2}{4\pi^2} \]

The orbital radius of the geo-stationary satellite is, \( r = \left( \frac{gR^2T^2}{4\pi^2} \right)^{1/3} \)

This orbit is called parking orbit of the satellite.

Substituting \( T = 86400 \) s, \( R = 6400 \) km and \( g = 9.8 \) m/s\(^2\), the radius of the orbit of geo-stationary satellite is calculated as 42400 km.

\[ \therefore \text{The height of the geo-stationary satellite above the surface of the Earth is } h = r - R = 36000 \text{ km}. \]

If a satellite is parked at this height, it appears to be stationary. Three satellites spaced at 120° intervals each above Atlantic, Pacific and Indian oceans provide a worldwide communication network.

### 4.9.5 Polar satellites

The polar satellites revolve around the Earth in a north–south orbit passing over the poles as the Earth spins about its north–south axis.

The polar satellites positioned nearly 500 to 800 km above the Earth travels pole to pole in 102 minutes. The polar orbit remains fixed in space as the Earth rotates inside the orbit. As a result, most of the earth’s surface crosses the satellite in a polar orbit. Excellent coverage of the Earth is possible with this polar orbit. The polar satellites are used for mapping and surveying.

### 4.9.6 Uses of satellites

**(i) Satellite communication**

Communication satellites are used to send radio, television and telephone signals over long distances. These satellites are fitted with devices which can receive signals from an Earth – station and transmit them in different directions.

**(ii) Weather monitoring**

Weather satellites are used to photograph clouds from space and measure the amount of heat reradiated from the Earth. With this information scientists can make better forecasts about weather. You
might have seen the aerial picture of our country taken by the satellites, which is shown daily in the news bulletin on the television and in the newspapers.

(iii) Remote sensing

Collecting of information about an object without physical contact with the object is known as remote sensing. Data collected by the remote sensing satellites can be used in agriculture, forestry, drought assessment, estimation of crop yields, detection of potential fishing zones, mapping and surveying.

(iv) Navigation satellites

These satellites help navigators to guide their ships or planes in all kinds of weather.

4.9.7 Indian space programme

India recognised the importance of space science and technology for the socio-economic development of the society soon after the launch of Sputnik by erstwhile USSR in 1957. The Indian space efforts started in 1960 with the establishment of Thumba Equatorial Rocket Launching Station near Thiruvananthapuram for the investigation of ionosphere. The foundation of space research in India was laid by Dr. Vikram Sarabhai, father of the Indian space programme. Initially, the space programme was carried out by the Department of Atomic Energy. A separate Department of Space (DOS) was established in June 1972. Indian Space Research Organisation (ISRO) under DOS executes space programme through its establishments located at different places in India (Mahendragiri in Tamil Nadu, Sriharikota in Andhra Pradesh, Thiruvananthapuram in Kerala, Bangalore in Karnataka, Ahmedabad in Gujarat, etc...). India is the sixth nation in the world to have the capability of designing, constructing and launching a satellite in an Earth orbit. The main events in the history of space research in India are given below:

Indian satellites

1. Aryabhata - The first Indian satellite was launched on April 19, 1975.
2. Bhaskara - 1
3. Rohini

4. APPLE - It is the abbreviation of Ariane Passenger Pay Load Experiment. APPLE was the first Indian communication satellite put in geo-stationary orbit.

5. Bhaskara - 2

6. INSAT - 1A, 1B, 1C, 1D, 2A, 2B, 2C, 2D, 3A, 3B, 3C, 3D, 3E (Indian National Satellite). Indian National Satellite System is a joint venture of Department of Space, Department of Telecommunications, Indian Meteoro-logical Department and All India Radio and Doordarshan.

7. SROSS - A, B, C and D (Stretched Rohini Satellite Series)

8. IRS - 1A, 1B, 1C, 1D, P2, P3, P4, P5, P6 (Indian Remote Sensing Satellite)

   Data from IRS is used for various applications like drought monitoring, flood damage assessment, flood risk zone mapping, urban planning, mineral prospecting, forest survey etc.

9. METSAT (Kalpana - I) - METSAT is the first exclusive meteorological satellite.

10. GSAT-1, GSAT-2 (Geo-stationary Satellites)

**Indian Launch Vehicles (Rockets)**

1. SLV - 3 - This was India’s first experimental Satellite Launch Vehicle. SLV - 3 was a 22 m long, four stage vehicle weighing 17 tonne. All its stages used solid propellant.
2. ASLV - Augmented Satellite Launch Vehicle. It was a five stage solid propellant vehicle, weighing about 40 tonnes and of about 23.8 m long.

3. PSLV - The Polar Satellite Launch Vehicle has four stages using solid and liquid propellant systems alternately. It is 44.4 m tall weighing about 294 tonnes.

4. GSLV - The Geosynchronous Satellite Launch Vehicle is a 49 m tall, three stage vehicle weighing about 414 tonnes capable of placing satellite of 1800 kg.

**India’s first mission to moon**

ISRO has a plan to send an unmanned spacecraft to moon in the year 2008. The spacecraft is named as CHANDRAYAAN-1. This programme will be much useful in expanding scientific knowledge about the moon, upgrading India’s technological capability and providing challenging opportunities for planetary research for the younger generation. This journey to moon will take 5½ days.

CHANDRAYAAN - 1 will probe the moon by orbiting it at the lunar orbit of altitude 100 km. This mission to moon will be carried by PSLV rocket.

**4.9.8 Weightlessness**

Television pictures and photographs show astronauts and objects floating in satellites orbiting the Earth. This apparent weightlessness is sometimes explained wrongly as zero-gravity condition. Then, what should be the reason?

Consider the astronaut standing on the ground. He exerts a force (his weight) on the ground. At the same time, the ground exerts an equal and opposite force of reaction on the astronaut. Due to this force of reaction, he has a feeling of weight.

When the astronaut is in an orbiting satellite, both the satellite and astronaut have the same acceleration towards the centre of the Earth. Hence, the astronaut does not exert any force on the floor of the satellite. So, the floor of the satellite also does not exert any force of reaction on the astronaut. As there is no reaction, the astronaut has a feeling of weightlessness.
4.9.9 **Rockets – principle**

A rocket is a vehicle which propels itself by ejecting a part of its mass. Rockets are used to carry the payloads (satellites). We have heard of the PSLV and GSLV rockets. All of them are based on Newton’s third law of motion.

Consider a hollow cylindrical vessel closed on both ends with a small hole at one end, containing a mixture of combustible fuels (Fig. 4.14). If the fuel is ignited, it is converted into a gas under high pressure. This high pressure pushes the gas through the hole with an enormous force. This force represents the action A. Hence an opposite force, which is the reaction R, will act on the vessel and make it to move forward.

The force \( F_m \) on the escaping mass of gases and hence the rocket is proportional to the product of the mass of the gases discharged per unit time \( \left( \frac{dm}{dt} \right) \) and the velocity with which they are expelled (\( v \))

\[
F_m \propto \frac{dm}{dt} \cdot v
\]

\[
\therefore F = m \frac{dv}{dt}
\]

This force is known as momentum thrust. If the pressure \( P_1 \) of the escaping gases differs from the pressure \( P_0 \) in the region outside the rocket, there is an additional thrust called the velocity thrust \( F_v \) acts. It is given by \( F_v = A (P_1 - P_0) \) where \( A \) is the area of the nozzle through which the gases escape. Hence, the total thrust on the rocket is \( F = F_m + F_v \).

4.9.10 **Types of fuels**

The hot gases which are produced by the combustion of a mixture of substances are called propellants. The mixture contains a fuel which burns and an oxidizer which supplies the oxygen necessary for the burning of the fuel. The propellants may be in the form of a solid or liquid.
4.9.11 Launching a satellite

To place a satellite at a height of 300 km, the launching velocity should at least be about 8.5 km s\(^{-1}\) or 30600 kmph. If this high velocity is given to the rocket at the surface of the Earth, the rocket will be burnt due to air friction. Moreover, such high velocities cannot be developed by single rocket. Hence, multistage rockets are used.

To be placed in an orbit, a satellite must be raised to the desired height and given the correct speed and direction by the launching rocket (Fig. 4.15).

At lift off, the rocket, with a manned or unmanned satellite on top, is held down by clamps on the launching pad. Now the exhaust gases built-up an upward thrust which exceeds the rocket’s weight. The clamps are then removed by remote control and the rocket accelerates upwards.
To penetrate the dense lower part of the atmosphere, initially the rocket rises vertically and then tilted by a guidance system. The first stage rocket, which may burn for about 2 minutes producing a speed of 3 km s\(^{-1}\), lifts the vehicle to a height of about 60 km and then separates and falls back to the Earth.

The vehicle now goes to its orbital height, say 160 km, where it moves horizontally for a moment. Then the second stage of the rocket fires and increases the speed that is necessary for a circular orbit. By firing small rockets with remote control system, the satellite is separated from the second stage and made to revolve in its orbit.

4.10 The Universe

The science which deals with the study of heavenly bodies in respect of their motions, positions and compositions is known as astronomy. The Sun around which the planets revolve is a star. It is one of the hundred billion stars that comprise our galaxy called the Milky Way. A vast collection of stars held together by mutual gravitation is called a galaxy. The billions of such galaxies form the universe. Hence, the Solar system, stars and galaxies are the constituents of the universe.

4.10.1 The Solar system

The part of the universe in which the Sun occupies the central position of the system holding together all the heavenly bodies such as planets, moons, asteroids, comets ... etc., is called Solar system. The gravitational attraction of the Sun primarily governs the motion of the planets and other heavenly bodies around it. Mercury, Venus, Earth, Mars, Jupiter, Saturn, Uranus, Neptune and Pluto are the nine planets that revolve around the Sun. We can see the planet Venus in the early morning in the eastern sky or in the early evening in the western sky. The planet Mercury can also be seen sometimes after the sunset in the West or just before sunrise in the East. From the Earth, the planet Mars was visibly seen on 27\(^{th}\) August 2003. The planet Mars came closer to the Earth after 60,000 years from a distance of 380 × 10\(^6\) km to a nearby distance of 55.7 × 10\(^6\) km. It would appear again in the year 2287.

Some of the well known facts about the solar system have been summarised in the Table 4.1.
### Table 4.1 Physical properties of the objects in the Solar system

<table>
<thead>
<tr>
<th>Objects</th>
<th>Mass in Earth unit</th>
<th>Semi-major axis of orbit (AU)</th>
<th>Period of revolution in years</th>
<th>Rotation period</th>
<th>Mean density (kg m(^{-3}))</th>
<th>Radius in Earth unit</th>
<th>g in Earth unit</th>
<th>Escape speed (km/s)</th>
<th>Atmosphere</th>
<th>Albedo</th>
<th>Number of satellites</th>
</tr>
</thead>
<tbody>
<tr>
<td>Mercury</td>
<td>0.056</td>
<td>0.387</td>
<td>0.241</td>
<td>58.6 days</td>
<td>5,400</td>
<td>0.38</td>
<td>0.367</td>
<td>4</td>
<td>Nil</td>
<td>0.06</td>
<td>0</td>
</tr>
<tr>
<td>Venus</td>
<td>0.815</td>
<td>0.723</td>
<td>0.615</td>
<td>243 days (E → W)</td>
<td>5100</td>
<td>0.96</td>
<td>0.886</td>
<td>10.5</td>
<td>CO(_2)</td>
<td>0.85</td>
<td>0</td>
</tr>
<tr>
<td>Earth</td>
<td>1.000</td>
<td>1.000</td>
<td>1.000</td>
<td>23 hours 56.1 minutes</td>
<td>5520</td>
<td>1.00</td>
<td>1.000</td>
<td>11.2</td>
<td>N(_2)O(_2)</td>
<td>0.40</td>
<td>1</td>
</tr>
<tr>
<td>Mars</td>
<td>0.107</td>
<td>1.524</td>
<td>1.881</td>
<td>24 hours 27.4 minutes</td>
<td>3970</td>
<td>0.53</td>
<td>0.383</td>
<td>5</td>
<td>CO(_2)</td>
<td>0.15</td>
<td>2</td>
</tr>
<tr>
<td>Ceres (Asteroid)</td>
<td>0.0001</td>
<td>2.767</td>
<td>4.603</td>
<td>90 hours</td>
<td>3340</td>
<td>0.055</td>
<td>0.18</td>
<td>–</td>
<td>–</td>
<td>–</td>
<td>–</td>
</tr>
<tr>
<td>Jupiter</td>
<td>317.9</td>
<td>5.203</td>
<td>11.864</td>
<td>9 hours 50.5 minutes</td>
<td>1330</td>
<td>11.23</td>
<td>2.522</td>
<td>60</td>
<td>He, CH(_4), NH(_3)</td>
<td>0.45</td>
<td>38</td>
</tr>
<tr>
<td>Saturn</td>
<td>95.2</td>
<td>9.540</td>
<td>29.46</td>
<td>10 hours 14 minutes</td>
<td>700</td>
<td>9.41</td>
<td>1.074</td>
<td>37</td>
<td>He, CH(_4)</td>
<td>0.61</td>
<td>30 + 3 rings</td>
</tr>
<tr>
<td>Uranus</td>
<td>14.6</td>
<td>19.18</td>
<td>84.01</td>
<td>10 hours 49 minutes (E → W)</td>
<td>1330</td>
<td>3.98</td>
<td>0.922</td>
<td>21</td>
<td>H(_2), He, CH(_4)</td>
<td>0.35</td>
<td>24</td>
</tr>
<tr>
<td>Neptune</td>
<td>17.2</td>
<td>30.07</td>
<td>164.1</td>
<td>15 hours</td>
<td>1660</td>
<td>3.88</td>
<td>1.435</td>
<td>22.5</td>
<td>H(_2), He, CH(_4)</td>
<td>0.35</td>
<td>2</td>
</tr>
<tr>
<td>Pluto</td>
<td>0.002</td>
<td>39.44</td>
<td>247</td>
<td>6.39 days</td>
<td>2030</td>
<td>0.179</td>
<td>0.051</td>
<td>1.1</td>
<td>–</td>
<td>0.14</td>
<td>0</td>
</tr>
<tr>
<td>Moon</td>
<td>0.0123</td>
<td>–</td>
<td>–</td>
<td>27.32 days</td>
<td>3340</td>
<td>0.27</td>
<td>0.170</td>
<td>2.5</td>
<td>Nil</td>
<td>0.07</td>
<td>–</td>
</tr>
</tbody>
</table>
4.10.2 Planetary motion

The ancient astronomers contributed a great deal by identifying the planets in the solar system and carefully plotting the variations in their positions of the sky over the periods of many years. These data eventually led to models and theories of the solar system.

The first major theory, called the Geo-centric theory was developed by a Greek astronomer, Ptolemy. The Earth is considered to be the centre of the universe, around which all the planets, the moons and the stars revolve in various orbits. The great Indian Mathematician and astronomer Aryabhata of the 5th century AD stated that the Earth rotates about its axis. Due to lack of communication between the scientists of the East and those of West, his observations did not reach the philosophers of the West.

Nicolaus Copernicus, a Polish astronomer proposed a new theory called Helio-centric theory. According to this theory, the Sun is at rest and all the planets move around the Sun in circular orbits. A Danish astronomer Tycho Brahe made very accurate observations of the motion of planets and a German astronomer Johannes Kepler analysed Brahe’s observations carefully and proposed the empirical laws of planetary motion.

Kepler’s laws of planetary motion

(i) The law of orbits

Each planet moves in an elliptical orbit with the Sun at one focus.

A is a planet revolving round the Sun. The position P of the planet where it is very close to the Sun is known as perigee and the position Q of the planet where it is farthest from the Sun is known as apogee.

(ii) The law of areas

The line joining the Sun and the planet (i.e radius vector) sweeps out equal areas in equal interval of times.

The orbit of the planet around the Sun is as shown in Fig. 4.17. The areas $A_1$ and $A_2$ are swept by the radius vector in equal times. The planet covers unequal distances $S_1$ and $S_2$ in equal time. This is due to
the variable speed of the planet. When the planet is closest to the Sun, it covers greater distance in a given time. Hence, the speed is maximum at the closest position. When the planet is far away from the Sun, it covers lesser distance in the same time. Hence the speed is minimum at the farthest position.

**Proof for the law of areas**

Consider a planet moving from A to B. The radius vector OA sweeps a small angle \( d\theta \) at the centre in a small interval of time \( dt \).

From the Fig. 4.18, \( AB = rd\theta \). The small area \( dA \) swept by the radius is,

\[
dA = \frac{1}{2} \times r \times rd\theta
\]

Dividing by \( dt \) on both sides

\[
\frac{dA}{dt} = \frac{1}{2} \times r^2 \times \frac{d\theta}{dt}
\]

\[
\frac{dA}{dt} = \frac{1}{2} r^2 \omega
\]

where \( \omega \) is the angular velocity.

The angular momentum is given by \( L = mr^2\omega \)

\[
\therefore r^2\omega = \frac{L}{m}
\]

Hence,

\[
\frac{dA}{dt} = \frac{1}{2} \times \frac{L}{m}
\]

Since the line of action of gravitational force passes through the axis, the external torque is zero. Hence, the angular momentum is conserved.

\[
\therefore \frac{dA}{dt} = \text{constant.}
\]

(i.e) the area swept by the radius vector in unit time is the same.

**(iii) The law of periods**

The square of the period of revolution of a planet around the Sun
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is directly proportional to the cube of the mean distance between the planet and the Sun.

\[ T^2 \propto r^3 \]

\[ \frac{T^2}{r^3} = \text{constant} \]

**Proof for the law of periods**

Let us consider a planet \( P \) of mass \( m \) moving with the velocity \( v \) around the Sun of mass \( M \) in a circular orbit of radius \( r \).

The gravitational force of attraction of the Sun on the planet is,

\[ F = \frac{GMm}{r^2} \]

The centripetal force is, \( F = \frac{mv^2}{r} \)

Equating the two forces

\[ \frac{mv^2}{r} = \frac{GMm}{r^2} \]

\[ v^2 = \frac{GM}{r} \]  

.....(1)

If \( T \) be the period of revolution of the planet around the Sun, then

\[ v = \frac{2\pi r}{T} \]  

.....(2)

Substituting (2) in (1) \( \frac{4\pi^2 r^2}{T^2} = \frac{GM}{r} \)

\[ \frac{r^3}{T^2} \sqrt{\frac{GM}{4\pi^2}} \]

\( GM \) is a constant for any planet

\[ \therefore T^2 \propto r^3 \]

4.10.3 **Distance of a heavenly body in the Solar system**

The distance of a planet can be accurately measured by the radar echo method. In this method, the radio signals are sent towards the planet from a radar. These signals are reflected back from the surface of a planet. The reflected signals or pulses are received and detected on
Earth. The time \( t \) taken by the signal in going to the planet and coming back to Earth is noted. The signal travels with the velocity of the light \( c \). The distance \( s \) of the planet from the Earth is given by \( s = \frac{ct}{2} \)

4.10.4 Size of a planet

It is possible to determine the size of any planet once we know the distance \( S \) of the planet. The image of every heavenly body is a disc when viewed through an optical telescope. The angle \( \theta \) between two extreme points A and B on the disc with respect to a certain point on the Earth is determined with the help of a telescope. The angle \( \theta \) is called the angular diameter of the planet. The linear diameter \( d \) of the planet is then given by

\[
d = \text{distance} \times \text{angular diameter}
\]

\[
d = s \times \theta
\]

4.10.5 Surface temperatures of the planets

The planets do not emit light of their own. They reflect the Sun’s light that falls on them. Only a fraction of the solar radiation is absorbed and it heats up the surface of the planet. Then it radiates energy. We can determine the surface temperature \( T \) of the planet using Stefan’s law of radiation \( E = \sigma T^4 \) where \( \sigma \) is the Stefan’s constant and \( E \) is the radiant energy emitted by unit area in unit time.

In general, the temperature of the planets decreases as we go away from the Sun, since the planets receive less and less solar energy according to inverse square law. Hence, the planets farther away from the Sun will be colder than those closer to it. Day temperature of Mercury is maximum (340°C) since it is a planet closest to the Sun and that of Pluto is minimum (−240°C). However Venus is an exception as it has very thick atmosphere of carbon-di-oxide. This acts as a blanket and keeps its surface hot. Thus the temperature of Venus is comparatively large of the order of 480°C.

4.10.6 Mass of the planets and the Sun

In the universe one heavenly body revolves around another massive heavenly body. (The Earth revolves around the Sun and the moon revolves
around the Earth). The centripetal force required by the lighter body to revolve around the heavier body is provided by the gravitational force of attraction between the two. For an orbit of given radius, the mass of the heavier body determines the speed with which the lighter body must revolve around it. Thus, if the period of revolution of the lighter body is known, the mass of the heavier body can be determined. For example, in the case of Sun – planet system, the mass of the Sun \( M \) can be calculated if the distance of the Sun from the Earth \( r \), the period of revolution of the Earth around the Sun \( T \) and the gravitational constant \( G \) are known using the relation

\[
M = \frac{4\pi^2 r^3}{GT^2}
\]

### 4.10.7 Atmosphere

The ratio of the amount of solar energy reflected by the planet to that incident on it is known as \textit{albedo}. From the knowledge of albedo, we get information about the existence of atmosphere in the planets. The albedo of Venus is 0.85. It reflects 85% of the incident light, the highest among the nine planets. It is supposed to be covered with thick layer of atmosphere. The planets Earth, Jupiter, Saturn, Uranus and Neptune have high albedoes, which indicate that they possess atmosphere. The planet Mercury and the moon reflect only 6% of the sunlight. It indicates that they have no atmosphere, which is also confirmed by recent space probes.

There are two factors which determine whether the planets have atmosphere or not. They are (i) acceleration due to gravity on its surface and (ii) the surface temperature of the planet.

The value of \( g \) for moon is very small (\( \frac{1}{4} \)th of the Earth). Consequently the escape speed for moon is very small. As the average velocity of the atmospheric air molecules at the surface temperature of the moon is greater than the escape speed, the air molecules escape.

Mercury has a larger value of \( g \) than moon. Yet there is no atmosphere on it. It is because, Mercury is very close to the Sun and hence its temperature is high. So the mean velocity of the gas molecules is very high. Hence the molecules overcome the gravitational attraction and escape.
4.10.8 Conditions for life on any planet

The following conditions must hold for plant life and animal life to exist on any planet.

(i) The planet must have a suitable living temperature range.
(ii) The planet must have a sufficient and right kind of atmosphere.
(iii) The planet must have considerable amount of water on its surface.

4.10.9 Other objects in the Solar system

(i) Asteroids

Asteroids are small heavenly bodies which orbit round the Sun between the orbits of Mars and Jupiter. They are the pieces of much larger planet which broke up due to the gravitational effect of Jupiter. About 1600 asteroids are revolving around the Sun. The largest among them has a diameter of about 700 km is called Ceres. It circles the Sun once in every 4½ years.

(ii) Comets

A comet consists of a small mass of rock-like material surrounded by large masses of substances such as water, ammonia and methane. These substances are easily vapourised. Comets move round the Sun in highly elliptical orbits and most of the time they keep far away from the Sun. As the comet approaches the Sun, it is heated by the Sun's radiant energy and vapourises and forms a head of about 10000 km in diameter. The comet also develops a tail pointing away from the Sun. Some comets are seen at a fixed regular intervals of time. Halley's comet is a periodic comet which made its appearance in 1910 and in 1986. It would appear again in 2062.

(iii) Meteors and Meteorites

The comets break into pieces as they approach very close to the Sun. When Earth's orbit cross the orbit of comet, these broken pieces fall on the Earth. Most of the pieces are burnt up by the heat generated due to friction in the Earth’s atmosphere. They are called meteors (shooting stars). We can see these meteors in the sky on a clear moonless night.
Some bigger size meteors may survive the heat produced by friction and may not be completely burnt. These blazing objects which manage to reach the Earth are called meteorites.

The formation of craters on the surface of the moon, Mercury and Mars is due to the fact that they have been bombarded by large number of meteorites.

**4.10.10 Stars**

A star is a huge, more or less spherical mass of glowing gas emitting large amount of radiant energy. Billions of stars form a galaxy. There are three types of stars. They are (i) double and multiple stars (ii) intrinsically variable stars and (iii) Novae and super novae.

In a galaxy, there are only a few single stars like the Sun. Majority of the stars are either double stars (binaries) or multiple stars. The binary stars are pairs of stars moving round their common centre of gravity in stable equilibrium. An intrinsically variable star shows variation in its apparent brightness. Some stars suddenly attain extremely large brightness, that they may be seen even during daytime and then they slowly fade away. Such stars are called novae. Supernovae is a large novae.

The night stars in the sky have been given names such as Sirius (Vyadha), Canopas (Agasti), Spica (Chitra), Arcturus (Swathi), Polaris (Dhruva) ... etc. After the Sun, the star Alpha Centauri is nearest to Earth.

**Sun**

The Sun is extremely hot and self-luminous body. It is made of hydrogenous matter. It is the star nearest to the Earth. Its mass is about $1.989 \times 10^{30}$ kg. Its radius is about $6.95 \times 10^{8}$ m. Its distance from the Earth is $1.496 \times 10^{11}$ m. This is known as astronomical unit (AU). Light of the sun takes 8 minutes 20 seconds to reach the Earth. The gravitational force of attraction on the surface of the Sun is about 28 times that on the surface of the Earth.

Sun rotates about its axis from East to West. The period of revolution is 34 days at the pole and 25 days at the equator. The density of material is one fourth that of the Earth. The inner part of the Sun
is a bright disc of temperature $14 \times 10^6$ K known as photosphere. The outer most layer of the Sun of temperature 6000 K is called chromosphere.

### 4.10.11 Constellations

Most of the stars appear to be grouped together forming interesting patterns in the sky. The configurations or groups of star formed in the patterns of animals and human beings are called constellations. There are 88 constellations into which the whole sky has been divided.

If we look towards the northern sky on a clear moonless night during the months of July and August, a group of seven bright stars resembling a bear, the four stars forming a quadrangle form the body, the remaining three stars make the tail and some other faint stars form the paws and head of the bear. This constellation is called Ursa Major or Saptarishi or Great Bear. The constellation Orion resembles the figure of a hunter and Taurus (Vrishabha) resembles the shape of a bull.

### 4.10.12 Galaxy

A large band of stars, gas and dust particles held together by gravitational forces is called a galaxy. Galaxies are really complex in nature consisting of billions of stars. Some galaxies emit a comparatively small amount of radio radiations compared to the total radiations emitted. They are called normal galaxies. Our galaxy Milky Way is a normal galaxy spiral in shape.

The nearest galaxy to us known as Andromeda galaxy, is also a normal galaxy. It is at a distance of $2 \times 10^6$ light years. (The distance travelled by the light in one year $[9.467 \times 10^{12}$ km] is called light year). Some galaxies are found to emit millions of times more radio waves compared to normal galaxies. They are called radio galaxies.

### 4.10.13 Milky Way galaxy

Milky Way looks like a stream of milk across the sky. Some of the important features are given below.

1. **Shape and size**
   Milky Way is thick at the centre and thin at the edges. The diameter of the disc is $10^5$ light years. The thickness of the Milky Way varies from 5000 light years at the centre to 1000 light years at the
position of the Sun and to 500 light years at the edges. The Sun is at a distance of about 27000 light years from the galactic centre.

(ii) Interstellar matter

The interstellar space in the Milky Way is filled with dust and gases called interstellar matter. It is found that about 90% of the matter is in the form of hydrogen.

(iii) Clusters

Groups of stars held by mutual gravitational force in the galaxy are called star clusters. A star cluster moves as a whole in the galaxy. A group of 100 to 1000 stars is called galactic cluster. A group of about 10000 stars is called globular cluster.

(iv) Rotation

The galaxy is rotating about an axis passing through its centre. All the stars in the Milky Way revolve around the centre and complete one revolution in about 300 million years. The Sun, one of the many stars revolves around the centre with a velocity of 250 km/s and its period of revolution is about 220 million years.

(v) Mass

The mass of the Milky Way is estimated to be $3 \times 10^{41}$ kg.

4.10.14 Origin of the Universe

The following three theories have been proposed to explain the origin of the Universe.

(i) Big Bang theory

According to the big bang theory all matter in the universe was concentrated as a single extremely dense and hot fire ball. An explosion occurred about 20 billion years ago and the matter was broken into pieces, thrown off in all directions in the form of galaxies. Due to
continuous movement more and more galaxies will go beyond the boundary and will be lost. Consequently, the number of galaxies per unit volume will go on decreasing and ultimately we will have an empty universe.

(ii) Pulsating theory

Some astronomers believe that if the total mass of the universe is more than a certain value, the expansion of the galaxies would be stopped by the gravitational pull. Then the universe may again contract. After it has contracted to a certain critical size, an explosion again occurs. The expansion and contraction repeat after every eight billion years. Thus we may have alternate expansion and contraction giving rise to a pulsating universe.

(iii) Steady state theory

According to this theory, new galaxies are continuously created out of empty space to fill up the gap caused by the galaxies which escape from the observable part of the universe. This theory, therefore suggests that the universe has always appeared as it does today and the rate of expansion has been the same in the past and will remain the same in future. So a steady state has been achieved so that the total number of galaxies in the universe remains constant.
Solved Problems

4.1 Calculate the force of attraction between two bodies, each of mass 200 kg and 2 m apart on the surface of the Earth. Will the force of attraction be different, if the same bodies are placed on the moon, keeping the separation same?

Data: \( m_1 = m_2 = 200 \text{ kg} ; r = 2 \text{ m} ; G = 6.67 \times 10^{-11} \text{ N m}^2 \text{ kg}^{-2} ; F = ? \)

Solution: 
\[
F = \frac{G m_1 m_2}{r^2} = \frac{6.67 \times 10^{-11} \times 200 \times 200}{(2)^2} \\
F = 6.67 \times 10^{-7} \text{ N}
\]

The force of attraction on the moon will remain same, since G is the universal constant and the masses do not change.

4.2 The acceleration due to gravity at the moon’s surface is 1.67 m s\(^{-2}\). If the radius of the moon is 1.74 \(\times 10^6\) m, calculate the mass of the moon.

Data: \( g = 1.67 \text{ m s}^{-2} ; R = 1.74 \times 10^6 \text{ m} ; G = 6.67 \times 10^{-11} \text{ N m}^2 \text{ kg}^{-2} ; M = ? \)

Solution: 
\[
M = \frac{gR^2}{G} = \frac{1.67 \times (1.74 \times 10^6)^2}{6.67 \times 10^{-11}} \\
M = 7.58 \times 10^{22} \text{ kg}
\]

4.3 Calculate the height above the Earth’s surface at which the value of acceleration due to gravity reduces to half its value on the Earth’s surface. Assume the Earth to be a sphere of radius 6400 km.

Data: \( h = ? ; g_h = \frac{g}{2} ; R = 6400 \times 10^3 \text{ m} \)

Solution: 
\[
\frac{g_h}{g} = \frac{R^2}{(R+h)^2} = \left( \frac{R}{R+h} \right)^2 \\
\frac{g}{2g} = \left( \frac{R}{R+h} \right)^2 \\
\frac{R}{R+h} = \frac{1}{\sqrt{2}}
\]

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\[ h = (\sqrt{2} - 1) \quad R = (1.414 - 1) \quad 6400 \times 10^3 \]
\[ h = 2649.6 \times 10^3 \text{ m} \]

At a height of 2649.6 km from the Earth’s surface, the acceleration due to gravity will be half of its value at the Earth’s surface.

4.4 Determine the escape speed of a body on the moon. Given: radius of the moon is \( 1.74 \times 10^6 \text{ m} \) and mass of the moon is \( 7.36 \times 10^{22} \text{ kg} \).

**Data:** \( G = 6.67 \times 10^{-11} \text{ N m}^2 \text{ kg}^{-2} \); \( R = 1.74 \times 10^6 \text{ m} \); \( M = 7.36 \times 10^{22} \text{ kg} \); \( v_e = ? \)

**Solution:**
\[ v_e = \sqrt{\frac{2GM}{R}} = \sqrt{\frac{2 \times 6.67 \times 10^{-11} \times 7.36 \times 10^{22}}{1.74 \times 10^6}} \]
\[ v_e = 2.375 \text{ km s}^{-1} \]

4.5 The mass of the Earth is 81 times that of the moon and the distance from the centre of the Earth to that of the moon is about \( 4 \times 10^5 \text{ km} \). Calculate the distance from the centre of the Earth where the resultant gravitational force becomes zero when a spacecraft is launched from the Earth to the moon.

**Solution:**

Let the mass of the spacecraft be \( m \). The gravitational force on the spacecraft at \( S \) due to the Earth is opposite in direction to that of the moon. Suppose the spacecraft \( S \) is at a distance \( x \) from the centre of the Earth and at a distance of \( (4 \times 10^5 - x) \) from the moon.

\[ \frac{GM_Em}{x^2} = \frac{GM_m m}{(4 \times 10^5 - x)^2} \]

\[ \frac{M_E}{M_m} = 81 = \frac{x^2}{(4 \times 10^5 - x)^2} \]

\[ x = 3.6 \times 10^5 \text{ km} \]

The resultant gravitational force is zero at a distance of \( 3.6 \times 10^5 \text{ km} \) from the centre of the Earth. The resultant force on \( S \) due to the Earth acts towards the Earth until \( 3.6 \times 10^5 \text{ km} \) is reached. Then it acts towards the moon.
4.6 A stone of mass 12 kg falls on the Earth’s surface. If the mass of the Earth is about \(6 \times 10^{24}\) kg and acceleration due to gravity is 9.8 m s\(^{-2}\), calculate the acceleration produced on the Earth by the stone.

**Data:** 
- \(m = 12\) kg; \(M = 6 \times 10^{24}\) kg; 
- \(g = a_s = 9.8\) m s\(^{-2}\); \(a_E = ?\)

**Solution:** Let \(F\) be the gravitational force between the stone and the Earth.

The acceleration of the stone \((g)\) \(a_s = F/m\)

The acceleration of the Earth, \(a_E = F/M\)

\[
\frac{a_E}{a_s} = \frac{m}{M} = \frac{12}{6 \times 10^{24}} = 2 \times 10^{-24}
\]

\[
a_E = 2 \times 10^{-24} \times 9.8 = 19.6 \times 10^{-24}\text{ m s}^{-2}
\]

4.7 The maximum height upto which astronaut can jump on the Earth is 0.75 m. With the same effort, to what height can he jump on the moon? The mean density of the moon is \((2/3)\) that of the Earth and the radius of the moon is \((1/4)\) that of the Earth.

**Data:**
- \(\rho_m = \frac{2}{3} \rho_E\); \(R_m = \frac{1}{4} R_E\)
- \(h_E = 0.75\) m; \(h_m = ?\)

**Solution:** The astronaut of mass \(m\) jumps a height \(h_E\) on the Earth and a height \(h_m\) on the moon. If he gives himself the same kinetic energy on the Earth and on the moon, the potential energy gained at \(h_E\) and \(h_m\) will be the same.

\[
\therefore \quad mgh = \text{constant}
\]

\[
mg_m h_m = mg_E h_E
\]

\[
\frac{h_m}{R_E} = \frac{g_E}{g_m}
\]

... (1)

For the Earth, \(g_E = \frac{GM_E}{R_E^2} = \frac{4}{3} \pi G R_E \rho_E\)
For the moon, \( g_m = \frac{GM_m}{R_m^2} = \frac{4}{3} \pi G R_m \rho_m \)

\[ \therefore \frac{g_E}{g_m} = \frac{R_E}{R_m} \frac{\rho_E}{\rho_m} \] ...

Equating (1) and (2)

\[ h_m = \frac{R_E}{R_m} \frac{\rho_E}{\rho_m} \times h_E \]

\[ h_m = \frac{R_E}{R_m} \times \frac{\rho_E}{\rho_m} \times 0.75 \]

\[ h_m = 4.5 \text{ m} \]

4.8 Three point masses, each of mass \( m \), are placed at the vertices of an equilateral triangle of side \( a \). What is the gravitational field and potential due to the three masses at the centroid of the triangle.

**Solution:**

The distance of each mass from the centroid \( O \) is \( OA = OB = OC \).

From the \( \Delta ODC \), \( \cos 30^\circ = \frac{a/2}{OC} \)

\[ \therefore OC = \frac{a/2}{\cos 30^\circ} = \frac{a}{\sqrt{3}} \]

Similarly, \( OB = \frac{a}{\sqrt{3}} \) and \( OA = \frac{a}{\sqrt{3}} \).

(i) The gravitational field \( E = \frac{GM}{r^2} \)

\[ \therefore \text{Field at } O \text{ due to } A \text{ is, } E_A = \frac{3GM}{a^2} \text{ (towards } A) \]

Field at \( O \) due to \( B \) is, \( E_B = \frac{3GM}{a^2} \text{ (towards } B) \)

Field at \( O \) due to \( C \) is, \( E_C = \frac{3GM}{a^2} \text{ (towards } C) \)
The resultant field due to $E_B$ and $E_C$ is

$$E_R = \sqrt{E_B^2 + E_C^2 + 2E_BE_C \cos 120^\circ}$$

$$E_R = \sqrt{E_B^2 - E_B^2 - E_B^2} = E_B \quad \therefore E_B = E_C$$

The resultant field $E_R = \frac{3GM}{a^2}$ acts along OD.

Since $E_A$ along OA and $E_B$ along OD are equal and opposite, the net gravitational field is zero at the centroid.

(ii) The gravitational potential is, $v = -\frac{GM}{r}$

Net potential at ‘O’ is

$$v = -\frac{GM}{a/\sqrt{3}} - \frac{GM}{a/\sqrt{3}} - \frac{GM}{a/\sqrt{3}} = -\sqrt{3}\left(\frac{GM}{a} + \frac{GM}{a} + \frac{GM}{a}\right) = -3\sqrt{3}\frac{GM}{a}$$

4.9 A geo-stationary satellite is orbiting the Earth at a height of 6R above the surface of the Earth. Here $R$ is the radius of the Earth. What is the time period of another satellite at a height of 2.5R from the surface of the Earth?

**Data :** The height of the geo-stationary satellite from the Earth’s surface, $h = 6R$

The height of another satellite from the Earth’s surface, $h = 2.5R$

**Solution :** The time period of a satellite is $T = 2\pi \sqrt{\frac{(R+h)^3}{GM}}$

$$\therefore T \propto (R+h)^3$$

For geo-stationary satellite,

$$T_1 \propto \sqrt{(R + 6R)^3} \quad \text{...(1)}$$

For another satellite,

$$T_2 \propto \sqrt{(R + 2.5R)^3} \quad \text{...(2)}$$

Dividing (2) by (1)

$$\frac{T_2}{T_1} = \frac{\sqrt{(3.5R)^3}}{\sqrt{(7R)^3}} = \frac{1}{\sqrt{2}}$$

$$T_2 = \frac{T_1}{2\sqrt{2}} = \frac{24}{2\sqrt{2}} = 8 \text{ hours } 29 \text{ minutes} \quad \therefore T_1 = 24 \text{ hours}$$
Self evaluation

(The questions and problems given in this self evaluation are only samples. In the same way any question and problem could be framed from the text matter. Students must be prepared to answer any question and problem from the text matter, not only from the self evaluation.)

4.1 If the distance between two masses is doubled, the gravitational attraction between them
   (a) is reduced to half    (b) is reduced to a quarter
   (c) is doubled           (d) becomes four times

4.2 The acceleration due to gravity at a height (1/20)th the radius of the Earth above the Earth’s surface is 9 m s\(^{-2}\). Its value at a point at an equal distance below the surface of the Earth is
   (a) 0                     (b) 9 m s\(^{-2}\)
   (c) 9.8 m s\(^{-2}\)       (d) 9.5 m s\(^{-2}\)

4.3 The weight of a body at Earth’s surface is W. At a depth half way to the centre of the Earth, it will be
   (a) W                     (b) W/2
   (c) W/4                   (d) W/8

4.4 Force due to gravity is least at a latitude of
   (a) 0°                    (b) 45°
   (c) 60°                   (d) 90°

4.5 If the Earth stops rotating, the value of g at the equator will
   (a) increase               (b) decrease
   (c) remain same           (d) become zero

4.6 The escape speed on Earth is 11.2 km s\(^{-1}\). Its value for a planet having double the radius and eight times the mass of the Earth is
   (a) 11.2 km s\(^{-1}\)   (b) 5.6 km s\(^{-1}\)
   (c) 22.4 km s\(^{-1}\)   (d) 44.8 km s\(^{-1}\)

4.7 If r represents the radius of orbit of satellite of mass m moving around a planet of mass M. The velocity of the satellite is given by
   (a) \(v^2 = \frac{GM}{r}\)   (b) \(v = \frac{GM}{r}\)
   (c) \(v^2 = \frac{GMm}{r}\)   (d) \(v = \frac{Gm}{r}\)

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4.8 If the Earth is at one fourth of its present distance from the Sun, the duration of the year will be
(a) one fourth of the present year
(b) half the present year
(c) one - eighth the present year
(d) one - sixth the present year

4.9 Which of the following objects do not belong to the solar system?
(a) Comets    (b) Nebulae
(c) Asteroids (d) Planets

4.10 According to Kepler’s law, the radius vector sweeps out equal areas in equal intervals of time. The law is a consequence of the conservation of
(a) angular momentum (b) linear momentum
(c) energy          (d) all the above

4.11 Why is the gravitational force of attraction between the two bodies of ordinary masses not noticeable in everyday life?

4.12 State the universal law of gravitation.

4.13 Define gravitational constant. Give its value, unit and dimensional formula.

4.14 The acceleration due to gravity varies with (i) altitude and (ii) depth. Prove.

4.15 Discuss the variation of g with latitude due to the rotation of the Earth.

4.16 The acceleration due to gravity is minimum at equator and maximum at poles. Give the reason.

4.17 What are the factors affecting the ‘g’ value?

4.18 Why a man can jump higher on the moon than on the Earth?

4.19 Define gravitational field intensity.

4.20 Define gravitational potential.

4.21 Define gravitational potential energy. Deduce an expression for it for a mass in the gravitational field of the Earth.

4.22 Obtain an expression for the gravitational potential at a point.

4.23 Differentiate between inertial mass and gravitational mass.
4.24 The moon has no atmosphere. Why?

4.25 What is escape speed? Obtain an expression for it.

4.26 What is orbital velocity? Obtain an expression for it.

4.27 What will happen to the orbiting satellite, if its velocity varies?

4.28 What are the called geo-stationary satellites?

4.29 Show that the orbital radius of a geo-stationary satellite is 36000 km.

4.30 Why do the astronauts feel weightlessness inside the orbiting spacecraft?

4.31 Deduce the law of periods from the law of gravitation.

4.32 State and prove the law of areas based on conservation of angular momentum.

4.33 State Heliocentric theory.

4.34 State Geocentric theory.

4.35 What is solar system?

4.36 State Kepler’s laws of planetary motion.

4.37 What is albedo?

4.38 What are asteroids?

4.39 What are constellations?

4.40 Write a note on Milky Way.

Problems

4.41 Two spheres of masses 10 kg and 20 kg are 5 m apart. Calculate the force of attraction between the masses.

4.42 What will be the acceleration due to gravity on the surface of the moon, if its radius is \( \frac{1}{4} \) th the radius of the Earth and its mass is \( \frac{1}{80} \) th the mass of the Earth? (Take \( g \) as 9.8 m \( \text{s}^{-2} \))

4.43 The acceleration due to gravity at the surface of the moon is 1.67 m \( \text{s}^{-2} \). The mass of the Earth is about 81 times more massive than the moon. What is the ratio of the radius of the Earth to that of the moon?

4.44 If the diameter of the Earth becomes two times its present value and its mass remains unchanged, then how would the weight of an object on the surface of the Earth be affected?
4.45 Assuming the Earth to be a sphere of uniform density, how much would a body weigh one fourth down to the centre of the Earth, if it weighed 250 N on the surface?

4.46 What is the value of acceleration due to gravity at an altitude of 500 km? The radius of the Earth is 6400 km.

4.47 What is the acceleration due to gravity at a distance from the centre of the Earth equal to the diameter of the Earth?

4.48 What should be the angular velocity of the Earth, so that bodies lying on equator may appear weightless? How many times this angular velocity is faster than the present angular velocity? (Given : \( g = 9.8 \text{ m s}^{-2}; R = 6400 \text{ km} \))

4.49 Calculate the speed with which a body has to be projected vertically from the Earth’s surface, so that it escapes the Earth’s gravitational influence. (\( R = 6.4 \times 10^3 \text{ km}; g = 9.8 \text{ m s}^{-2} \))

4.50 Jupiter has a mass 318 times that of the Earth and its radius is 11.2 times the radius of the Earth. Calculate the escape speed of a body from Jupiter’s surface. (Given : escape speed on Earth is 11.2 km/s)

4.51 A satellite is revolving in circular orbit at a height of 1000 km from the surface of the Earth. Calculate the orbital velocity and time of revolution. The radius of the Earth is 6400 km and the mass of the Earth is \( 6 \times 10^{24} \text{ kg} \).

4.52 An artificial satellite revolves around the Earth at a distance of 3400 km. Calculate its orbital velocity and period of revolution. Radius of the Earth = 6400 km : \( g = 9.8 \text{ m s}^{-2} \).

4.53 A satellite of 600 kg orbits the Earth at a height of 500 km from its surface. Calculate its (i) kinetic energy (ii) potential energy and (iii) total energy (\( M = 6 \times 10^{24} \text{ kg}; R = 6.4 \times 10^6 \text{ m} \))

4.54 A satellite revolves in an orbit close to the surface of a planet of density 6300 kg m\(^{-3}\). Calculate the time period of the satellite. Take the radius of the planet as 6400 km.

4.55 A spaceship is launched into a circular orbit close to the Earth’s surface. What additional velocity has to be imparted to the spaceship in the orbit to overcome the gravitational pull. (\( R = 6400 \text{ km}, g = 9.8 \text{ m s}^{-2} \)).
Answers

4.1  (b)  4.2  (d)  4.3  (b)
4.4  (a)  4.5  (a)  4.6  (c)
4.7  (a)  4.8  (c)  4.9  (b)
4.10 (a)

4.41  $53.36 \times 10^{11}$ N  4.42  1.96 m s$^{-2}$
4.43  3.71  4.44  W/4
4.45  187.5 N  4.46  8.27 m s$^{-2}$
4.47  2.45 m s$^{-2}$  4.48  $1.25 \times 10^3$ rad s$^{-1}$; 17
4.49  11.2 km s$^{-1}$  4.50  59.67 km s$^{-1}$
4.51  7.35 km s$^{-1}$; 1 hour 45 minutes 19 seconds
4.52  6.4 km s$^{-1}$; 9614 seconds
4.53  $1.74 \times 10^{10}$ J; $-3.48 \times 10^{10}$ J; $-1.74 \times 10^{10}$ J
4.54  4734 seconds  4.55  3.28 km s$^{-1}$
5. Mechanics of Solids and Fluids

Matter is a substance, which has certain mass and occupies some volume. Matter exists in three states namely solid, liquid and gas. A fourth state of matter consisting of ionised matter of bare nuclei is called plasma. However in our forth coming discussions, we restrict ourselves to the first three states of matter. Each state of matter has some distinct properties. For example a solid has both volume and shape. It has elastic properties. A gas has the volume of the closed container in which it is kept. A liquid has a fixed volume at a given temperature, but no shape. These distinct properties are due to two factors: (i) interatomic or intermolecular forces (ii) the agitation or random motion of molecules due to temperature.

In solids, the atoms and molecules are free to vibrate about their mean positions. If this vibration increases sufficiently, molecules will shake apart and start vibrating in random directions. At this stage, the shape of the material is no longer fixed, but takes the shape of its container. This is liquid state. Due to increase in their energy, if the molecules vibrate at even greater rates, they may break away from one another and assume gaseous state. Water is the best example for this changing of states. Ice is the solid form of water. With increase in temperature, ice melts into water due to increase in molecular vibration. If water is heated, a stage is reached where continued molecular vibration results in a separation among the water molecules and therefore steam is produced. Further continued heating causes the molecules to break into atoms.

5.1 Intermolecular or interatomic forces

Consider two isolated hydrogen atoms moving towards each other as shown in Fig. 5.1.

As they approach each other, the following interactions are observed.

Fig. 5.1 Electrical origin of interatomic forces
(i) Attractive force $A$ between the nucleus of one atom and electron of the other. This attractive force tends to decrease the potential energy of the atomic system.

(ii) Repulsive force $R$ between the nucleus of one atom and the nucleus of the other atom and electron of one atom with the electron of the other atom. These repulsive forces always tend to increase the energy of the atomic system.

There is a universal tendency of all systems to acquire a state of minimum potential energy. This stage of minimum potential energy corresponds to maximum stability.

If the net effect of the forces of attraction and repulsion leads to decrease in the energy of the system, the two atoms come closer to each other and form a covalent bond by sharing of electrons. On the other hand, if the repulsive forces are more and there is increase in the energy of the system, the atoms will repel each other and do not form a bond.

The variation of potential energy with interatomic distance between the atoms is shown in Fig. 5.2.

![Graph showing the variation of potential energy with interatomic distance](image-url)
It is evident from the graph that as the atoms come closer i.e. when the interatomic distance between them decreases, a stage is reached when the potential energy of the system decreases. When the two hydrogen atoms are sufficiently closer, sharing of electrons takes place between them and the potential energy is minimum. This results in the formation of covalent bond and the interatomic distance is $r_o$.

In solids the interatomic distance is $r_o$, and in the case of liquids it is greater than $r_o$. For gases, it is much greater than $r_o$.

The forces acting between the atoms due to electrostatic interaction between the charges of the atoms are called interatomic forces. Thus, interatomic forces are electrical in nature. The interatomic forces are active if the distance between the two atoms is of the order of atomic size $\approx 10^{-10} \text{ m}$. In the case of molecules, the range of the force is of the order of $10^{-9} \text{ m}$.

### 5.2 Elasticity

When an external force is applied on a body, which is not free to move, there will be a relative displacement of the particles. Due to the property of elasticity, the particles tend to regain their original position. The external forces may produce change in length, volume and shape of the body. This external force which produces these changes in the body is called deforming force. A body which experiences such a force is called deformed body. When the deforming force is removed, the body regains its original state due to the force developed within the body. This force is called restoring force. The property of a material to regain its original state when the deforming force is removed is called elasticity. The bodies which possess this property are called elastic bodies. Bodies which do not exhibit the property of elasticity are called plastic. The study of mechanical properties helps us to select the material for specific purposes. For example, springs are made of steel because steel is highly elastic.

### Stress and strain

In a deformed body, restoring force is set up within the body which tends to bring the body back to the normal position. The magnitude of these restoring force depends upon the deformation caused. This restoring force per unit area of a deformed body is known as stress.
\[ \text{Stress} = \frac{\text{restoring force}}{\text{area}} \quad N \text{ m}^{-2} \]

Its dimensional formula is \( ML^{-1}T^{-2} \).

Due to the application of deforming force, length, volume or shape of a body changes. Or in other words, the body is said to be strained. Thus, \textit{strain produced in a body is defined as the ratio of change in dimension of a body to the original dimension.}

\[ \therefore \text{Strain} = \frac{\text{change in dimension}}{\text{original dimension}} \]

Strain is the ratio of two similar quantities. Therefore it has no unit.

\textit{Elastic limit}

If an elastic material is stretched or compressed beyond a certain limit, it will not regain its original state and will remain deformed. The limit beyond which permanent deformation occurs is called the elastic limit.

\textit{Hooke's law}

English Physicist Robert Hooke (1635 - 1703) in the year 1676 put forward the relation between the extension produced in a wire and the restoring force developed in it. The law formulated on the basis of this study is known as Hooke’s law. According to Hooke’s law, \textit{within the elastic limit, strain produced in a body is directly proportional to the stress that produces it.}

\[ \text{(i.e) stress } \alpha \text{ strain} \]

\[ \frac{\text{Stress}}{\text{Strain}} = \text{a constant, known as modulus of elasticity.} \]

Its unit is \( N \text{ m}^{-2} \) and its dimensional formula is \( ML^{-1}T^{-2} \).

\textbf{5.2.1 Experimental verification of Hooke’s law}

A spring is suspended from a rigid support as shown in the Fig. 5.3. A weight hanger and a light pointer is attached at its lower end such
that the pointer can slide over a scale graduated in millimeters. The initial reading on the scale is noted. A slotted weight of \( m \) kg is added to the weight hanger and the pointer position is noted. The same procedure is repeated with every additional \( m \) kg weight. It will be observed that the extension of the spring is proportional to the weight. This verifies Hooke’s law.

5.2.2 Study of stress - strain relationship

Let a wire be suspended from a rigid support. At the free end, a weight hanger is provided on which weights could be added to study the behaviour of the wire under different load conditions. The extension of the wire is suitably measured and a stress - strain graph is plotted as in Fig. 5.4.

(i) In the figure the region OP is linear. Within a normal stress, strain is proportional to the applied stress. This is Hooke’s law. Upto P, when the load is removed the wire regains its original length along PO. The point P represents the elastic limit, PO represents the elastic range of the material and OB is the elastic strength.

(ii) Beyond P, the graph is not linear. In the region PQ the material is partly elastic and partly plastic. From Q, if we start decreasing the load, the graph does not come to O via P, but traces a straight line QA. Thus a permanent strain OA is caused in the wire. This is called permanent set.

(iii) Beyond Q addition of even a very small load causes enormous strain. This point Q is called the yield point. The region QR is the plastic range.

(iv) Beyond R, the wire loses its shape and becomes thinner and thinner in diameter and ultimately breaks, say at S. Therefore S is the breaking point. The stress corresponding to S is called breaking stress.

Fig. 5.4 Stress - Strain relationship
5.2.3 Three moduli of elasticity

Depending upon the type of strain in the body there are three different types of modulus of elasticity. They are

(i) Young's modulus

(ii) Bulk modulus

(iii) Rigidity modulus

(i) Young's modulus of elasticity

Consider a wire of length \( l \) and cross-sectional area \( A \) stretched by a force \( F \) acting along its length. Let \( dl \) be the extension produced.

\[
\text{Longitudinal stress} = \frac{\text{Force}}{\text{Area}} = \frac{F}{A}
\]

\[
\text{Longitudinal strain} = \frac{\text{change in length}}{\text{original length}} = \frac{dl}{l}
\]

Young's modulus of the material of the wire is defined as the ratio of longitudinal stress to longitudinal strain. It is denoted by \( q \).

\[
\text{Young's modulus} = \frac{\text{longitudinal stress}}{\text{longitudinal strain}} \quad (\text{i.e.} \quad q = \frac{F/A}{dl/l} \quad \text{or} \quad q = \frac{F}{A \cdot dl})
\]

(ii) Bulk modulus of elasticity

Suppose equal forces act perpendicular to the six faces of a cube of volume \( V \) as shown in Fig. 5.6. Due to the action of these forces, let the decrease in volume be \( dV \).

Now, Bulk stress = \( \frac{\text{Force}}{\text{Area}} = \frac{F}{A} \)

Bulk Strain = \( \frac{\text{change in volume}}{\text{original volume}} = \frac{-dV}{V} \)

(The negative sign indicates that volume decreases.)

\[
\text{Fig. 5.5 Young's modulus of elasticity}
\]

\[
\text{Fig. 5.6 Bulk modulus of elasticity}
\]
Bulk modulus of the material of the object is defined as the ratio of bulk stress to bulk strain.

It is denoted by \( k \).

\[
\therefore \text{Bulk modulus} = \frac{\text{Bulk stress}}{\text{Bulk strain}}
\]

(i.e) \( k = \frac{F/A}{\frac{dV}{V}} = \frac{P}{\frac{dV}{V}} \quad \therefore P = \frac{F}{A} \) or \( k = \frac{PV}{dV} \)

(iii) **Rigidity modulus or shear modulus**

Let us apply a force \( F \) tangential to the top surface of a block whose bottom AB is fixed, as shown in Fig. 5.7.

Under the action of this tangential force, the body suffers a slight change in shape, its volume remaining unchanged. The side AD of the block is sheared through an angle \( \theta \) to the position AD'.

If the area of the top surface is \( A \), then shear stress = \( F/A \).

Shear modulus or rigidity modulus of the material of the object is defined as the ratio of shear stress to shear strain. It is denoted by \( n \).

\[
\text{Rigidity modulus} = \frac{\text{shear stress}}{\text{shear strain}}
\]

(i.e) \( n = \frac{F/A}{\theta} = \frac{F}{A\theta} \)

**Table 5.1 Values for the moduli of elasticity**

<table>
<thead>
<tr>
<th>Material</th>
<th>Modulus of elasticity (x 10^{11} Pa)</th>
</tr>
</thead>
<tbody>
<tr>
<td></td>
<td>( q )</td>
</tr>
<tr>
<td>Aluminium</td>
<td>0.70</td>
</tr>
<tr>
<td>Copper</td>
<td>1.1</td>
</tr>
<tr>
<td>Iron</td>
<td>1.9</td>
</tr>
<tr>
<td>Steel</td>
<td>2.0</td>
</tr>
<tr>
<td>Tungsten</td>
<td>3.6</td>
</tr>
</tbody>
</table>

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5.2.4 Relation between the three moduli of elasticity

Suppose three stresses $P$, $Q$ and $R$ act perpendicular to the three faces $ABCD$, $ADHE$ and $ABFE$ of a cube of unit volume (Fig. 5.8). Each of these stresses will produce an extension in its own direction and a compression along the other two perpendicular directions. If $\lambda$ is the extension per unit stress, then the elongation along the direction of $P$ will be $\lambda P$. If $\mu$ is the contraction per unit stress, then the contraction along the direction of $P$ due to the other two stresses will be $\mu Q$ and $\mu R$.

∴ The net change in dimension along the direction of $P$ due to all the stresses is $e = \lambda P - \mu Q - \mu R$.

Similarly the net change in dimension along the direction of $Q$ is $f = \lambda Q - \mu P - \mu R$ and the net change in dimension along the direction of $R$ is $g = \lambda R - \mu P - \mu Q$.

Case (i)

If only $P$ acts and $Q = R = 0$ then it is a case of longitudinal stress.
∴ Linear strain $e = \lambda P$
∴ Young’s modulus $q = \frac{\text{linear stress}}{\text{linear strain}} = \frac{P}{\lambda P}$
(i.e) $q = \frac{1}{\lambda}$ or $\lambda = \frac{1}{q}$ ......(1)

Case (ii)

If $R = 0$ and $P = -Q$, then the change in dimension along $P$ is

(i.e) $e = (\lambda + \mu) P$

Angle of shear $\theta = 2e^* = 2(\lambda + \mu) P$
∴ Rigidity modulus

\[
\frac{p}{\theta} = \frac{P}{2(\lambda + \mu)P} \quad \text{(or)} \quad 2(\lambda + \mu) = \frac{1}{n} \quad \text{......(2)}
\]

* The proof for this is not given here
Case (iii)

If \( P = Q = R \), the increase in volume is \( e + f + g \)

\[ = 3e = 3(\lambda - 2\mu)P \quad \text{(since} \ e = f = g) \]

\[ \therefore \text{Bulk strain} = 3(\lambda - 2\mu)P \]

Bulk modulus \( k = \frac{P}{3(\lambda - 2\mu)P} \quad \text{or} \quad (\lambda - 2\mu) = \frac{1}{3k} \quad \ldots(3) \)

From (2), \( 2(\lambda + \mu) = \frac{1}{n} \)

\[ 2\lambda + 2\mu = \frac{1}{n} \quad \ldots(4) \]

From (3), \( (\lambda - 2\mu) = \frac{1}{3k} \quad \ldots(5) \)

Adding (4) and (5),

\[ 3\lambda = \frac{1}{n} + \frac{1}{3k} \]

\[ \lambda = \frac{1}{3n} + \frac{1}{9k} \]

\[ \therefore \text{From} \ (1), \ \frac{1}{q} = \frac{1}{3n} + \frac{1}{9k} \]

or \[ \frac{9}{q} = \frac{3}{n} + \frac{1}{k} \]

This is the relation between the three moduli of elasticity.

5.2.5 Determination of Young’s modulus by Searle’s method

The Searle’s apparatus consists of two rectangular steel frames A and B as shown in Fig. 5.9. The two frames are hinged together by means of a frame F. A spirit level L is provided such that one of its ends is pivoted to one of the frame B whereas the other end rests on top of a screw working through a nut in the other frame. The bottom
of the screw has a circular scale C which can move along a vertical scale V graduated in mm. This vertical scale and circular scale arrangement act as pitch scale and head scale respectively of a micrometer screw.

The frames A and B are suspended from a fixed support by means of two wires PQ and RS respectively. The wire PQ attached to the frame A is the experimental wire. To keep the reference wire RS taut, a constant weight W is attached to the frame B. To the frame A, a weight hanger is attached in which slotted weights can be added.

To begin with, the experimental wire PQ is brought to the elastic mood by loading and unloading the weights in the hanger in the frame A four or five times, in steps of 0.5 kg. Then with the dead load, the micrometer screw is adjusted to ensure that both the frames are at the same level. This is done with the help of the spirit level. The reading of the micrometer is noted by taking the readings of the pitch scale and head scale. Weights are added to the weight hanger in steps of 0.5 kg upto 4 kg and in each case the micrometer reading is noted by adjusting the spirit level. The readings are again noted during unloading and are tabulated in Table 5.2. The mean extension \( dl \) for M kg of load is found out.

<table>
<thead>
<tr>
<th>Load in weight hanger kg</th>
<th>Micrometer reading</th>
<th>Extension for M kg weight</th>
</tr>
</thead>
<tbody>
<tr>
<td></td>
<td>Loading</td>
<td>Unloading</td>
</tr>
<tr>
<td>W</td>
<td></td>
<td></td>
</tr>
<tr>
<td>W + 0.5</td>
<td></td>
<td></td>
</tr>
<tr>
<td>W + 1.0</td>
<td></td>
<td></td>
</tr>
<tr>
<td>W + 1.5</td>
<td></td>
<td></td>
</tr>
<tr>
<td>W + 2.0</td>
<td></td>
<td></td>
</tr>
<tr>
<td>W + 2.5</td>
<td></td>
<td></td>
</tr>
<tr>
<td>W + 3.0</td>
<td></td>
<td></td>
</tr>
<tr>
<td>W + 3.5</td>
<td></td>
<td></td>
</tr>
<tr>
<td>W + 4.0</td>
<td></td>
<td></td>
</tr>
</tbody>
</table>
If \( l \) is the original length and \( r \) the mean radius of the experimental wire, then Young’s modulus of the material of the wire is given by

\[
q = \frac{F/A}{dl/l} = \frac{F/πr^2}{dl/l}
\]

(i.e.) \( q = \frac{F}{πr^2dl} \)

5.2.6 Applications of modulus of elasticity

Knowledge of the modulus of elasticity of materials helps us to choose the correct material, in right dimensions for the right application. The following examples will throw light on this.

(i) Most of us would have seen a crane used for lifting and moving heavy loads. The crane has a thick metallic rope. The maximum load that can be lifted by the rope must be specified. This maximum load under any circumstances should not exceed the elastic limit of the material of the rope. By knowing this elastic limit and the extension per unit length of the material, the area of cross section of the wire can be evaluated. From this the radius of the wire can be calculated.

(ii) While designing a bridge, one has to keep in mind the following factors (1) traffic load (2) weight of bridge (3) force of winds. The bridge is so designed that it should neither bend too much nor break.

5.3 Fluids

A fluid is a substance that can flow when external force is applied on it. The term fluids include both liquids and gases. Though liquids and gases are termed as fluids, there are marked differences between them. For example, gases are compressible whereas liquids are nearly incompressible. We only use those properties of liquids and gases, which are linked with their ability to flow, while discussing the mechanics of fluids.

5.3.1 Pressure due to a liquid column

Let \( h \) be the height of the liquid column in a cylinder of cross sectional area \( A \). If \( ρ \) is the density of the liquid, then weight of the
liquid column $W$ is given by

$$W = \text{mass of liquid column} \times g = Ah\rho g$$

By definition, pressure is the force acting per unit area.

$$\therefore \text{Pressure} = \frac{\text{weight of liquid column}}{\text{area of cross-section}} = \frac{Ah\rho g}{A} = h\rho g$$

$$\therefore P = h\rho g$$

5.3.2 Pascal’s law

One of the most important facts about fluid pressure is that a change in pressure at one part of the liquid will be transmitted without any change to other parts. This was put forward by Blaise Pascal (1623 - 1662), a French mathematician and physicist. This rule is known as Pascal’s law.

Pascal’s law states that if the effect of gravity can be neglected then the pressure in a fluid in equilibrium is the same everywhere.

Consider any two points A and B inside the fluid. Imagine a cylinder such that points A and B lie at the centre of the circular surfaces at the top and bottom of the cylinder (Fig. 5.11). Let the fluid inside this cylinder be in equilibrium under the action of forces from outside the fluid. These forces act everywhere perpendicular to the surface of the cylinder. The forces acting on the circular, top and bottom surfaces are perpendicular to the forces acting on the cylindrical surface. Therefore the forces acting on the faces at A and B are equal and opposite and hence add to zero. As the areas of these two faces are equal, we can conclude that pressure at A is equal to pressure at B. This is the proof of Pascal’s law when the effect of gravity is not taken into account.

Pascal’s law and effect of gravity

When gravity is taken into account, Pascal’s law is to be modified. Consider a cylindrical liquid column of height $h$ and density $\rho$ in a
vessel as shown in the Fig. 5.12.

If the effect of gravity is neglected, then pressure at \( M \) will be equal to pressure at \( N \). But, if force due to gravity is taken into account, then they are not equal.

As the liquid column is in equilibrium, the forces acting on it are balanced. The vertical forces acting are

(i) Force \( P_1A \) acting vertically down on the top surface.

(ii) Weight \( mg \) of the liquid column acting vertically downwards.

(iii) Force \( P_2A \) at the bottom surface acting vertically upwards.

where \( P_1 \) and \( P_2 \) are the pressures at the top and bottom faces, \( A \) is the area of cross section of the circular face and \( m \) is the mass of the cylindrical liquid column.

\[
\text{At equilibrium, } P_1A + mg - P_2A = 0 \quad \text{or} \quad P_1A + mg = P_2A
\]

\[
P_2 = P_1 + \frac{mg}{A} \quad \text{But} \quad m = Ah\rho
\]

\[
\therefore \quad P_2 = P_1 + \frac{Ah\rho g}{A} \quad \text{(i.e)} \quad P_2 = P_1 + h\rho g
\]

This equation proves that the pressure is the same at all points at the same depth. This results in another statement of Pascal’s law which can be stated as change in pressure at any point in an enclosed fluid at rest is transmitted undiminished to all points in the fluid and act in all directions.

**5.3.3 Applications of Pascal’s law**

**(i) Hydraulic lift**

An important application of Pascal’s law is the hydraulic lift used to lift heavy objects. A schematic diagram of a hydraulic lift is shown in the Fig. 5.13. It consists of a liquid container which has pistons fitted into the small and large opening cylinders. If \( a_1 \) and \( a_2 \) are the areas of the pistons A and B respectively, \( F \) is the force applied on A and \( W \) is the load on B, then
\[ \frac{F}{a_1} = \frac{W}{a_2} \quad \text{or} \quad W = F \frac{a_2}{a_1} \]

This is the load that can be lifted by applying a force \( F \) on \( A \). In the above equation \( \frac{a_2}{a_1} \) is called mechanical advantage of the hydraulic lift. One can see such a lift in many automobile service stations.

**(ii) Hydraulic brake**

When brakes are applied suddenly in a moving vehicle, there is every chance of the vehicle to skid because the wheels are not retarded uniformly. In order to avoid this danger of skidding when the brakes are applied, the brake mechanism must be such that each wheel is equally and simultaneously retarded. A hydraulic brake serves this purpose. It works on the principle of Pascal’s law.

Fig. 5.14 shows the schematic diagram of a hydraulic brake system. The brake system has a main cylinder filled with brake oil. The main cylinder is provided with a piston \( P \) which is connected to the brake pedal.
pedal through a lever assembly. A T shaped tube is provided at the other end of the main cylinder. The wheel cylinder having two pistons P₁ and P₂ is connected to the T tube. The pistons P₁ and P₂ are connected to the brake shoes S₁ and S₂ respectively.

When the brake pedal is pressed, piston P is pushed due to the lever assembly operation. The pressure in the main cylinder is transmitted to P₁ and P₂. The pistons P₁ and P₂ push the brake shoes away, which in turn press against the inner rim of the wheel. Thus the motion of the wheel is arrested. The area of the pistons P₁ and P₂ is greater than that of P. Therefore a small force applied to the brake pedal produces a large thrust on the wheel rim.

The main cylinder is connected to all the wheels of the automobile through pipe line for applying equal pressure to all the wheels.

5.4 Viscosity

Let us pour equal amounts of water and castor oil in two identical funnels. It is observed that water flows out of the funnel very quickly whereas the flow of castor oil is very slow. This is because of the frictional force acting within the liquid. This force offered by the adjacent liquid layers is known as viscous force and the phenomenon is called viscosity.

Viscosity is the property of the fluid by virtue of which it opposes relative motion between its different layers. Both liquids and gases exhibit viscosity but liquids are much more viscous than gases.

**Co-efficient of viscosity**

Consider a liquid to flow steadily through a pipe as shown in the Fig. 5.15. The layers of the liquid which are in contact with the walls of the pipe have zero velocity. As we move towards the axis, the velocity of the liquid layer increases and the centre layer has the maximum velocity v. Consider any two layers P and Q separated by a distance \(dx\). Let \(dv\) be the difference in velocity between the two layers.

![Fig. 5.15 Steady flow of a liquid](image-url)
The viscous force $F$ acting tangentially between the two layers of the liquid is proportional to (i) area $A$ of the layers in contact (ii) velocity gradient $\frac{dv}{dx}$ perpendicular to the flow of liquid.

$$ \therefore F \alpha A \frac{dv}{dx} $$

$$ F = \eta A \frac{dv}{dx} $$

where $\eta$ is the coefficient of viscosity of the liquid.

This is known as Newton’s law of viscous flow in fluids.

If $A = 1m^2$ and $\frac{dv}{dx} = 1s^{-1}$ then $F = \eta$

The coefficient of viscosity of a liquid is numerically equal to the viscous force acting tangentially between two layers of liquid having unit area of contact and unit velocity gradient normal to the direction of flow of liquid.

The unit of $\eta$ is $Ns^{-1}m^{-2}$. Its dimensional formula is $ML^{-1}T^{-1}$.

5.4.1 Streamline flow

The flow of a liquid is said to be steady, streamline or laminar if every particle of the liquid follows exactly the path of its preceding particle and has the same velocity of its preceding particle at every point.

Let $abc$ be the path of flow of a liquid and $v_1$, $v_2$ and $v_3$ be the velocities of the liquid at the points $a$, $b$ and $c$ respectively. During a streamline flow, all the particles arriving at ‘a’ will have the same velocity $v_1$ which is directed along the tangent at the point ‘a’. A particle arriving at $b$ will always have the same velocity $v_2$. This velocity $v_2$ may or may not be equal to $v_1$. Similarly all the particles arriving at the point $c$ will always have the same velocity $v_3$. In other words, in the streamline flow of a liquid, the velocity of every particle crossing a particular point is the same.
The streamline flow is possible only as long as the velocity of the fluid does not exceed a certain value. This limiting value of velocity is called critical velocity.

5.4.2 Turbulent flow

When the velocity of a liquid exceeds the critical velocity, the path and velocities of the liquid become disorderly. At this stage, the flow loses all its orderliness and is called turbulent flow. Some examples of turbulent flow are:

(i) After rising a short distance, the smooth column of smoke from an incense stick breaks up into irregular and random patterns.

(ii) The flash-flood after a heavy rain.

Critical velocity of a liquid can be defined as that velocity of liquid upto which the flow is streamlined and above which its flow becomes turbulent.

5.4.3 Reynolds number

Reynolds number is a pure number which determines the type of flow of a liquid through a pipe. It is denoted by \( N_R \).

It is given by the formula

\[
N_R = \frac{v_c \rho D}{\eta}
\]

where \( v_c \) is the critical velocity, \( \rho \) is the density, \( \eta \) is the co-efficient of viscosity of the liquid and \( D \) is the diameter of the pipe.

If \( N_R \) lies between 0 and 2000, the flow of a liquid is said to be streamline. If the value of \( N_R \) is above 3000, the flow is turbulent. If \( N_R \) lies between 2000 and 3000, the flow is neither streamline nor turbulent, it may switch over from one type to another.

Narrow tubes and highly viscous liquids tend to promote streamline motion while wider tubes and liquids of low viscosity lead to turbulence.

5.4.4 Stoke’s law (for highly viscous liquids)

When a body falls through a highly viscous liquid, it drags the layer of the liquid immediately in contact with it. This results in a relative motion between the different layers of the liquid. As a result of this, the falling body experiences a viscous force \( F \). Stoke performed
many experiments on the motion of small spherical bodies in different fluids and concluded that the viscous force $F$ acting on the spherical body depends on

(i) Coefficient of viscosity $\eta$ of the liquid
(ii) Radius $a$ of the sphere and
(iii) Velocity $v$ of the spherical body.

Dimensionally it can be proved that

$$ F = k \eta av $$

Experimentally Stoke found that

$$ k = 6\pi $$

$$ \therefore F = 6\pi \eta av $$

This is Stoke’s law.

### 5.4.5 Expression for terminal velocity

Consider a metallic sphere of radius ‘$a$’ and density $\rho$ to fall under gravity in a liquid of density $\sigma$. The viscous force $F$ acting on the metallic sphere increases as its velocity increases. A stage is reached when the weight $W$ of the sphere becomes equal to the sum of the upward viscous force $F$ and the upward thrust $U$ due to buoyancy (Fig. 5.17). Now, there is no net force acting on the sphere and it moves down with a constant velocity $v$ called terminal velocity.

$$ \therefore W - F - U = 0 \quad \ldots (1) $$

*Terminal velocity of a body is defined as the constant velocity acquired by a body while falling through a viscous liquid.*

From (1),

$$ W = F + U \quad \ldots (2) $$

According to Stoke’s law, the viscous force $F$ is given by $F = 6\pi \eta av$.

The buoyant force $U = \text{Weight of liquid displaced by the sphere}$

$$ = \frac{4}{3} \pi a^3 \sigma g $$

The weight of the sphere

$$ W = \frac{4}{5} \pi a^3 \rho g $$
Substituting in equation (2).
\[
\frac{4}{3} \pi a^3 \rho g = 6\pi \eta a v + \frac{4}{3} \pi a^3 \sigma g
\]
or
\[
6\pi \eta a v = \frac{4}{3} \pi a^3 (\rho - \sigma) g
\]
\[
\therefore \, v = \frac{2a^2(\rho - \sigma)g}{9\eta}
\]

5.4.6 Experimental determination of viscosity of highly viscous liquids

The coefficient of highly viscous liquid like castor oil can be determined by Stoke’s method. The experimental liquid is taken in a tall, wide jar. Two marking B and C are marked as shown in Fig. 5.18. A steel ball is gently dropped in the jar.

The marking B is made well below the free surface of the liquid so that by the time ball reaches B, it would have acquired terminal velocity \( v \).

When the ball crosses B, a stopwatch is switched on and the time taken \( t \) to reach C is noted. If the distance BC is s, then terminal velocity \( v = \frac{s}{t} \).

The expression for terminal velocity is
\[
v = \frac{2a^2(\rho - \sigma)g}{9\eta}
\]
\[
\therefore \, \frac{s}{t} = \frac{2a^2(\rho - \sigma)g}{9\eta} \quad \text{or} \quad \eta = \frac{2a^2(\rho - \sigma)g}{9 \frac{t}{s}}
\]

Knowing \( a, \rho \) and \( \sigma \), the value of \( \eta \) of the liquid is determined.

Application of Stoke’s law

Falling of rain drops: When the water drops are small in size, their terminal velocities are small. Therefore they remain suspended in air in the form of clouds. But as the drops combine and grow in size, their terminal velocities increases because \( v \propto a^2 \). Hence they start falling as rain.
5.4.7 Poiseuille’s equation

Poiseuille investigated the steady flow of a liquid through a capillary tube. He derived an expression for the volume of the liquid flowing per second through the tube.

Consider a liquid of co-efficient of viscosity $\eta$ flowing, steadily through a horizontal capillary tube of length $l$ and radius $r$. If $P$ is the pressure difference across the ends of the tube, then the volume $V$ of the liquid flowing per second through the tube depends on $\eta$, $r$ and the pressure gradient $\left(\frac{P}{l}\right)$.

(i.e) $V \propto \eta^x r^y \left(\frac{P}{l}\right)^z$

$$V = k \eta^x r^y \left(\frac{P}{l}\right)^z \quad \ldots (1)$$

where $k$ is a constant of proportionality. Rewriting equation (1) in terms of dimensions,

$$[L^3T^{-1}] = [ML^{-1}T^{-1}]^x \cdot [L]^y \cdot \left[\frac{ML^{-1}T^{-2}}{L}\right]^z$$

Equating the powers of $L$, $M$ and $T$ on both sides we get $x = -1$, $y = 4$ and $z = 1$

Substituting in equation (1),

$$V = k \eta^{-1} r^4 \left(\frac{P}{l}\right)$$

$$V = \frac{k \eta^{-1} r^4 P}{l}$$

Experimentally $k$ was found to be equal to $\frac{\pi}{8}$.

$$\therefore \quad V = \frac{\pi \eta r^4}{8l}$$

This is known as Poiseuille’s equation.
5.4.8 Determination of coefficient of viscosity of water by Poiseuille's flow method

A capillary tube of very fine bore is connected by means of a rubber tube to a burette kept vertically. The capillary tube is kept horizontal as shown in Fig. 5.19. The burette is filled with water and the pinch-stopper is removed. The time taken for water level to fall from A to B is noted. If \( V \) is the volume between the two levels A and B, then volume of liquid flowing per second is \( \frac{V}{t} \). If \( l \) and \( r \) are the length and radius of the capillary tube respectively, then

\[
\frac{V}{t} = \frac{\pi Pr^4}{8\eta l} \quad \ldots(1)
\]

If \( \rho \) is the density of the liquid then the initial pressure difference between the ends of the tube is \( P_1 = h_1\rho g \) and the final pressure difference \( P_2 = h_2\rho g \). Therefore the average pressure difference during the flow of water is \( P \) where

\[
P = \frac{P_1 + P_2}{2} = \frac{h_1 + h_2}{2}\rho g = h\rho g \quad \left[ \therefore h = \frac{h_1 + h_2}{2} \right]
\]

Substituting in equation (1), we get

\[
\frac{V}{t} = \frac{\pi h\rho gr^4}{8\eta l} \quad \text{or} \quad \eta = \frac{\pi h\rho gr^4 t}{8lV}
\]

5.4.9 Viscosity - Practical applications

The importance of viscosity can be understood from the following examples.

(i) The knowledge of coefficient of viscosity of organic liquids is used to determine their molecular weights.
(ii) The knowledge of coefficient of viscosity and its variation with temperature helps us to choose a suitable lubricant for specific machines. In light machinery thin oils (example, lubricant oil used in clocks) with low viscosity is used. In heavy machinery, highly viscous oils (example, grease) are used.

5.5 Surface tension

**Intermolecular forces**

The force between two molecules of a substance is called intermolecular force. This intermolecular force is basically electric in nature. When the distance between two molecules is greater, the distribution of charges is such that the mean distance between opposite charges in the molecule is slightly less than the distance between their like charges. So a force of attraction exists. When the intermolecular distance is less, there is overlapping of the electron clouds of the molecules resulting in a strong repulsive force.

The intermolecular forces are of two types. They are (i) cohesive force and (ii) adhesive force.

**Cohesive force**

*Cohesive force is the force of attraction between the molecules of the same substance.* This cohesive force is very strong in solids, weak in liquids and extremely weak in gases.

**Adhesive force**

*Adhesive force is the force of attraction between the molecules of two different substances.* For example due to the adhesive force, ink sticks to paper while writing. Fevicol, gum etc exhibit strong adhesive property.

Water wets glass because the cohesive force between water molecules is less than the adhesive force between water and glass molecules. Whereas, mercury does not wet glass because the cohesive force between mercury molecules is greater than the adhesive force between mercury and glass molecules.

**Molecular range and sphere of influence**

*Molecular range is the maximum distance upto which a molecule can exert force of attraction on another molecule.* It is of the order of $10^{-9}$ m for solids and liquids.
Sphere of influence is a sphere drawn around a particular molecule as centre and molecular range as radius. The central molecule exerts a force of attraction on all the molecules lying within the sphere of influence.

**5.5.1 Surface tension of a liquid**

Surface tension is the property of the free surface of a liquid at rest to behave like a stretched membrane in order to acquire minimum surface area.

Imagine a line AB in the free surface of a liquid at rest (Fig. 5.20). The force of surface tension is measured as the force acting per unit length on either side of this imaginary line AB. The force is perpendicular to the line and tangential to the liquid surface. If F is the force acting on the length l of the line AB, then surface tension is given by

\[ T = \frac{F}{l}. \]

Surface tension is defined as the force per unit length acting perpendicular on an imaginary line drawn on the liquid surface, tending to pull the surface apart along the line. Its unit is \( N\,m^{-1} \) and dimensional formula is \( MT^{-2} \).

**Experiments to demonstrate surface tension**

(i) When a painting brush is dipped into water, its hair gets separated from each other. When the brush is taken out of water, it is observed that its hair will cling together. This is because the free surface of water films tries to contract due to surface tension.

![Fig. 5.20 Force on a liquid surface](image)

![Fig. 5.21 Practical examples for surface tension](image)
(ii) When a sewing needle is gently placed on water surface, it floats. The water surface below the needle gets depressed slightly. The force of surface tension acts tangentially. The vertical component of the force of surface tension balances the weight of the needle.

5.5.2 Molecular theory of surface tension

Consider two molecules P and Q as shown in Fig. 5.22. Taking them as centres and molecular range as radius, a sphere of influence is drawn around them.

The molecule P is attracted in all directions equally by neighbouring molecules. Therefore net force acting on P is zero. The molecule Q is on the free surface of the liquid. It experiences a net downward force because the number of molecules in the lower half of the sphere is more and the upper half is completely outside the surface of the liquid. Therefore all the molecules lying on the surface of a liquid experience only a net downward force.

If a molecule from the interior is to be brought to the surface of the liquid, work must be done against this downward force. This work done on the molecule is stored as potential energy. For equilibrium, a system must possess minimum potential energy. So, the free surface will have minimum potential energy. The free surface of a liquid tends to assume minimum surface area by contracting and remains in a state of tension like a stretched elastic membrane.

5.5.3 Surface energy and surface tension

The potential energy per unit area of the surface film is called surface energy. Consider a metal frame ABCD in which AB is movable. The frame is dipped in a soap solution. A film is formed which pulls AB inwards due to surface tension. If $T$ is the surface tension of the film and $l$ is the length...
of the wire AB, this inward force is given by \(2 \times Tl\). The number 2 indicates the two free surfaces of the film.

If AB is moved through a small distance \(x\) as shown in Fig. 5.23 to the position \(A'B'\), then work done is

\[ W = 2Tlx \]

Work done per unit area = \(\frac{W}{2lx}\)

\[ \therefore \text{Surface energy} = \frac{T \cdot 2lx}{2lx} \]

Surface energy = \(T\)

Surface energy is numerically equal to surface tension.

### 5.5.4 Angle of contact

When the free surface of a liquid comes in contact with a solid, it becomes curved at the point of contact. *The angle between the tangent to the liquid surface at the point of contact of the liquid with the solid and the solid surface inside the liquid is called angle of contact.*

In Fig. 5.24, QR is the tangent drawn at the point of contact Q. The angle PQR is called the angle of contact. When a liquid has concave meniscus, the angle of contact is acute. When it has a convex meniscus, the angle of contact is obtuse.

The angle of contact depends on the nature of liquid and solid in contact. For water and glass, \(\theta\) lies between 8° and 18°. For pure water and clean glass, it is very small and hence it is taken as zero. The angle of contact of mercury with glass is 138°.

### 5.5.5 Pressure difference across a liquid surface

If the free surface of a liquid is plane, then the surface tension acts horizontally (Fig. 5.25a). It has no component perpendicular to the horizontal surface. As a result, there is no pressure difference between the liquid side and the vapour side.

If the surface of the liquid is concave (Fig. 5.25b), then the resultant
force $R$ due to surface tension on a molecule on the surface act vertically upwards. To balance this, an excess of pressure acting downward on the concave side is necessary. On the other hand if the surface is convex (Fig. 5.25c), the resultant $R$ acts downward and there must be an excess of pressure on the concave side acting in the upward direction.

Thus, there is always an excess of pressure on the concave side of a curved liquid surface over the pressure on its convex side due to surface tension.

### 5.5.6 Excess pressure inside a liquid drop

Consider a liquid drop of radius $r$. The molecules on the surface of the drop experience a resultant force acting inwards due to surface tension. Therefore, the pressure inside the drop must be greater than the pressure outside it. The excess of pressure $P$ inside the drop provides a force acting outwards perpendicular to the surface, to balance the resultant force due to surface tension. Imagine the drop to be divided into two equal halves. Considering the equilibrium of the upper hemisphere of the drop, the upward force on the plane face ABCD due to excess pressure $P$ is $P\pi r^2$ (Fig. 5.26).

If $T$ is the surface tension of the liquid, the force due to surface tension acting downward along the circumference of the circle ABCD is $T2\pi r$.

At equilibrium, $P\pi r^2 = T2\pi r$

$$\therefore P = \frac{2T}{r}$$
**Excess pressure inside a soap bubble**

A soap bubble has two liquid surfaces in contact with air, one inside the bubble and the other outside the bubble. Therefore the force due to surface tension is

\[ 2 \times 2\pi rT \]

\[ \therefore \text{At equilibrium, } P\pi r^2 = 2 \times 2\pi rT \]

(i.e) \[ P = \frac{4T}{r} \]

Thus the excess of pressure inside a drop is inversely proportional to its radius i.e. \( P \propto \frac{1}{r} \). As \( P \propto \frac{1}{r} \), the pressure needed to form a very small bubble is high. This explains why one needs to blow hard to start a balloon growing. Once the balloon has grown, less air pressure is needed to make it expand more.

**5.5.7 Capillarity**

The property of surface tension gives rise to an interesting phenomenon called capillarity. When a capillary tube is dipped in water, the water rises up in the tube. The level of water in the tube is above the free surface of water in the beaker (capillary rise). When a capillary tube is dipped in mercury, mercury also rises in the tube. But the level of mercury is depressed below the free surface of mercury in the beaker (capillary fall).

The rise of a liquid in a capillary tube is known as capillarity. The height \( h \) in Fig. 5.27 indicates the capillary rise (for water) or capillary fall (for mercury).

**Illustrations of capillarity**

(i) A blotting paper absorbs ink by capillary action. The pores in the blotting paper act as capillaries.
(ii) The oil in a lamp rises up the wick through the narrow spaces between the threads of the wick.
(iii) A sponge retains water due to capillary action.
(iv) Walls get damped in rainy season due to absorption of water by bricks.
5.5.8 Surface tension by capillary rise method

Let us consider a capillary tube of uniform bore dipped vertically in a beaker containing water. Due to surface tension, water rises to a height \( h \) in the capillary tube as shown in Fig. 5.28. The surface tension \( T \) of the water acts inwards and the reaction of the tube \( R \) outwards. \( R \) is equal to \( T \) in magnitude but opposite in direction. This reaction \( R \) can be resolved into two rectangular components.

(i) Horizontal component \( R \sin \theta \) acting radially outwards

(ii) Vertical component \( R \cos \theta \) acting upwards.

The horizontal component acting all along the circumference of the tube cancel each other whereas the vertical component balances the weight of water column in the tube.

Total upward force = \( R \cos \theta \times \) circumference of the tube

(i.e) \( F = 2\pi r \cos \theta \) \quad or \quad \( F = 2\pi r T \cos \theta \) \quad \( ... (1) \)

\[ \therefore \quad R = T \]

This upward force is responsible for the capillary rise. As the water column is in equilibrium, this force acting upwards is equal to weight of the water column acting downwards.

(i.e) \( F = W \) \quad \( ... (2) \)

Now, volume of water in the tube is assumed to be made up of (i) a cylindrical water column of height \( h \) and (ii) water in the meniscus above the plane CD.

Volume of cylindrical water column = \( \pi r^2 h \)

Volume of water in the meniscus = (Volume of cylinder of height \( r \) and radius \( r \)) – (Volume of hemisphere)
Volume of water in the meniscus = \( (\pi r^2 \times \rho) - \left( \frac{2}{3} \pi r^3 \right) \)

\[ = \frac{1}{3} \pi r^3 \]

Total volume of water in the tube = \( \pi r^2 h + \frac{1}{3} \pi r^3 \)

\[ = \pi r^2 \left( h + \frac{r}{3} \right) \]

If \( \rho \) is the density of water, then weight of water in the tube is

\[ W = \pi r^2 \left( h + \frac{r}{3} \right) \rho g \]

...(3)

Substituting (1) and (3) in (2),

\[ \pi r^2 \left( h + \frac{r}{3} \right) \rho g = 2\pi r T \cos \theta \]

\[ T = \frac{\left( h + \frac{r}{3} \right) \rho g}{2 \cos \theta} \]

Since \( r \) is very small, \( \frac{r}{3} \) can be neglected compared to \( h \).

\[ \therefore T = \frac{h \rho g}{2 \cos \theta} \]

For water, \( \theta \) is small, therefore \( \cos \theta = 1 \)

\[ \therefore T = \frac{h \rho g}{2} \]

### 5.5.9 Experimental determination of surface tension of water by capillary rise method

A clean capillary tube of uniform bore is fixed vertically with its lower end dipping into water taken in a beaker. A needle \( N \) is also fixed with the capillary tube as shown in the Fig. 5.30. The tube is raised or lowered until the tip of the needle just touches the water surface. A travelling microscope \( M \) is focussed on the meniscus of the
water in the capillary tube. The reading $R_1$ corresponding to the lower meniscus is noted. The microscope is lowered and focused on the tip of the needle and the corresponding reading is taken as $R_2$. The difference between $R_1$ and $R_2$ gives the capillary rise $h$.

The radius of the capillary tube is determined using the travelling microscope. If $\rho$ is the density of water then the surface tension of water is given by $T = \frac{\rho g h r}{2}$ where $g$ is the acceleration due to gravity.

### 5.5.10 Factors affecting surface tension

*Impurities present in a liquid appreciably affect surface tension.* A highly soluble substance like salt increases the surface tension whereas sparingly soluble substances like soap decreases the surface tension.

*The surface tension decreases with rise in temperature.* The temperature at which the surface tension of a liquid becomes zero is called critical temperature of the liquid.

### 5.5.11 Applications of surface tension

(i) During stormy weather, oil is poured into the sea around the ship. As the surface tension of oil is less than that of water, it spreads on water surface. Due to the decrease in surface tension, the velocity of the waves decreases. This reduces the wrath of the waves on the ship.

(ii) Lubricating oils spread easily to all parts because of their low surface tension.

(iii) Dirty clothes cannot be washed with water unless some detergent is added to water. When detergent is added to water, one end of the hairpin shaped molecules of the detergent get attracted to water and the other end, to molecules of the dirt. Thus the dirt is suspended surrounded by detergent molecules and this can be easily removed. This detergent action is due to the reduction of surface tension of water when soap or detergent is added to water.

(iv) Cotton dresses are preferred in summer because cotton dresses have fine pores which act as capillaries for the sweat.
5.6 Total energy of a liquid

A liquid in motion possesses pressure energy, kinetic energy and potential energy.

(i) Pressure energy

It is the energy possessed by a liquid by virtue of its pressure.

Consider a liquid of density $\rho$ contained in a wide tank $T$ having a side tube near the bottom of the tank as shown in Fig. 5.31. A frictionless piston of cross sectional area $'a'$ is fitted to the side tube. Pressure exerted by the liquid on the piston is $P = h \rho g$ where $h$ is the height of liquid column above the axis of the side tube. If $x$ is the distance through which the piston is pushed inwards, then

Volume of liquid pushed into the tank = $ax$

$\therefore$ Mass of the liquid pushed into the tank = $ax \rho$

As the tank is wide enough and a very small amount of liquid is pushed inside the tank, the height $h$ and hence the pressure $P$ may be considered as constant.

Work done in pushing the piston through the distance $x = \text{Force on the piston} \times \text{distance moved}$

(i.e) $W = Pax$

This work done is the pressure energy of the liquid of mass $ax\rho$.

$\therefore$ Pressure energy per unit mass of the liquid = $P\frac{ax}{\rho}$

(ii) Kinetic energy

It is the energy possessed by a liquid by virtue of its motion.

If $m$ is the mass of the liquid moving with a velocity $v$, the kinetic energy of the liquid = $\frac{1}{2} mv^2$.

Kinetic energy per unit mass = $\frac{1}{2} \frac{mv^2}{m} = \frac{v^2}{2}$
(iii) **Potential energy**

It is the energy possessed by a liquid by virtue of its height above the ground level.

If \( m \) is the mass of the liquid at a height \( h \) from the ground level, the potential energy of the liquid = \( mgh \)

Potential energy per unit mass = \( \frac{mgh}{m} = gh \)

Total energy of the liquid in motion = pressure energy + kinetic energy + potential energy.

∴ Total energy per unit mass of the flowing liquid = \( \frac{P}{\rho} + \frac{v^2}{2} + gh \)

### 5.6.1 Equation of continuity

Consider a non-viscous liquid in streamline flow through a tube AB of varying cross section as shown in Fig. 5.32 Let \( a_1 \) and \( a_2 \) be the area of cross section, \( v_1 \) and \( v_2 \) be the velocity of flow of the liquid at A and B respectively.

∴ Volume of liquid entering per second at A = \( a_1v_1 \).

If \( \rho \) is the density of the liquid, then mass of liquid entering per second at A = \( a_1v_1\rho \).

Similarly, mass of liquid leaving per second at B = \( a_2v_2\rho \)

If there is no loss of liquid in the tube and the flow is steady, then mass of liquid entering per second at A = mass of liquid leaving per second at B

\( a_1v_1\rho = a_2v_2\rho \) or \( a_1v_1 = a_2v_2 \)

i.e. \( av = \text{constant} \)

This is called as the equation of continuity. From this equation \( v \propto \frac{1}{a} \).

i.e. the larger the area of cross section the smaller will be the velocity of flow of liquid and vice-versa.
5.6.2 Bernoulli’s theorem

In 1738, Daniel Bernoulli proposed a theorem for the streamline flow of a liquid based on the law of conservation of energy. According to Bernoulli’s theorem, for the streamline flow of a non-viscous and incompressible liquid, the sum of the pressure energy, kinetic energy and potential energy per unit mass is a constant.

\[
\frac{P}{\rho} + \frac{v^2}{2} + gh = \text{constant}
\]

This equation is known as Bernoulli’s equation.

Consider streamline flow of a liquid of density \( \rho \) through a pipe AB of varying cross section. Let \( P_1 \) and \( P_2 \) be the pressures and \( a_1 \) and \( a_2 \) the cross sectional areas at A and B respectively. The liquid enters A normally with a velocity \( v_1 \) and leaves B normally with a velocity \( v_2 \). The liquid is accelerated against the force of gravity while flowing from A to B, because the height of B is greater than that of A from the ground level. Therefore \( P_1 \) is greater than \( P_2 \). This is maintained by an external force.

The mass \( m \) of the liquid crossing per second through any section of the tube in accordance with the equation of continuity is

\[
a_1v_1\rho = a_2v_2\rho = m
\]

or \( a_1v_1 = a_2v_2 = \frac{m}{\rho} = V \) .... (1)

As \( a_1 > a_2 \), \( v_1 < v_2 \)

The force acting on the liquid at A = \( P_1a_1 \)
The force acting on the liquid at B = \( P_2a_2 \)

Work done per second on the liquid at A = \( P_1a_1 \times v_1 = P_1V \)

Work done by the liquid at B = \( P_2a_2 \times v_2 = P_2V \)

:. Net work done per second on the liquid by the pressure energy in moving the liquid from A to B is = \( P_1V - P_2V \) ....(2)
If the mass of the liquid flowing in one second from A to B is \( m \), then increase in potential energy per second of liquid from A to B is \( mgh_2 - mgh_1 \).

Increase in kinetic energy per second of the liquid

\[
= \frac{1}{2}mv_2^2 - \frac{1}{2}mv_1^2
\]

According to work-energy principle, work done per second by the pressure energy = Increase in potential energy per second + Increase in kinetic energy per second

\[
(P_1V - P_2V) = (mgh_2 - mgh_1) + \left(\frac{1}{2}mv_2^2 - \frac{1}{2}mv_1^2\right)
\]

\[
P_1V + mgh_1 + \frac{1}{2}mv_1^2 = P_2V + mgh_2 + \frac{1}{2}mv_2^2
\]

\[
\frac{P_1}{\rho} + gh_1 + \frac{1}{2}v_1^2 = \frac{P_2}{\rho} + gh_2 + \frac{1}{2}v_2^2
\]

or

\[
\frac{P}{\rho} + gh + \frac{1}{2}v^2 = \text{constant} \quad \cdots (3)
\]

This is Bernoulli’s equation. Thus the total energy of unit mass of liquid remains constant.

Dividing equation (3) by \( g \),

\[
\frac{P}{\rho g} + \frac{v^2}{2g} + h = \text{constant}
\]

Each term in this equation has the dimension of length and hence is called head. \( \frac{P}{\rho g} \) is called pressure head, \( \frac{v^2}{2g} \) is velocity head and \( h \) is the gravitational head.

**Special case**:

If the liquid flows through a horizontal tube, \( h_1 = h_2 \). Therefore there is no increase in potential energy of the liquid i.e. the gravitational head becomes zero.

\[
\therefore \text{equation (3) becomes}
\]

\[
\frac{P}{\rho} + \frac{1}{2}v^2 = \text{a constant}
\]

This is another form of Bernoulli’s equation.
5.6.3 Application of Bernoulli’s theorem

(i) Lift of an aircraft wing

A section of an aircraft wing and the flow lines are shown in Fig. 5.34. The orientation of the wing relative to the flow direction causes the flow lines to crowd together above the wing. This corresponds to increased velocity in this region and hence the pressure is reduced. But below the wing, the pressure is nearly equal to the atmospheric pressure. As a result of this, the upward force on the underside of the wing is greater than the downward force on the topside. Thus there is a net upward force or lift.

(ii) Blowing of roofs

During a storm, the roofs of huts or tinned roofs are blown off without any damage to other parts of the hut. The blowing wind creates a low pressure $P_1$ on top of the roof. The pressure $P_2$ under the roof is however greater than $P_1$. Due to this pressure difference, the roof is lifted and blown off with the wind.

(iii) Bunsen burner

In a Bunsen burner, the gas comes out of the nozzle with high velocity. Due to this the pressure in the stem of the burner decreases. So, air from the atmosphere rushes into the burner.

(iv) Motion of two parallel boats

When two boats separated by a small distance row parallel to each other along the same direction, the velocity of water between the boats becomes very large compared to that on the outer sides. Because of this, the pressure in between the two boats gets reduced. The high pressure on the outer side pushes the boats inwards. As a result of this, the boats come closer and may even collide.
Solved problems

5.1 A 50 kg mass is suspended from one end of a wire of length 4 m and diameter 3 mm whose other end is fixed. What will be the elongation of the wire? Take \( q = 7 \times 10^{10} \text{ N} \text{ m}^{-2} \) for the material of the wire.

**Data:** \( l = 4 \text{ m}; \ d = 3 \text{ mm} = 3 \times 10^{-3} \text{ m}; \ m = 50 \text{ kg}; \ q = 7 \times 10^{10} \text{ N} \text{ m}^{-2} \)

**Solution:**

\[
q = \frac{F_l}{A dl} \]

\[
\therefore \ dl = \frac{F_l}{\pi^2 q} = \frac{50 \times 9.8 \times 4}{3.14 \times (1.5 \times 10^{-3})^2 \times 7 \times 10^{10}} = 3.96 \times 10^{-3} \text{ m}
\]

5.2 A sphere contracts in volume by 0.01% when taken to the bottom of sea 1 km deep. If the density of sea water is \( 10^3 \text{ kg} \text{ m}^{-3} \), find the bulk modulus of the material of the sphere.

**Data:** \( dV = 0.01\% \) i.e. \( \frac{dV}{V} = \frac{0.01}{100} \); \( h = 1 \text{ km}; \ \rho = 10^3 \text{ kg} \text{ m}^{-3} \)

**Solution:**

\[
dP = 10^3 \times 10^3 \times 9.8 = 9.8 \times 10^6
\]

\[
\therefore \ k = \frac{dP}{dV/V} = \frac{9.8 \times 10^6 \times 100}{0.01} = 9.8 \times 10^{10} \text{ N} \text{ m}^{-2}
\]

5.3 A hydraulic automobile lift is designed to lift cars with a maximum mass of 3000 kg. The area of cross-section of the piston carrying the load is \( 425 \times 10^{-4} \text{ m}^2 \). What maximum pressure would the piston have to bear?

**Data:** \( m = 3000 \text{ kg}, A = 425 \times 10^{-4} \text{ m}^2 \)

**Solution:** Pressure on the piston = \( \frac{\text{Weight of car}}{\text{Area of piston}} = \frac{mg}{A} \)

\[
= \frac{3000 \times 9.8}{425 \times 10^{-4}} = 6.92 \times 10^5 \text{ N} \text{ m}^{-2}
\]

5.4 A square plate of 0.1 m side moves parallel to another plate with a velocity of 0.1 m s\(^{-1}\), both plates being immersed in water. If the viscous force is \( 2 \times 10^{-3} \text{ N} \) and viscosity of water is \( 10^{-3} \text{ N} \text{ s} \text{ m}^{-2} \), find their distance of separation.
Data: Area of plate A = 0.1 × 0.1 = 0.01 m²
Viscous force F = 2 × 10⁻³ N
Velocity dv = 0.1 m s⁻¹
Coefficient of viscosity η = 10⁻³ N s m⁻²

Solution: Distance dx = \( \frac{ηAdv}{F} \)
\[ = \frac{10² × 0.01 × 0.1}{2 × 10⁻³} = 5 × 10⁻⁴ \text{ m} \]

5.5 Determine the velocity for air flowing through a tube of 10⁻² m radius. For air \( ρ = 1.3 \text{ kg m}⁻³ \) and \( η = 187 × 10⁻⁷ \text{ N s m}⁻² \).

Data: \( r = 10⁻² \text{ m} ; ρ = 1.3 \text{ kg m}⁻³ ; η = 187 × 10⁻⁷ \text{ N s m}⁻² ; N_R = 2000 \)

Solution: velocity \( v = \frac{N_Rη}{ρD} \)
\[ = \frac{2000 × 187 × 10⁻⁷}{1.3 × 2 × 10⁻²} = 1.44 \text{ m s}⁻¹ \]

5.6 Fine particles of sand are shaken up in water contained in a tall cylinder. If the depth of water in the cylinder is 0.3 m, calculate the size of the largest particle of sand that can remain suspended after 40 minutes. Assume density of sand = 2600 kg m⁻³ and viscosity of water = 10⁻³ N s m⁻².

Data: \( s = 0.3 \text{ m} , t = 40 \text{ minutes} = 40 × 60 \text{ s} , \rho = 2600 \text{ kg m}⁻³ \)

Solution: Let us assume that the sand particles are spherical in shape and are of different size.
Let \( r \) be the radius of the largest particle.

Terminal velocity \( v = \frac{0.3}{40 × 60} = 1.25 × 10⁻⁴ \text{ m s}⁻¹ \)

Radius \( r = \sqrt{\frac{9ηv}{2(ρ - σ)g}} \)
\[ = \sqrt{\frac{9 × 10⁻³ × 1.25 × 10⁻⁴}{2(2600 - 1000)9.8}} \]
\[ = 5.989 × 10⁻⁶ \text{ m} \]

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5.7 A circular wire loop of 0.03 m radius is rested on the surface of a liquid and then raised. The pull required is 0.003 kg wt greater than the force acting after the film breaks. Find the surface tension of the liquid.

**Solution:** The additional pull $F$ of 0.003 kg wt is the force due to surface tension.

\[ F = T \times \text{length of ring in contact with liquid} \]

(i.e) \[ F = T \times 2 \times \pi r = 4\pi Tr \]

(i.e) \[ 4\pi Tr = F \]

\[ 4\pi Tr = 0.003 \times 9.81 \]

or \[ T = \frac{0.003 \times 9.81}{4 \times 3.14 \times 0.03} = 0.078 \text{ N m}^{-1} \]

5.8 Calculate the diameter of a capillary tube in which mercury is depressed by 2.219 mm. Given $T$ for mercury is 0.54 N m$^{-1}$, angle of contact is 140° and density of mercury is 13600 kg m$^{-3}$

**Data :** $h = -2.219 \times 10^{-3}$ m; $T = 0.54 \text{ N m}^{-1}$; $\theta = 140^\circ$; $\rho = 13600 \text{ kg m}^{-3}$

**Solution :** \[ h\rho g = 2T \cos \theta \]

\[ \therefore r = \frac{2T \cos \theta}{h\rho g} \]

\[ = \frac{2 \times 0.54 \times \cos 140^\circ}{(-2.219 \times 10^{-3}) \times 13600 \times 9.8} \]

\[ = 2.79 \times 10^{-3} \text{ m} \]

**Diameter** = $2r = 2 \times 2.79 \times 10^{-3} = 5.58 \text{ mm}$

5.9 Calculate the energy required to split a water drop of radius $1 \times 10^{-3}$ m into one thousand million droplets of same size. Surface tension of water = 0.072 N m$^{-1}$

**Data :** Radius of big drop $R = 1 \times 10^{-3}$ m

Number of drops $n = 10^3 \times 10^6 = 10^9$; $T = 0.072 \text{ N m}^{-1}$

**Solution :** Let $r$ be the radius of droplet.
Volume of $10^9$ drops = Volume of big drop

$$10^9 \times \frac{4}{3} \pi r^3 = \frac{4}{3} \pi R^3$$

$$10^9 r^3 = R^3 = (10^{-3})^3$$

$$(10^9 r)^3 = (10^{-3})^3$$

$$r = \frac{10^{-3}}{10^3} = 10^{-6} \text{ m}$$

Increase in surface area $ds = 10^9 \times 4 \pi r^2 - 4 \pi R^2$

(i.e) $ds = 4 \pi [ 10^9 \times (10^{-6})^2 - (10^{-3})^2 ] = 4 \pi [10^{-3} - 10^{-6}] \text{ m}^2$

∴ $ds = 0.01254 \text{ m}^2$

Work done $W = T ds = 0.072 \times 0.01254 = 9.034 \times 10^{-4} \text{ J}$

5.10 Calculate the minimum pressure required to force the blood from the heart to the top of the head (a vertical distance of 0.5 m). Given density of blood = 1040 kg m$^{-3}$. Neglect friction.

Data : $h_2 - h_1 = 0.5 \text{ m}$, $\rho = 1040 \text{ kg m}^{-3}$, $P_1 - P_2 = ?$

Solution : According to Bernoulli’s theorem

$$P_1 - P_2 = \rho g (h_2 - h_1) + \frac{1}{2} \rho \left( v_2^2 - v_1^2 \right)$$

If $v_2 = v_1$, then

$$P_1 - P_2 = \rho g (h_2 - h_1)$$

$$P_1 - P_2 = 1040 \times 9.8 \times (0.5)$$

$$P_1 - P_2 = 5.096 \times 10^3 \text{ N m}^{-2}$$
Self evaluation
(The questions and problems given in this self evaluation are only samples. In the same way any question and problem could be framed from the text matter. Students must be prepared to answer any question and problem from the text matter, not only from the self evaluation.)

5.1 If the length of the wire and mass suspended are doubled in a Young’s modulus experiment, then, Young’s modulus of the wire
(a) remains unchanged (b) becomes double
(c) becomes four times (d) becomes sixteen times

5.2 For a perfect rigid body, Young’s modulus is
(a) zero (b) infinity
(c) 1 (d) –1

5.3 Two wires of the same radii and material have their lengths in the ratio 1 : 2. If these are stretched by the same force, the strains produced in the two wires will be in the ratio
(a) 1 : 4 (b) 1 : 2
(c) 2 : 1 (d) 1 : 1

5.4 If the temperature of a liquid is raised, then its surface tension is
(a) decreased (b) increased
(c) does not change (d) equal to viscosity

5.5 The excess of pressure inside two soap bubbles of diameters in the ratio 2 : 1 is
(a) 1 : 4 (b) 2 : 1
(c) 1 : 2 (d) 4 : 1

5.6 A square frame of side l is dipped in a soap solution. When the frame is taken out, a soap film is formed. The force on the frame due to surface tension T of the soap solution is
(a) 8 Tl (b) 4 Tl
(c) 10 Tl (d) 12 Tl
5.7 The rain drops falling from the sky neither hit us hard nor make holes on the ground because they move with
(a) constant acceleration  (b) variable acceleration
(c) variable speed     (d) constant velocity

5.8 Two hail stones whose radii are in the ratio of 1 : 2 fall from a height of 50 km. Their terminal velocities are in the ratio of
(a) 1 : 9           (b) 9 : 1
(c) 4 : 1           (d) 1 : 4

5.9 Water flows through a horizontal pipe of varying cross-section at the rate of 0.2 m$^3$ s$^{-1}$. The velocity of water at a point where the area of cross-section of the pipe is 0.01 m$^2$ is
(a) 2 ms$^{-1}$     (b) 20 ms$^{-1}$
(c) 200 ms$^{-1}$   (d) 0.2 ms$^{-1}$

5.10 An object entering Earth’s atmosphere at a high velocity catches fire due to
(a) viscosity of air  (b) the high heat content of atmosphere
(c) pressure of certain gases  (d) high force of g.

5.11 Define: i) elastic body ii) plastic body iii) stress iv) strain v) elastic limit vi) restoring force

5.12 State Hooke’s law.

5.13 Explain the three moduli of elasticity.

5.14 Describe Searle’s Experiment.

5.15 Which is more elastic, rubber or steel? Support your answer.

5.16 State and prove Pascal’s law without considering the effect of gravity.

5.17 Taking gravity into account, explain Pascal’s law.

5.18 Explain the principle, construction and working of hydraulic brakes.

5.19 What is Reynold’s number?

5.20 What is critical velocity of a liquid?

5.21 Why aeroplanes and cars have streamline shape?

5.22 Describe an experiment to determine viscosity of a liquid.
5.23 What is terminal velocity?
5.24 Explain Stoke’s law.
5.25 Derive an expression for terminal velocity of a small sphere falling through a viscous liquid.
5.26 Define cohesive force and adhesive force. Give examples.
5.27 Define i) molecular range ii) sphere of influence iii) surface tension.
5.28 Explain surface tension on the basis of molecular theory.
5.29 Establish the relation between surface tension and surface energy.
5.30 Give four examples of practical application of surface tension.
5.31 How do insects run on the surface of water?
5.32 Why hot water is preferred to cold water for washing clothes?
5.33 Derive an expression for the total energy per unit mass of a flowing liquid.
5.34 State and prove Bernoulli’s theorem.
5.35 Why the blood pressure in humans is greater at the feet than at the brain?
5.36 Why two holes are made to empty an oil tin?
5.37 A person standing near a speeding train has a danger of falling towards the train. Why?
5.38 Why a small bubble rises slowly through a liquid whereas the bigger bubble rises rapidly?

Problems

5.39 A wire of diameter 2.5 mm is stretched by a force of 980 N. If the Young’s modulus of the wire is $12.5 \times 10^{10} \text{ N m}^{-2}$, find the percentage increase in the length of the wire.

5.40 Two wires are made of same material. The length of the first wire is half of the second wire and its diameter is double that of second wire. If equal loads are applied on both the wires, find the ratio of increase in their lengths.

5.41 The diameter of a brass rod is 4 mm. Calculate the stress and strain when it is stretched by 0.25% of its length. Find the force exerted. Given $q = 9.2 \times 10^{10} \text{ N m}^{-2}$ for brass.
5.42 Calculate the volume change of a solid copper cube, 40 mm on each side, when subjected to a pressure of $2 \times 10^7$ Pa. Bulk modulus of copper is $1.25 \times 10^{11}$ N m$^{-2}$.

5.43 In a hydraulic lift, the piston $P_2$ has a diameter of 50 cm and that of $P_1$ is 10 cm. What is the force on $P_2$ when 1 N of force is applied on $P_1$?

5.44 Calculate the mass of water flowing in 10 minutes through a tube of radius $10^{-2}$ m and length 1 m having a constant pressure of 0.2 m of water. Assume coefficient of viscosity of water = $9 \times 10^{-4}$ N s m$^{-2}$ and $g = 9.8$ m s$^{-2}$.

5.45 A liquid flows through a pipe of $10^{-3}$ m radius and 0.1 m length under a pressure of $10^3$ Pa. If the coefficient of viscosity of the liquid is $1.25 \times 10^{-3}$ N s m$^{-2}$, calculate the rate of flow and the speed of the liquid coming out of the pipe.

5.46 For cylindrical pipes, Reynold’s number is nearly 2000. If the diameter of a pipe is 2 cm and water flows through it, determine the velocity of the flow. Take $\eta$ for water = $10^{-3}$ N s m$^{-2}$.

5.47 In a Poiseuille’s flow experiment, the following are noted.

i) Volume of liquid discharged per minute = $15 \times 10^{-6}$ m$^3$

ii) Head of liquid = 0.30 m

iii) Length of tube = 0.25 m

iv) Diameter = $2 \times 10^{-3}$ m

v) Density of liquid = 2300 kg m$^{-3}$.

Calculate the coefficient of viscosity.

5.48 An air bubble of 0.01 m radius raises steadily at a speed of $5 \times 10^{-3}$ m s$^{-1}$ through a liquid of density 800 kg m$^{-3}$. Find the coefficient of viscosity of the liquid. Neglect the density of air.

5.49 Calculate the viscous force on a ball of radius 1 mm moving through a liquid of viscosity 0.2 N s m$^{-2}$ at a speed of 0.07 m s$^{-1}$.

5.50 A U shaped wire is dipped in soap solution. The thin soap film formed between the wire and a slider supports a weight of $1.5 \times 10^{-2}$ N. If the length of the slider is 30 cm, calculate the surface tension of the film.
5.51 Calculate the force required to remove a flat circular plate of radius 0.02 m from the surface of water. Assume surface tension of water is 0.07 N m$^{-1}$.

5.52 Find the work done in blowing up a soap bubble from an initial surface area of $0.5 \times 10^{-4}$ m$^2$ to an area $1.1 \times 10^{-4}$ m$^2$. The surface tension of soap solution is 0.03 N m$^{-1}$.

5.53 Determine the height to which water will rise in a capillary tube of $0.5 \times 10^{-3}$ m diameter. Given for water, surface tension is 0.074 N m$^{-1}$.

5.54 A capillary tube of inner diameter 4 mm stands vertically in a bowl of mercury. The density of mercury is 13,500 kg m$^{-3}$ and its surface tension is 0.544 N m$^{-1}$. If the level of mercury in the tube is 2.33 mm below the level outside, find the angle of contact of mercury with glass.

5.55 A capillary tube of inner radius $5 \times 10^{-4}$ m is dipped in water of surface tension 0.075 N m$^{-1}$. To what height is the water raised by the capillary action above the water level outside. Calculate the weight of water column in the tube.

5.56 What amount of energy will be liberated if 1000 droplets of water, each of diameter $10^{-8}$ m, coalesce to form a big drop. Surface tension of water is 0.075 N m$^{-1}$.

5.57 Water flows through a horizontal pipe of varying cross-section. If the pressure of water equals $2 \times 10^{-2}$ m of mercury where the velocity of flow is $32 \times 10^{-2}$ m s$^{-1}$ find the pressure at another point, where the velocity of flow is $40 \times 10^{-2}$ m s$^{-1}$. 
Answers

<table>
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<tr>
<th>Question</th>
<th>Answer</th>
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**Mathematical Notes**

(Not for examination)

**Logarithm**

In physics, a student is expected to do the calculation by using logarithm tables. The logarithm of any number to a given base is the power to which the base must be raised in order to obtain the number. For example, we know that 2 raised to power 3 is equal to 8 (i.e) \(2^3 = 8\). In the logarithm form this fact is stated as the logarithm of 8 to the base 2 is equal to 3. (i.e.) \(\log_2 8 = 3\).

In general, if \(a^x = N\), then \(\log_a N = x\).

We use “common logarithm” for calculation purposes. Common logarithm of a number is the power to which 10 must be raised in order to obtain that number. The base 10 is usually not mentioned. In other words, when base is not mentioned, it is understood as base of 10.

For doing calculations with log tables, the following formulae should be kept in mind.

(i) **Product formula**: \(\log mn = \log m + \log n\)

(ii) **Quotient formula**: \(\log \frac{m}{n} = \log m - \log n\)

(iii) **Power formula**: \(\log m^n = n \log m\)

(iv) **Base changing formula**: \(\log_a m = \log_b m \times \log_a b\)

Logarithm of a number consists of two parts called characteristic and Mantissa. The integral part of the logarithm of a number after expressing the decimal part as a positive is called characteristic. The positive decimal part is called Mantissa.

**To find the characteristic of a number**

(i) The characteristic of a number greater than one or equal to one is lesser by one (i.e) \((n - 1)\) than the number of digits \((n)\) present to the left of the decimal point in a given number.

(ii) The characteristic of a number less than one is a negative number whose numerical value is more by one i.e. \((n+1)\) than
the number of zeroes \((n)\) between the decimal point and the first significant figure of the number.

<table>
<thead>
<tr>
<th>Example</th>
<th>Number</th>
<th>Characteristic</th>
</tr>
</thead>
<tbody>
<tr>
<td>5678.9</td>
<td></td>
<td>3</td>
</tr>
<tr>
<td>567.89</td>
<td></td>
<td>2</td>
</tr>
<tr>
<td>56.789</td>
<td></td>
<td>1</td>
</tr>
<tr>
<td>5.6789</td>
<td></td>
<td>0</td>
</tr>
<tr>
<td>0.56789</td>
<td></td>
<td>1</td>
</tr>
<tr>
<td>0.056789</td>
<td></td>
<td>2</td>
</tr>
<tr>
<td>0.0056789</td>
<td></td>
<td>3</td>
</tr>
</tbody>
</table>

**To find the Mantissa of a number**

We have to find out the Mantissa from the logarithm table. The position of a decimal point is immaterial for finding the Mantissa. (i.e) \(\log 39\), \(\log 0.39\), \(\log 0.039\) all have same Mantissa. We use the following procedure for finding the Mantissa.

(i) For finding the Mantissa of \(\log 56.78\), the decimal point is ignored. We get 5678. It can be noted that the first two digits from the left form 56, the third digit is 7 and the fourth is 8.

(ii) In the log tables, proceed in the row 56 and in this row find the number written under the column headed by the third digit 7. (i.e) 7536. To this number the mean difference written under the fourth digit 8 in the same row is added (i.e) \(7536 + 6 = 7542\). Hence logarithm of 56.78 is 1.7542. 1 is the characteristic and 0.7542 is the Mantissa.

(iii) To find out the Mantissa of 567, find the number in the row headed by 56 and under the column 7. It is 7536. Hence the logarithm of 567 is 2. 7536. Here 2 is the characteristic and 0.7536 is the Mantissa.

(iv) To find out the Mantissa of 56, find the number in the row headed by 56 and under the column 0. It is 7482. Hence the logarithm of 56 is 1.7482. Here 1 is the characteristic and 0.7482 is the Mantissa.

(v) To find out the Mantissa of 5, find the number in the row headed by 50 and under the column 0. It is 6990. Hence the logarithm
of 5 is 0.6990. Here 0 is the characteristic and 0.6990 is the Mantissa.

**Antilogarithm**

To find out the antilogarithm of a number, we use the decimal part of a number and read the antilogarithm table in the same manner as in the case of logarithm.

(i) If the characteristic is \( n \), then the decimal point is fixed after \((n+1)\)th digit.

(ii) If the characteristic is \( n \), then add \((n-1)\) zeroes to the left side and then fix the decimal point.

(iii) In general if the characteristic is \( n \) or \( \bar{n} \), then fix the decimal point right side of the first digit and multiply the whole number by \( 10^n \) or \( 10^{-n} \).

<table>
<thead>
<tr>
<th>Example</th>
<th>Number</th>
<th>Antilogarithm</th>
</tr>
</thead>
<tbody>
<tr>
<td>0.9328</td>
<td>8.567 or 8.567 × 10^0</td>
<td></td>
</tr>
<tr>
<td>1.9328</td>
<td>85.67 or 8.567 × 10^1</td>
<td></td>
</tr>
<tr>
<td>2.9328</td>
<td>856.7 or 8.567 × 10^2</td>
<td></td>
</tr>
<tr>
<td>3.9328</td>
<td>8567.0 or 8.567 × 10^3</td>
<td></td>
</tr>
<tr>
<td>1.9328</td>
<td>0.8567 or 8.567 × 10^{-1}</td>
<td></td>
</tr>
<tr>
<td>2.9328</td>
<td>0.08567 or 8.567 × 10^{-2}</td>
<td></td>
</tr>
<tr>
<td>3.9328</td>
<td>0.008567 or 8.567 × 10^{-3}</td>
<td></td>
</tr>
</tbody>
</table>

**EXERCISE - 1**

1. Expand by using logarithm formula
   
   (i) \( T = 2\pi \sqrt{\frac{L}{g}} \)  
   (ii) \( v_c = \sqrt{2gR} \)  
   (iii) \( q = \frac{mg}{\pi r^2 x} \)  
   (iv) \( \log_2 \)

2. Multiply \( \frac{5}{7} .5670 \) by 3

3. Divide \( \frac{5}{7} .6990 \) by 2

4. Evaluate using logarithm
   
   (i) \( \frac{2\times22\times6400}{7\times7918.4} \)  
   (ii) \( \sqrt[3]{9.8\times6370\times10^3} \)
Some commonly used formulae of algebra

(i) \((a+b)^2 = a^2 + 2ab + b^2\)

(ii) \((a-b)^2 = a^2 - 2ab + b^2\)

(iii) \((a+b+c)^2 = a^2 + b^2 + c^2 + 2ab + 2bc + 2ca\)

(iv) \((a+b)^3 = a^3 + b^3 + 3a^2b + 3ab^2\)

(v) \((a - b)^3 = a^3 - b^3 - 3a^2b + 3ab^2\)

**Quadratic equation**

An algebraic equation in the form \(ax^2 + bx + c = 0\) is called quadratic equation. Here \(a\) is the coefficient of \(x^2\), \(b\) is the coefficient of \(x\) and \(c\) is the constant. The solution of the quadratic equation is

\[
x = \frac{-b \pm \sqrt{b^2 - 4ac}}{2a}
\]

**Binomial theorem**

The theorem states that \((1 + x)^n = 1 + nx + \frac{n(n-1)}{2!} x^2 + \frac{n(n-1)(n-2)}{3!} x^3 + \ldots\) where \(x\) is less than 1 and \(n\) is any number. If \(n\) is a positive integer the expansion will have \((n+1)\) terms and if \(n\) is negative or fraction, the expansion will have infinite terms.

Factorial 2 = 2! = 2 \times 1

Factorial 3 = 3! = 3 \times 2 \times 1

Factorial \(n\) = \(n!\) = \(n(n-1)(n-2)\) ....

If \(x\) is very small, then the terms with higher powers of \(x\) can be neglected.

(i.e) \((1 + x)^n = 1 + nx\) \quad (1 + x)^{-n} = 1 - nx

(1 - x)^n = 1 - nx \quad (1 - x)^{-n} = 1 + nx

**Exercise- 2**

1. Find the value of \(x\) in \(4x^2 + 5x - 2 = 0\)

2. Expand Binomially \((i) \left[1 + \frac{h}{R}\right]^2 \quad \text{(ii)} (1 - 2x)^3\)
**Trigonometry**

Let the line AC moves in anticlockwise direction from the initial position AB. The amount of revolution that the moving line makes with its initial position is called angle. From the figure \( \theta = \frac{\angle CAB}{\text{AB}} \). The angle is measured with degree and radian. Radian is the angle subtended at the centre of a circle by an arc of the circle, whose length is equal to the radius of the circle.

1 radian = 57° 17′ 45″
1 right angle = \( \frac{\pi}{2} \) radian
1° = 60′ (sixty minutes), 1′ = 60″ (sixty seconds)

\[
1 \text{ radian} = 57° 17′ 45″
1 \text{ right angle} = \frac{\pi}{2} \text{ radian}
1° = 60′ (sixty minutes), 1′ = 60″ (sixty seconds)
\]

**Triangle laws of sine and cosine**

\[
a^2 = b^2 + c^2 - 2bc \cos \alpha
\]

\[
\frac{a}{\sin \alpha} = \frac{b}{\sin \beta} = \frac{c}{\sin \gamma}
\]

**Trigonometrical ratios (T – ratios)**

Consider the line OA making an angle \( \theta \) in anticlockwise direction with OX.

From A, draw the perpendicular AB to OX.

The longest side of the right angled triangle, OA is called hypotenuse. The side AB is called perpendicular or opposite side. The side OB is called base or adjacent side.

1. Sine of angle \( \theta \) = \( \sin \theta = \frac{\text{perpendicular}}{\text{hypotenuse}} \)
2. Cosine of angle \( \theta \) = \( \cos \theta = \frac{\text{base}}{\text{hypotenuse}} \)
3. Tangent of angle \( \theta \) = \( \tan \theta = \frac{\text{perpendicular}}{\text{base}} \)
4. Cotangent of \( \theta \) = \( \cot \theta = \frac{\text{base}}{\text{perpendicular}} \)
5. Secant of $\theta = \sec \theta = \frac{\text{hypotenuse}}{\text{base}}$

6. Cosecant of $\theta = \csc \theta = \frac{\text{hypotenuse}}{\text{perpendicular}}$

**Sign of trigonometrical ratios**

<table>
<thead>
<tr>
<th>II quadrant</th>
<th>I quadrant</th>
</tr>
</thead>
<tbody>
<tr>
<td>$\sin \theta$ and $\csc \theta$</td>
<td>All positive</td>
</tr>
<tr>
<td>only positive</td>
<td></td>
</tr>
</tbody>
</table>

<table>
<thead>
<tr>
<th>III quadrant</th>
<th>IV quadrant</th>
</tr>
</thead>
<tbody>
<tr>
<td>$\tan \theta$ and $\cot \theta$</td>
<td>$\cos \theta$ and $\sec \theta$ only positive</td>
</tr>
<tr>
<td>only positive</td>
<td></td>
</tr>
</tbody>
</table>

**$T =$ ratios of allied angles**

$- \theta, 90^\circ - \theta, 90^\circ + \theta, 180^\circ - \theta, 180^\circ + \theta, 270^\circ - \theta, 270^\circ + \theta$ are called allied angles to the angle $\theta$. The allied angles are always integral multiples of $90^\circ$.

1. (a) $\sin (-\theta) = -\sin \theta$  
   (b) $\cos (-\theta) = \cos \theta$  
   (c) $\tan (-\theta) = -\tan \theta$

2. (a) $\sin (90^\circ - \theta) = \cos \theta$  
   (b) $\cos (90^\circ - \theta) = \sin \theta$  
   (c) $\tan (90^\circ - \theta) = \cot \theta$

3. (a) $\sin (90^\circ + \theta) = \cos \theta$  
   (b) $\cos (90^\circ + \theta) = -\sin \theta$  
   (c) $\tan (90^\circ + \theta) = -\cot \theta$

4. (a) $\sin (180^\circ - \theta) = \sin \theta$  
   (b) $\cos (180^\circ - \theta) = -\cos \theta$  
   (c) $\tan (180^\circ - \theta) = -\tan \theta$

5. (a) $\sin (180^\circ + \theta) = -\sin \theta$  
   (b) $\cos (180^\circ + \theta) = -\cos \theta$  
   (c) $\tan (180^\circ + \theta) = tan \theta$

6. (a) $\sin (270^\circ - \theta) = -\cos \theta$  
   (b) $\cos (270^\circ - \theta) = -\sin \theta$  
   (c) $\tan (270^\circ - \theta) = \cot \theta$

7. (a) $\sin (270^\circ + \theta) = -\cos \theta$  
   (b) $\cos (270^\circ + \theta) = \sin \theta$  
   (c) $\tan (270^\circ + \theta) = -\cot \theta$
### T-ratios of some standard angles

<table>
<thead>
<tr>
<th>Angle</th>
<th>0°</th>
<th>30°</th>
<th>45°</th>
<th>60°</th>
<th>90°</th>
<th>120°</th>
<th>180°</th>
</tr>
</thead>
<tbody>
<tr>
<td>sin θ</td>
<td>0</td>
<td>1/2</td>
<td>1/√2</td>
<td>√3/2</td>
<td>1</td>
<td>√3/2</td>
<td>0</td>
</tr>
<tr>
<td>cos θ</td>
<td>1</td>
<td>√3/2</td>
<td>1/√2</td>
<td>1/2</td>
<td>0</td>
<td>-1/2</td>
<td>-1</td>
</tr>
<tr>
<td>tan θ</td>
<td>0</td>
<td>√3</td>
<td>1</td>
<td>∞</td>
<td>-1/3</td>
<td>-1</td>
<td>0</td>
</tr>
</tbody>
</table>

### Some trigonometric formulae

1. \( \sin (A + B) = \sin A \cos B + \cos A \sin B \)
2. \( \cos (A + B) = \cos A \cos B - \sin A \sin B \)
3. \( \sin (A - B) = \sin A \cos B - \cos A \sin B \)
4. \( \cos (A - B) = \cos A \cos B + \sin A \sin B \)
5. \( \sin A + \sin B = 2 \sin \frac{A + B}{2} \cos \frac{A - B}{2} \)
6. \( \sin A - \sin B = 2 \cos \frac{A + B}{2} \sin \frac{A - B}{2} \)
7. \( \cos A + \cos B = 2 \cos \frac{A + B}{2} \cos \frac{A - B}{2} \)
8. \( \cos A - \cos B = 2 \sin \frac{A + B}{2} \sin \frac{A - B}{2} \)
9. \( \sin 2A = 2 \sin A \cos A = \frac{2 \tan A}{1 + \tan^2 A} \)
10. \( 2 \sin A \cos B = \sin (A + B) + \sin (A - B) \)
11. \( 2 \cos A \sin B = \sin (A + B) - \sin (A - B) \)
12. \( 2 \sin A \cos B = \cos (A - B) - \cos (A + B) \)
13. \( 2 \cos A \cos B = \cos (A + B) + \cos (A - B) \)
14. \( \cos 2A = 1 - 2 \sin^2 A \)

### Differential calculus

Let \( y \) be the function of \( x \)

(i.e) \( y = f(x) \) \hspace{1cm} \text{....(1)}

The function \( y \) depends on variable \( x \). If the variable \( x \) is changed to \( x + \Delta x \), then the function is also changed to \( y + \Delta y \)
\[ y + \Delta y = f(x + \Delta x) \quad \text{....(2)} \]

Subtracting equation (1) from (2)
\[ \Delta y = f(x + \Delta x) - f(x) \]

dividing on both sides by \( \Delta x \), we get
\[ \frac{\Delta y}{\Delta x} = \frac{f(x + \Delta x) - f(x)}{\Delta x} \]

Taking limits on both sides of equation, when \( \Delta x \) approaches zero, we get
\[ \lim_{\Delta x \to 0} \left( \frac{\Delta y}{\Delta x} \right) = \lim_{\Delta x \to 0} \frac{f(x + \Delta x) - f(x)}{\Delta x} \]

In calculus, \( \lim_{\Delta x \to 0} \frac{\Delta y}{\Delta x} \) is denoted by \( \frac{dy}{dx} \) and is called differentiation of \( y \) with respect to \( x \).

The differentiation of a function with respect to a variable means the instantaneous rate of change of the function with respect to the variable.

**Some theorems and formulae**

1. \( \frac{d}{dx} (c) = 0 \), if \( c \) is a constant.

2. If \( y = cu \), where \( c \) is a constant and \( u \) is a function of \( x \) then
   \[ \frac{dy}{dx} = \frac{d}{dx} (cu) = c \frac{du}{dx} \]

3. If \( y = u \pm v \pm w \) where \( u, v \) and \( w \) are functions of \( x \) then
   \[ \frac{dy}{dx} = \frac{d}{dx} (u \pm v \pm w) = \frac{du}{dx} \pm \frac{dv}{dx} \pm \frac{dw}{dx} \]

4. If \( y = x^n \), where \( n \) is the real number then
   \[ \frac{dy}{dx} = \frac{d}{dx} (x^n) = nx^{n-1} \]

5. If \( y = uv \) where \( u \) and \( v \) are functions of \( x \) then
   \[ \frac{dy}{dx} = \frac{d}{dx} (uv) = u \frac{dv}{dx} + v \frac{du}{dx} \]
6. If \( y \) is a function of \( x \), then \( \frac{dy}{dx} = \frac{dy}{dx} \cdot \frac{dx}{dx} \).

7. \( \frac{d}{dx}(e^x) = e^x \)

8. \( \frac{d}{dx}(\log x) = \frac{1}{x} \)

9. \( \frac{d}{d\theta}(\sin \theta) = \cos \theta \)

10. \( \frac{d}{d\theta}(\cos \theta) = -\sin \theta \)

11. If \( y \) is a trigonometrical function of \( \theta \) and \( \theta \) is the function of \( t \), then

   \[
   \frac{d}{dt}(\sin \theta) = \cos \theta \frac{d\theta}{dt}
   \]

12. If \( y \) is a trigonometrical function of \( \theta \) and \( \theta \) is the function of \( t \), then

   \[
   \frac{d}{dt}(\cos \theta) = -\sin \theta \frac{d\theta}{dt}
   \]

**EXERCISE - 3**

1. If \( y = \sin 3\theta \) find \( \frac{dy}{d\theta} \)

2. If \( y = x^{5/7} \) find \( \frac{dy}{dx} \)

3. If \( y = \frac{1}{x^7} \) find \( \frac{dy}{dx} \)

4. If \( y = 4x^3 + 3x^2 + 2 \), find \( \frac{dy}{dx} \)

5. Differentiate : (i) \( ax^2 + bx + c \)

6. If \( s = 2t^3 - 5t^2 + 4t - 2 \), find the position \( (s) \), velocity \( \left( \frac{ds}{dt} \right) \) and acceleration \( \left( \frac{dv}{dt} \right) \) of the particle at the end of 2 seconds.

**Integration**

It is the reverse process of differentiation. In other words integration is the process of finding a function whose derivative is given. The integral of a function \( y \) with respect to \( x \) is given by \( \int y \, dx \). Integration is represented by the elongated S. The letter S represents the summation of all differential parts.
**Indefinite integral**

We know that
\[
\frac{d}{dx}(x^3) = 3x^2
\]
\[
\frac{d}{dx}(x^3 + 4) = 3x^2
\]
\[
\frac{d}{dx}(x^3 + c) = 3x^2
\]

The result in the above three equations is the same. Hence the question arises as to which of the above results is the integral of \(3x^2\). To overcome this difficulty the integral of \(3x^2\) is taken as \((x^3 + c)\), where \(c\) is an arbitrary constant and can have any value. It is called the constant of integration and is indefinite. The integral containing \(c\), (i.e) \((x^3 + c)\) is called indefinite integral. In practice ‘\(c\)’ is generally not written, though it is always implied.

**Some important formulae**

1. \[
\int dx = x \quad \therefore \frac{d}{dx}(x) = 1
\]
2. \[
\int x^n dx = \left(\frac{x^{n+1}}{n+1}\right) \quad \therefore \frac{d}{dx}\left(\frac{x^{n+1}}{n+1}\right) = x^n
\]
3. \[
\int cu \ dx = c\int u \ dx \quad \text{where} \ c \ \text{is a constant}
\]
4. \[
\int (u \pm v \pm w) \ dx = \int u \ dx \pm \int v \ dx \pm \int w \ dx
\]
5. \[
\int \frac{1}{x} \ dx = \log_c x
\]
6. \[
\int e^t \ dx = e^t
\]
7. \[
\int \cos \theta \ d\theta = \sin \theta
\]
8. \[
\int \sin \theta \ d\theta = -\cos \theta
\]

**Definite integrals**

When a function is integrated between a lower limit and an upper limit, it is called a definite integral.
is a definite integral. Here a and b are lower and upper limits of the variable x.

**EXERCISE – 4**

1. Integrate the following with respect to x

   (i) $4x^3$  
   (ii) $\frac{1}{x^2}$  
   (iii) $3x^2 + 7x - 4$

   (iv) $\frac{5}{7x^{2/3}}$  
   (v) $-\frac{2}{x^3}$  
   (vi) $12x^2 + 6x$

2. Evaluate

   (i) $\int_{\frac{3}{2}}^{2} x^2 \, dx$  
   (ii) $\int_{1}^{\sqrt{x}} x \, dx$

   (iii) $\int_{\frac{4}{2}}^{x} x \, dx$  
   (iv) $\int_{-\pi/2}^{\pi/2} \cos \theta \, d\theta$

**ANSWERS**

**Exercise - 1**

1. (i) $\log 2 + \log 3.14 + \frac{1}{2} \log l - \frac{1}{2} \log g$

   (ii) $\frac{1}{2} (\log 2 + \log g + \log R)$

   (iii) $\log m + \log g + \log l - \log 3.14 - 2 \log r - \log x$

   (iv) 0.6931

2. 14.7010

3. 5.8495

4. (i) 5.080  
   (ii) $7.9 \times 10^3$

   (iii) $1.764 \times 10^{-4}$  
   (iv) $2.836 \times 10^{-1}$
Exercise - 2

(1) \( \frac{-5 \pm \sqrt{57}}{8} \)  
(2) (i) \( 1 - \frac{2h}{R} \)  
               (ii) \( 1 - 6x \)

Exercise - 3

(1) \( 3 \cos 3\theta \)  
(2) \( \frac{5}{7}x^{-2/7} \)  
(3) \( \frac{-2}{x^3} \)  
(4) \( 12x^2 + 6x \)  
(5) \( 2ax + b \)  
(6) \( 2, 8, 14 \)

Exercise - 4

1. (i) \( x^4 \)  
     (ii) \( \frac{1}{x} \)  
     (iii) \( x^3 + \frac{7}{2}x^2 - 4x \)  
     (iv) \( x^{5/7} \)  
     (v) \( \frac{1}{x^2} \)  
     (vi) \( 4x^3 + 3x^2 \)  

2. (i) \( \frac{19}{3} \)  
     (ii) \( \frac{14}{3} \)  
     (iii) \( 6 \)  
     (iv) \( 2 \)
Proof for Lami’s theorem

Let forces \( \vec{P}, \vec{Q} \) and \( \vec{R} \) acting at a point \( O \) be in equilibrium. Let \( OA \) and \( OB(=AD) \) represent the forces \( \vec{P} \) and \( \vec{Q} \) in magnitude and direction. By the parallelogram law of forces \( OD \) will represent the resultant of the forces \( \vec{P} \) and \( \vec{Q} \). Since the forces are in equilibrium \( DO \) will represent the third force \( R \).

In the triangle \( OAD \), using law of sines,

\[
\frac{OA}{\sin \angle ODA} = \frac{AD}{\sin \angle AOD} = \frac{OD}{\sin \angle OAD}
\]

From Fig. 2.35,

\[
\angle ODA = \angle BOD = 180^\circ - \angle BOC
\]

\[
\angle AOD = 180^\circ - \angle AOC
\]

\[
\angle AOD = 180^\circ - \angle AOB
\]

Therefore,

\[
\frac{OA}{\sin (180^\circ - \angle BOC)} = \frac{AD}{\sin (180^\circ - \angle AOC)} = \frac{OD}{\sin (180^\circ - \angle AOB)}
\]

(i.e) \[ \frac{OA}{\sin \angle BOC} = \frac{AD}{\sin \angle AOC} = \frac{OD}{\sin \angle AOB} \]

If \( \angle BOC = \alpha, \angle AOC = \beta, \angle AOB = \gamma \)

\[ \frac{P}{\sin \alpha} = \frac{Q}{\sin \beta} = \frac{R}{\sin \gamma} \] which proves Lami’s theorem.

ANNEXURE

(NOT FOR EXAMINATION)
1. Moment of inertia of a thin uniform rod

(i) About an axis passing through its centre of gravity and perpendicular to its length

Consider a thin uniform rod AB of mass \( M \) and length \( l \) as shown in Fig. 1. Its mass per unit length will be \( \frac{M}{l} \). Let, \( YY' \) be the axis passing through the centre of gravity \( G \) of the rod (and perpendicular to the length \( AB \)).

Consider a small element of length \( dx \) of the rod at a distance \( x \) from \( G \).

The mass of the element

\[
\text{mass of the element} = \text{mass per unit length} \times \text{length of the element} = \frac{M}{l} \times dx \quad \ldots(1)
\]

The moment of inertia of the element \( dx \) about the axis \( YY' \) is,

\[
dl = (\text{mass}) \times (\text{distance})^2 = \left( \frac{M}{l} \right) (x^2) \quad \ldots(2)
\]

Therefore the moment of inertia of the whole rod about \( YY' \) is obtained by integrating equation (2) within the limits \(-\frac{l}{2}\) to \(+\frac{l}{2}\).

\[
I_{CG} = \int_{-\frac{l}{2}}^{\frac{l}{2}} (x^2) dx = \frac{M}{l} \left[ \frac{x^3}{3} \right]_{-\frac{l}{2}}^{\frac{l}{2}} = \frac{M}{l} \left[ \frac{l^3}{3} - \left( \frac{l}{2} \right)^3 - \left( \frac{l}{2} \right)^3 \right] = \frac{M}{3l} \left[ \frac{l^3}{3} - \frac{8}{8} \right] = \frac{M}{3l} \left[ 2\frac{l^3}{8} \right] = \frac{ML^3}{12l} \quad \ldots(3)
\]
(ii) About an axis passing through the end and perpendicular to its length

The moment of inertia \( I \) about a parallel axis \( Y_1Y_1' \) passing through one end \( A \) can be obtained by using parallel axes theorem

\[
\therefore \quad I = I_{CG} + M \left( \frac{l}{2} \right)^2 = \frac{Ml^2}{12} + \frac{Ml^2}{4}
\]

\[
I = \frac{Ml^2}{3}
\]

2 Moment of inertia of a thin circular ring

(i) About an axis passing through its centre and perpendicular to its plane

Let us consider a thin ring of mass \( M \) and radius \( R \) with \( O \) as centre, as shown in Fig. 2. As the ring is thin, each particle of the ring is at a distance \( R \) from the axis \( XOY \) passing through \( O \) and perpendicular to the plane of the ring.

For a particle of mass \( m \) on the ring, its moment of inertia about the axis \( XOY \) is \( mR^2 \). Therefore the moment of inertia of the ring about the axis is,

\[
I = \Sigma mR^2 = (\Sigma m) R^2 = MR^2
\]

(ii) About its diameter

\( AB \) and \( CD \) are the diameters of the ring perpendicular to each other (Fig. 3). Since, the ring is symmetrical about any diameter, its moment of inertia about \( AB \) will be equal to that about \( CD \). Let it be \( I_d \). If \( I \) is the moment of inertia of the ring about an axis passing through the centre and perpendicular to its plane then applying perpendicular axes theorem,

\[
\therefore \quad I = I_d + I_d = MR^2 \quad \text{(or)} \quad I_d = \frac{1}{2} MR^2
\]
(iii) About a tangent

The moment of inertia of the ring about a tangent $EF$ parallel to $AB$ is obtained by using the parallel axes theorem. The moment of inertia of the ring about any tangent is,

$$I_T = I_d + M R^2 = \frac{1}{2} MR^2 + MR^2$$

$$I_T = \frac{3}{2} MR^2$$

3 Moment of inertia of a circular disc

(i) About an axis passing through its centre and perpendicular to its plane

Consider a circular disc of mass $M$ and radius $R$ with its centre at $O$ as shown in Fig. 4. Let $\sigma$ be the mass per unit area of the disc. The disc can be imagined to be made up of a large number of concentric circular rings of radii varying from $O$ to $R$. Let us consider one such ring of radius $r$ and width $dr$.

The circumference of the ring = $2\pi r$.

The area of the elementary ring = $2\pi r dr$

Mass of the ring = $2\pi r \sigma dr$ ... (1)

Moment of inertia of this elementary ring about the axis passing through its centre and perpendicular to its plane is

$$dl = \text{mass} \times (\text{distance})^2$$

$$= (2\pi r \sigma dr) r^2$$ ... (2)

The moment of inertia of the whole disc about an axis passing through its centre and perpendicular to its plane is,

$$I = \int_0^R 2\pi r^3 \sigma dr = 2\pi \sigma \left[ \frac{r^4}{4} \right]_0^R$$

(or) $$I = \frac{2\pi \sigma R^4}{4} = (\pi R^2 \sigma) \frac{1}{2} R^2 = \frac{1}{2} MR^2$$ ... (3)

where $M = \pi R^2 \sigma$ is the mass of the disc.
(ii) **About a diameter**

Since, the disc is symmetrical about any diameter, the moment of inertia about the diameter $AB$ will be same as its moment of inertia about the diameter $CD$. Let it be $I_d$ (Fig. 5). According to perpendicular axes theorem, the moment of inertia $I$ of the disc, about an axis perpendicular to its plane and passing through the centre will be equal to the sum of its moment of inertia about two mutually perpendicular diameters $AB$ and $CD$.

Hence, $I = I_d + I_d = \frac{1}{2} MR^2 + \frac{1}{4} MR^2$

(iii) **About a tangent in its plane**

The moment of inertia of the disc about the tangent $EF$ in the plane of the disc and parallel to $AB$ can be obtained by using the theorem of parallel axes (Fig. 3.15).

$$I_T = I_d + MR^2 = \frac{1}{4} MR^2 + MR^2$$

$$\therefore \quad I_T = \frac{5}{4} MR^2$$

4 **Moment of inertia of a sphere**

(i) **About a diameter**

Let us consider a homogeneous solid sphere of mass $M$, density $\rho$ and radius $R$ with centre $O$ (Fig. 6). $AB$ is the diameter about which the moment of inertia is to be determined. The sphere may be considered as made up of a large number of coaxial circular discs with their centres lying on $AB$ and their planes perpendicular to $AB$. Consider a disc of radius $PO' = y$ and thickness $dx$ with centre $O'$ and at a distance $x$ from $O$. 

![Fig 6 Moment of inertia of a sphere about a diameter](image)
Its volume = $\pi y^2 \, dx$ \hspace{1cm} \ldots(1)

Mass of the disc = $\pi y^2 \, dx \cdot \rho$ \hspace{1cm} \ldots(2)

From Fig. 6, $R^2 = y^2 + x^2$ \hspace{0.5cm} \text{(or)} \hspace{0.5cm} y^2 = R^2 - x^2 \hspace{1cm} \ldots(3)

Using (3) in (2),

Mass of the circular disc = $\pi (R^2 - x^2) \, dx \cdot \rho$ \hspace{1cm} \ldots(4)

The moment of inertia of the disc about the diameter AB is,

\[
dI = \frac{1}{2} \text{(mass)} \times \text{(radius)}^2 = \frac{1}{2} \pi (R^2 - x^2) dx \cdot \rho \left( R^2 - x^2 \right) dx = \frac{1}{2} \pi \rho \left( R^2 - x^2 \right)^2 dx \hspace{1cm} \ldots(5)
\]

The moment of inertia of the entire sphere about the diameter AB is obtained by integrating eqn (5) within the limits $x = -R$ to $x = +R$.

\[
I = \int_{-R}^{R} \frac{1}{2} \pi R^2 \rho \left( R^2 - x^2 \right)^2 dx
\]

\[
I = 2 \times \frac{1}{2} \pi R^2 \rho \int_0^R (R^4 - x^4) \, dx = \pi \rho \left[ R^5 + \frac{R^5}{5} - \frac{2R^5}{3} \right] = \pi \rho \left( \frac{8}{15} R^5 \right) = \frac{4}{3} \pi R^3 \rho \left( \frac{2}{5} R^2 \right) = M \left( \frac{2}{5} R^2 \right) = \frac{2}{5} MR^2
\]

where $M = \frac{4}{3} \pi R^3 \rho$ = mass of the solid sphere

\[
\therefore \quad I = \frac{2}{5} MR^2
\]

\(\text{(ii) About a tangent}\)

The moment of inertia of a solid sphere about a tangent EF parallel to the diameter AB (Fig. 7) can be determined using the parallel axes theorem,
\[ I_t = I_{AB} + MR^2 = \frac{2}{5}MR^2 + MR^2 \]

\[ \therefore I_t = \frac{7}{5}MR^2 \]

5. Moment of inertia of a solid cylinder

(i) about its own axis

Let us consider a solid cylinder of mass \( M \), radius \( R \) and length \( l \). It may be assumed that it is made up of a large number of thin circular discs each of mass \( m \) and radius \( R \) placed one above the other.

Moment of inertia of a disc about an axis passing through its centre but perpendicular to its plane = \( \frac{mR^2}{2} \)

\[ \therefore \] Moment of inertia of the cylinder about its axis \( I = \sum \frac{mR^2}{2} \)

\[ I = \frac{R^2}{2} \left( \sum m \right) = \frac{R^2}{2} M = \frac{MR^2}{2} \]

(ii) About an axis passing through its centre and perpendicular to its length

Mass per unit length of the cylinder = \( \frac{M}{l} \) \hspace{1cm} ...(1)

Let \( O \) be the centre of gravity of the cylinder and \( \text{YOY'} \) be the axis passing through the centre of gravity and perpendicular to the length of the cylinder (Fig. 8).

Consider a small circular disc of width \( dx \) at a distance \( x \) from the axis \( \text{YY'} \).

\[ \therefore \] Mass of the disc = mass per unit length \( \times \) width

\[ = \left( \frac{M}{l} \right) dx \] \hspace{1cm} ...(2)
Moment of inertia of the disc about an axis parallel to YY’ (i.e) about its diameter = (mass) \( \frac{radius^2}{4} \)

\[ = \left( \frac{M}{l} \right) dx \left( \frac{R^2}{4} \right) = \frac{MR^2}{4l} dx \]  

...(3)

By parallel axes theorem, the moment of inertia of this disc about an axis parallel to its diameter and passing through the centre of the cylinder (i.e. about YY’) is

\[ dl = \left( \frac{MR^2}{4l} \right) dx + \left( \frac{M}{l} \right) dx (x^2) \]  

...(4)

Hence the moment of inertia of the cylinder about YY’ is,

\[ I = \int_{-l/2}^{l/2} \left( \frac{MR^2}{4l} \right) dx + \left( \frac{M}{l} \right) \int_{-l/2}^{l/2} x^2 dx \]

\[ I = \frac{MR^2}{4l} \left[ x \right]_{-l/2}^{l/2} + \left( \frac{M}{l} \right) \left[ \frac{x^3}{3} \right]_{-l/2}^{l/2} \]

\[ I = \frac{MR^2}{4l} \left[ \left( \frac{l}{2} \right) - \left( -\frac{l}{2} \right) \right] + \left( \frac{M}{l} \right) \left[ \frac{\left( \frac{l}{2} \right)^3 - \left( \frac{-l}{2} \right)^3}{3} \right] \]

\[ I = \frac{MR^2}{4l} \left( l \right) + \left( \frac{M}{l} \right) \left[ \frac{2l^3}{24} \right] \]

\[ I = \frac{MR^2}{4} + \frac{Ml^2}{12} \]

\[ I = M \left( \frac{R^2}{4} + \frac{l^2}{12} \right) \]  

...(5)
Untouchability is a sin
Untouchability is a crime
Untouchability is inhuman
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Preface

The most important and crucial stage of school education is the higher secondary level. This is the transition level from a generalised curriculum to a discipline-based curriculum.

In order to pursue their career in basic sciences and professional courses, students take up Physics as one of the subjects. To provide them sufficient background to meet the challenges of academic and professional streams, the Physics textbook for Std. XI has been reformed, updated and designed to include basic information on all topics.

Each chapter starts with an introduction, followed by subject matter. All the topics are presented with clear and concise treatments. The chapters end with solved problems and self-evaluation questions.

Understanding the concepts is more important than memorising. Hence it is intended to make the students understand the subject thoroughly so that they can put forth their ideas clearly. In order to make the learning of Physics more interesting, application of concepts in real life situations are presented in this book.

Due importance has been given to develop in the students, experimental and observation skills. Their learning experience would make them to appreciate the role of Physics towards the improvement of our society.

The following are the salient features of the text book.

- The data has been systematically updated.
- Figures are neatly presented.
- Self-evaluation questions (only samples) are included to sharpen the reasoning ability of the student.
- As Physics cannot be understood without the basic knowledge of Mathematics, few basic ideas and formulae in Mathematics are given.

While preparing for the examination, students should not restrict themselves, only to the questions/problems given in the self evaluation. They must be prepared to answer the questions and problems from the text/syllabus.

Sincere thanks to Indian Space Research Organisation (ISRO) for providing valuable information regarding the Indian satellite programme.

– Dr. S. Gunasekaran
Chairperson
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6. Oscillations

Any motion that repeats itself after regular intervals of time is known as a periodic motion. The examples of periodic motion are the motion of planets around the Sun, motion of hands of a clock, motion of the balance wheel of a watch, motion of Halley’s comet around the Sun observable on the Earth once in 76 years.

If a body moves back and forth repeatedly about a mean position, it is said to possess oscillatory motion. Vibrations of guitar strings, motion of a pendulum bob, vibrations of a tuning fork, oscillations of mass suspended from a spring, vibrations of diaphragm in telephones and speaker system and freely suspended springs are few examples of oscillatory motion. In all the above cases of vibrations of bodies, the path of vibration is always directed towards the mean or equilibrium position.

The oscillations can be expressed in terms of simple harmonic functions like sine or cosine function. A harmonic oscillation of constant amplitude and single frequency is called simple harmonic motion (SHM).

6.1 Simple harmonic motion

A particle is said to execute simple harmonic motion if its acceleration is directly proportional to the displacement from a fixed point and is always directed towards that point.

Consider a particle P executing SHM along a straight line between A and B about the mean position O (Fig. 6.1). The acceleration of the particle is always directed towards a fixed point on the line and its magnitude is proportional to the displacement of the particle from this point.

(i.e) \( a \propto y \)

By definition \( a = -\omega^2 y \)

where \( \omega \) is a constant known as angular frequency of the simple harmonic motion. The negative sign indicates that the acceleration is opposite to the direction of displacement. If \( m \) is the mass of the particle, restoring force that tends to bring
back the particle to the mean position is given by

\[ F = -m \omega^2 y \]

or \[ F = -ky \]

The constant \( k = m \omega^2 \), is called force constant or spring constant. Its unit is \( N \ m^{-1} \). The restoring force is directed towards the mean position.

Thus, simple harmonic motion is defined as oscillatory motion about a fixed point in which the restoring force is always proportional to the displacement and directed always towards that fixed point.

6.1.1 The projection of uniform circular motion on a diameter is SHM

Consider a particle moving along the circumference of a circle of radius \( a \) and centre \( O \), with uniform speed \( v \), in anticlockwise direction as shown in Fig. 6.2. Let \( XX' \) and \( YY' \) be the two perpendicular diameters.

Suppose the particle is at \( P \) after a time \( t \). If \( \omega \) is the angular velocity, then the angular displacement \( \theta \) in time \( t \) is given by \( \theta = \omega t \). From \( P \) draw \( PN \) perpendicular to \( YY' \). As the particle moves from \( X \) to \( Y \), foot of the perpendicular \( N \) moves from \( O \) to \( Y \). As it moves further from \( Y \) to \( X' \), then from \( X' \) to \( Y' \) and back again to \( X \), the point \( N \) moves from \( Y \) to \( O \), from \( O \) to \( Y' \) and back again to \( O \). When the particle completes one revolution along the circumference, the point \( N \) completes one vibration about the mean position \( O \). The motion of the point \( N \) along the diameter \( YY' \) is simple harmonic.

Hence, the projection of a uniform circular motion on a diameter of a circle is simple harmonic motion.

Displacement in SHM

The distance travelled by the vibrating particle at any instant of time \( t \) from its mean position is known as displacement. When the particle is at \( P \) the displacement of the particle along \( Y \) axis is \( y \) (Fig. 6.3).
Then, in $\Delta$ OPN, $\sin \theta = \frac{ON}{OP}$

$$ON = y = OP \sin \theta$$

$$y = OP \sin \omega t \quad (\because \theta = \omega t)$$

since $OP = a$, the radius of the circle, the displacement of the vibrating particle is

$$y = a \sin \omega t \quad ...(1)$$

The amplitude of the vibrating particle is defined as its maximum displacement from the mean position.

**Velocity in SHM**

The rate of change of displacement is the velocity of the vibrating particle.

Differentiating eqn. (1) with respect to time $t$

$$\frac{dy}{dt} = \frac{d}{dt} (a \sin \omega t)$$

$$\therefore \; v = a \omega \cos \omega t \quad ...(2)$$

The velocity $v$ of the particle moving along the circle can also be obtained by resolving it into two components as shown in Fig. 6.4.

(i) $v \cos \theta$ in a direction parallel to OY

(ii) $v \sin \theta$ in a direction perpendicular to OY

The component $v \sin \theta$ has no effect along $YOY'$ since it is perpendicular to OY.

\[\therefore \text{Velocity} = v \cos \theta\]

\[= v \cos \omega t\]

We know that, linear velocity = radius $\times$ angular velocity

\[\therefore \; v = a \omega\]

\[\therefore \text{Velocity} = a \omega \cos \omega t\]

\[\therefore \text{Velocity} = a \omega \sqrt{1 - \sin^2 \omega t}\]
Velocity = $a\omega \sqrt{1 - \left(\frac{y}{a}\right)^2}$ $\therefore \sin \theta = \frac{y}{a}$

Velocity = $\omega \sqrt{a^2 - y^2}$ ...(3)

**Special cases**

(i) When the particle is at mean position, (i.e) $y = 0$. Velocity is $ao$ and is maximum. $v = \pm ao$ is called *velocity amplitude.*

(ii) When the particle is in the extreme position, (i.e) $y = \pm a$, the velocity is zero.

**Acceleration in SHM**

The rate of change of velocity is the acceleration of the vibrating particle.

$$\frac{d^2 y}{dt^2} = \frac{d}{dt} \left( \omega \cos \omega t \right) = -\omega^2 \sin \omega t.$$  

∴ acceleration = $\frac{d^2 y}{dt^2} = -\omega^2 y$ ... (4)

The acceleration of the particle can also be obtained by component method.

The centripetal acceleration of the particle $P$ acting along $PO$ is $\frac{v^2}{a}$. This acceleration is resolved into two components as shown in Fig. 6.5.

(i) $\frac{v^2}{a} \cos \theta$ along $PN$ perpendicular to $OY$

(ii) $\frac{v^2}{a} \sin \theta$ in a direction parallel to $YO$
The component $\frac{v^2}{a} \cos \theta$ has no effect along $YOY'$ since it is perpendicular to $OY$.

Hence acceleration $= -\frac{v^2}{a} \sin \theta$

$= -a \omega^2 \sin \omega t \quad (\therefore v = a \omega)$

$= -\omega^2 y \quad (\therefore y = a \sin \omega t)$

$\therefore$ acceleration $= -\omega^2 y$

The negative sign indicates that the acceleration is always opposite to the direction of displacement and is directed towards the centre.

**Special Cases**

(i) When the particle is at the mean position (i.e) $y = 0$, the acceleration is zero.

(ii) When the particle is at the extreme position (i.e) $y = \pm a$, acceleration is $\mp a \omega^2$ which is called as acceleration amplitude.

The differential equation of simple harmonic motion from eqn. (4) is

$$\frac{d^2y}{dt^2} + \omega^2 y = 0 \quad \ldots(5)$$

Using the above equations, the values of displacement, velocity and acceleration for the SHM are given in the Table 6.1.

It will be clear from the above, that at the mean position $y = 0$, velocity of the particle is maximum but acceleration is zero. At extreme

<table>
<thead>
<tr>
<th>Time</th>
<th>$\omega t$</th>
<th>Displacement $a \sin \omega t$</th>
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<td>$t = 0$</td>
<td>0</td>
<td>0</td>
<td>$a\omega$</td>
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<td>$t = \frac{T}{4}$</td>
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<td>$-a$</td>
<td>0</td>
<td>$+a\omega^2$</td>
</tr>
<tr>
<td>$t = T$</td>
<td>$2\pi$</td>
<td>0</td>
<td>$+a\omega$</td>
<td>0</td>
</tr>
</tbody>
</table>
position \( y = \pm a \), the velocity is zero but the acceleration is maximum \( \mp a \omega^2 \) acting in the opposite direction.

**Graphical representation of SHM**

Graphical representation of displacement, velocity and acceleration of a particle vibrating simple harmonically with respect to time \( t \) is shown in Fig. 6.6.

(i) Displacement graph is a sine curve. Maximum displacement of the particle is \( y = \pm a \).

(ii) The velocity of the vibrating particle is maximum at the mean position i.e \( v = \pm a \omega \) and it is zero at the extreme position.

(iii) The acceleration of the vibrating particle is zero at the mean position and maximum at the extreme position (i.e) \( \mp a \omega^2 \).

The velocity is ahead of displacement by a phase angle of \( \frac{\pi}{2} \). The acceleration is ahead of the velocity by a phase angle \( \frac{\pi}{2} \) or by a phase \( \pi \) ahead of displacement. (i.e) when the displacement has its greatest positive value, acceleration has its negative maximum value or vice versa.

### 6.2 Important terms in simple harmonic motion

(i) **Time period**

The time taken by a particle to complete one oscillation is called the time period \( T \).

In the Fig. 6.2, as the particle P completes one revolution with angular velocity \( \omega \), the foot of the perpendicular \( N \) drawn to the vertical diameter completes one vibration. Hence \( T \) is the time period.
Then $\omega = \frac{2\pi}{T}$ or $T = \frac{2\pi}{\omega}$

The displacement of a particle executing simple harmonic motion may be expressed as

$$y(t) = a \sin \left( \frac{2\pi}{T} t \right)$$

...(1)

and

$$y(t) = a \cos \left( \frac{2\pi}{T} t \right)$$

...(2)

where $T$ represents the time period, $a$ represents maximum displacement (amplitude).

These functions repeat when $t$ is replaced by $(t + T)$.

$$y(t + T) = a \sin \left( \frac{2\pi}{T} (t + T) \right)$$

$$= a \sin \left( 2\pi + \frac{2\pi}{T} t \right)$$

$$= a \sin \left( 2\pi \left( \frac{t}{T} + 1 \right) \right)$$

In general $y(t + nT) = y(t)$

Above functions are examples of periodic function with time period $T$. It is clear that the motion repeats after a time $T = \frac{2\pi}{\omega}$ where $\omega$ is the angular frequency of the motion. In one revolution, the angle covered by a particle is $2\pi$ in time $T$.

**Frequency and angular frequency**

The number of oscillations produced by the body in one second is known as frequency. It is represented by $n$. The time period to complete one oscillation is $\frac{1}{n}$.

$$T = \frac{1}{n}$$ shows the time period is the reciprocal of the frequency. Its unit is hertz, $\omega = 2\pi n$, is called as angular frequency. It is expressed in rad s$^{-1}$. 

7
(iii) Phase

The phase of a particle vibrating in SHM is the state of the particle as regards to its direction of motion and position at any instant of time. In the equation \( y = a \sin (\omega t + \phi_0) \) the term \((\omega t + \phi_0) = \phi\) is known as the phase of the vibrating particle.

**Epoch**

It is the initial phase of the vibrating particle (i.e) phase at \( t = 0 \).

\[ \therefore \phi = \phi_0 \quad (\because \phi = \omega t + \phi_0) \]

The phase of a vibrating particle changes with time but the epoch is phase constant.

![Fig. 6.7 Phase](image)

(a) Phase \( \phi = (\omega t - \phi_0) \) 
(b) Phase \( \phi = (\omega t + \phi_0) \)

(i) If the particle \( P \) starts from the position \( X \), the phase of the particle is Zero.

(ii) Instead of counting the time from the instant the particle is at \( X \), it is counted from the instant when the reference particle is at \( A \) (Fig. 6.7a). Then \[ \angle XOP = (\omega t - \phi_0) \].

Here \((\omega t - \phi_0) = \phi \) is called the phase of the vibrating particle. \((-\phi_0)\) is initial phase or epoch.

(iii) If the time is counted from the instant the particle \( P \) is above \( X \) (i.e) at \( B \), [Fig. 6.7b] then \((\omega t + \phi_0) = \phi \). Here \((+\phi_0)\) is the initial phase.
**Phase difference**

If two vibrating particles executing SHM with same time period, cross their respective mean positions at the same time in the same direction, they are said to be in phase.

If the two vibrating particles cross their respective mean position at the same time but in the opposite direction, they are said to be out of phase (i.e they have a phase difference of $\pi$).

If the vibrating motions are represented by equations

\[ y_1 = a \sin \omega t \]
\[ y_2 = a \sin (\omega t - \phi) \]

then the phase difference between their phase angles is equal to the phase difference between the two motions.

\[ \therefore \text{phase difference} = \omega t - \phi - \omega t = -\phi \]

negative sign indicates that the second motion lags behind the first.

If \[ y_2 = a \sin (\omega t + \phi), \]
\[ \text{phase difference} = \omega t + \phi - \omega t = \phi \]

Here the second motion leads the first motion.

We have discussed the SHM without taking into account the cause of the motion which can be a force (linear SHM) or a torque (angular SHM).

**Some examples of SHM**

(i) Horizontal and vertical oscillations of a loaded spring.
(ii) Vertical oscillation of water in a U-tube
(iii) Oscillations of a floating cylinder
(iv) Oscillations of a simple pendulum
(v) Vibrations of the prongs of a tuning fork.

**6.3 Dynamics of harmonic oscillations**

The oscillations of a physical system results from two basic properties namely elasticity and inertia. Let us consider a body displaced from a mean position. The restoring force brings the body to the mean position.

(i) At extreme position when the displacement is maximum, velocity is zero. The acceleration becomes maximum and directed towards the mean position.
(ii) Under the influence of restoring force, the body comes back to the mean position and overshoots because of negative velocity gained at the mean position.

(iii) When the displacement is negative maximum, the velocity becomes zero and the acceleration is maximum in the positive direction. Hence the body moves towards the mean position. Again when the displacement is zero in the mean position velocity becomes positive.

(iv) Due to inertia the body overshoots the mean position once again. This process repeats itself periodically. Hence the system oscillates.

The restoring force is directly proportional to the displacement and directed towards the mean position.

(i.e) \( F \propto y \)

\[ F = -ky \] ... (1)

where \( k \) is the force constant. It is the force required to give unit displacement. It is expressed in \( \text{N m}^{-1} \).

From Newton’s second law, \( F = ma \)  

\[ \therefore -ky = ma \]

or \( a = \frac{k}{m} y \)  

...(3)

From definition of SHM acceleration \( a = -\omega^2 y \)

The acceleration is directly proportional to the negative of the displacement.

Comparing the above equations we get,

\[ \omega = \sqrt{\frac{k}{m}} \] ... (4)

Therefore the period of SHM is

\[ T = \frac{2\pi}{\omega} = 2\pi \sqrt{\frac{m}{k}} \]

\[ T = 2\pi \sqrt{\frac{\text{inertial factor}}{\text{spring factor}}} \] ... (5)
6.4 Angular harmonic oscillator

Simple harmonic motion can also be angular. In this case, the restoring torque required for producing SHM is directly proportional to the angular displacement and is directed towards the mean position.

Consider a wire suspended vertically from a rigid support. Let some weight be suspended from the lower end of the wire. When the wire is twisted through an angle $\theta$ from the mean position, a restoring torque acts on it tending to return it to the mean position. Here restoring torque is proportional to angular displacement $\theta$.

Hence $\tau = -C\theta$ \hspace{1cm} ...(1)

where $C$ is called torque constant.

It is equal to the moment of the couple required to produce unit angular displacement. Its unit is N m rad$^{-1}$.

The negative sign shows that torque is acting in the opposite direction to the angular displacement. This is the case of angular simple harmonic motion.

Examples: Torsional pendulum, balance wheel of a watch.

But $\tau = I\alpha$ \hspace{1cm} ...(2)

where $\tau$ is torque, $I$ is the moment of inertia and $\alpha$ is angular acceleration

:. Angular acceleration, $\alpha = \frac{\tau}{I} = \frac{C\theta}{I}$ \hspace{1cm} ...(3)

This is similar to $a = -\omega^2 y$

Replacing $y$ by $\theta$, and $a$ by $\alpha$ we get

$\alpha = -\omega^2 \theta = -\frac{C}{I} \theta$

:. $\omega = \sqrt{\frac{C}{I}}$

:. Period of SHM $T = 2\pi\sqrt{\frac{I}{C}}$

:. Frequency $f = \frac{1}{T} = \frac{1}{2\pi\sqrt{\frac{I}{C}}} = \frac{1}{2\pi} \sqrt{\frac{C}{I}}$

\[\text{Fig. 6.8 Torsional Pendulum}\]
6.5 Linear simple harmonic oscillator

The block – spring system is a linear simple harmonic oscillator. All oscillating systems like diving board, violin string have some element of springiness, \( k \) (spring constant) and some element of inertia, \( m \).

6.5.1 Horizontal oscillations of spring

Consider a mass \( (m) \) attached to an end of a spiral spring (which obeys Hooke’s law) whose other end is fixed to a support as shown in Fig. 6.9. The body is placed on a smooth horizontal surface. Let the body be displaced through a distance \( x \) towards right and released. It will oscillate about its mean position. The restoring force acts in the opposite direction and is proportional to the displacement.

\[ \therefore \text{Restoring force } F = -kx. \]

From Newton’s second law, we know that \( F = ma \)

\[ \therefore ma = -kx \]

\[ a = \frac{-k}{m} x \]

Comparing with the equation of SHM \( a = -\omega^2 x \), we get

\[ \omega^2 = \frac{k}{m} \]

or \( \omega = \sqrt{\frac{k}{m}} \)

But \( T = \frac{2\pi}{\omega} \)

Time period \( T = 2\pi \sqrt{\frac{m}{k}} \)

\[ \therefore \text{Frequency } n = \frac{1}{T} = \frac{1}{2\pi} \sqrt{\frac{k}{m}} \]
6.5.2 Vertical oscillations of a spring

Fig 6.10a shows a light, elastic spiral spring suspended vertically from a rigid support in a relaxed position. When a mass \( m \) is attached to the spring as in Fig. 6.10b, the spring is extended by a small length \( dl \) such that the upward force \( F \) exerted by the spring is equal to the weight \( mg \).

The restoring force 
\[
F = k \, dl \quad ; \quad k \, dl = mg \quad \cdots (1)
\]
where \( k \) is spring constant. If we further extend the given spring by a small distance by applying a small force by our finger, the spring oscillates up and down about its mean position. Now suppose the body is at a distance \( y \) above the equilibrium position as in Fig. 6.10c. The extension of the spring is \( (dl - y) \). The upward force exerted on the body is \( k \, (dl - y) \) and the resultant force \( F \) on the body is

\[
F = k \, (dl - y) - mg = -ky \quad \cdots (2)
\]

The resultant force is proportional to the displacement of the body from its equilibrium position and the motion is simple harmonic.

If the total extension produced is \( (dl + y) \) as in Fig. 6.10d the restoring force on the body is \( k \, (dl + y) \) which acts upwards.
So, the increase in the upward force on the spring is
\[ k (dl + y) - mg = ky \]

Therefore if we produce an extension downward then the restoring force in the spring increases by \( ky \) in the upward direction. As the force acts in the opposite direction to that of displacement, the restoring force is \(-ky\) and the motion is SHM.

\[ F = -ky \]
\[ ma = -ky \]
\[ a = -\frac{k}{m} y \]  \( \text{...(3)} \)
\[ a = -\omega^2 y \]  (expression for SHM)

Comparing the above equations, \( \omega = \sqrt{\frac{k}{m}} \)  \( \text{...(4)} \)

But \( T = \frac{2\pi}{\omega} = 2\pi \sqrt{\frac{m}{k}} \)  \( \text{...(5)} \)

From equation (1) \( mg = k \, dl \)
\[ \frac{m}{k} = \frac{dl}{g} \]

Therefore time period \( T = 2\pi \sqrt{\frac{dl}{g}} \)  \( \text{...(6)} \)

Frequency \( n = \frac{1}{2\pi} \sqrt{\frac{g}{dl}} \)

**Case 1 : When two springs are connected in parallel**

Two springs of spring factors \( k_1 \) and \( k_2 \) are suspended from a rigid support as shown in Fig. 6.11. A load \( m \) is attached to the combination.

Let the load be pulled downwards through a distance \( y \) from its equilibrium position. The increase in length is \( y \) for both the springs but their restoring forces are different.
If $F_1$ and $F_2$ are the restoring forces

\[ F_1 = -k_1y, \quad F_2 = -k_2y \]

\[ \therefore \text{Total restoring force} = (F_1 + F_2) = -(k_1 + k_2)\ y \]

So, time period of the body is given by

\[ T = 2\pi \sqrt{\frac{m}{k_1 + k_2}} \]

If $k_1 = k_2 = k$

Then, \[ T = 2\pi \sqrt{\frac{m}{2k}} \]

\[ \therefore \text{frequency} \ n = \frac{1}{2\pi} \sqrt{\frac{2k}{m}} \]

**Case 2 : When two springs are connected in series.**

Two springs are connected in series in two different ways.

This arrangement is shown in Fig. 6.12a and 6.12b.

In this system when the combination of two springs is displaced to a distance $y$, it produces extension $y_1$ and $y_2$ in two springs of force constants $k_1$ and $k_2$.

\[ F = -k_1\ y_1 \quad : \quad F = -k_2\ y_2 \]
where $F$ is the restoring force.

Total extension, \[ y = y_1 + y_2 = -F \left[ \frac{1}{k_1} + \frac{1}{k_2} \right] \]

We know that $F = -ky$

\[ \therefore y = -\frac{F}{k} \]

From the above equations,

\[ -\frac{F}{k} = -F \left[ \frac{1}{k_1} + \frac{1}{k_2} \right] \]

or \[ k = \frac{k_1 k_2}{k_1 + k_2} \]

\[ \therefore \text{Time period} = T = 2\pi \sqrt{\frac{m(k_1 + k_2)}{k_1 k_2}} \]

frequency \[ n = \frac{1}{2\pi} \sqrt{\frac{k_1 k_2}{(k_1 + k_2) m}} \]

If both the springs have the same spring constant,

\[ k_1 = k_2 = k. \]

\[ \therefore n = \frac{1}{2\pi} \sqrt{\frac{k}{2m}} \]

### 6.5.3 Oscillation of liquid column in a U-tube

Consider a non-viscous liquid column of length $l$ of uniform cross-sectional area $A$ (Fig. 6.13). Initially the level of liquid in the limbs is the same. If the liquid on one side of the tube is depressed by blowing gently the levels of the liquid oscillates for a short time about their initial positions $O$ and $C$, before coming to rest.

If the liquid in one of the limbs is depressed by $y$, there will be a difference of $2y$ in the liquid levels in the two limbs. At some instant, suppose the level of the liquid on the left side of the tube is
at D, at a height y above its original position O, the level B of the liquid on the other side is then at a depth y below its original position C. So the excess pressure P on the liquid due to the restoring force is excess height × density × g

(i.e) pressure = 2 y ρ g

∴ Force on the liquid = pressure × area
of the cross-section of the tube

= – 2 y ρ g × A

...... (1)

The negative sign indicates that the force towards O is opposite to the displacement measured from O at that instant.

The mass of the liquid column of length l is volume × density

(i.e) m = l A ρ

∴ F = l A ρ a

...... (2)

From equations (1) and (2) l A ρ a = – 2 y A ρ g

∴ a = – \frac{2g}{l} y

...... (3)

We know that a = –ω² y

(i.e) a = – \frac{2g}{l} y = –ω² y

where ω = \frac{2g}{\sqrt{l}}

Here, the acceleration is proportional to the displacement, so the motion is simple harmonic and the period T is

T = \frac{2\pi}{\omega} = 2\pi \sqrt{\frac{l}{2g}}

6.5.4 Oscillations of a simple pendulum

A simple pendulum consists of massless and inelastic thread whose one end is fixed to a rigid support and a small bob of mass m is suspended from the other end of the thread. Let l be the length of the pendulum. When the bob is slightly displaced and released, it oscillates about its equilibrium position. Fig.6.14 shows the displaced position of the pendulum.
Suppose the thread makes an angle $\theta$ with the vertical. The distance of the bob from the equilibrium position $A$ is $AB$. At $B$, the weight $mg$ acts vertically downwards. This force is resolved into two components.

(i) The component $mg \cos \theta$ is balanced by the tension in the thread acting along the length towards the fixed point $O$.

(ii) $mg \sin \theta$ which is unbalanced, acts perpendicular to the length of thread. This force tends to restore the bob to the mean position. If the amplitude of oscillation is small, then the path of the bob is a straight line.

\[ F = -mg \sin \theta \]  

...(1)

If the angular displacement is small $\sin \theta \approx \theta$

\[ F = -mg \theta \]  

...(2)

But $\theta = \frac{x}{l}$.

\[ F = -mg \frac{x}{l} \]

Comparing this equation with Newton’s second law, $F = ma$ we get, acceleration $a = -\frac{gx}{l}$  

...(3)

(negative sign indicates that the direction of acceleration is opposite to the displacement) Hence the motion of simple pendulum is SHM.

We know that $a = -\omega^2 x$

Comparing this with (3)

\[ \omega^2 = \frac{g}{l} \]

or $\omega = \sqrt{\frac{g}{l}}$  

...(4)

\[ \therefore \quad \text{Time period } T = \frac{2\pi}{\omega} \]

\[ T = 2\pi \sqrt{\frac{l}{g}} \]  

...(5)

\[ \therefore \quad \text{frequency } n = \frac{1}{2\pi} \sqrt{\frac{g}{l}} \]  

...(6)
Laws of pendulum

From the expression for the time period of oscillations of a pendulum the following laws are enunciated.

(i) The law of length

The period of a simple pendulum varies directly as the square root of the length of the pendulum.

\[ T \propto \sqrt{l} \]

(ii) The law of acceleration

The period of a simple pendulum varies inversely as the square root of the acceleration due to gravity.

\[ T \propto \frac{1}{\sqrt{g}} \]

(iii) The law of mass

The time period of a simple pendulum is independent of the mass and material of the bob.

(iv) The law of amplitude

The period of a simple pendulum is independent of the amplitude provided the amplitude is small.

Note: The length of a seconds pendulum is 0.99 m whose period is two seconds.

\[ 2 = 2\pi \sqrt{\frac{\ell}{g}} \]

\[ \therefore \ell = \frac{9.81 \times 4}{4\pi^2} = 0.99 \text{ m} \]

Oscillations of simple pendulum can also be regarded as a case of angular SHM.

Let \( \theta \) be the angular displacement of the bob B at an instant of time. The bob makes rotation about the horizontal line which is perpendicular to the plane of motion as shown in Fig. 6.15.

Restoring torque about O is \( \tau = -mg \ell \sin \theta \)
\[ \tau = -mg l \theta \quad [\because \theta \text{ is small}] \quad \text{...(1)} \]

Moment of inertia about the axis = \( m l^2 \) \quad \text{...(2)}

If the amplitude is small, motion of the bob is angular simple harmonic. Therefore angular acceleration of the system about the axis of rotation is

\[ \alpha = \frac{\tau}{I} = \frac{-mg l \theta}{m l^2} \]

\[ \alpha = -\frac{g}{l} \theta \quad \text{...(3)} \]

We know that \( \alpha = -\omega^2 \theta \) \quad \text{...(4)}

Comparing (3) and (4)

\[ -\omega^2 \theta = -\frac{g}{l} \theta \]

angular frequency \( \omega = \sqrt{\frac{g}{l}} \)

Time period \( T = \frac{2\pi}{\omega} = 2\pi \sqrt{\frac{l}{g}} \quad \text{...(5)} \)

Frequency \( n = \frac{1}{2\pi} \sqrt{\frac{g}{l}} \quad \text{...(6)} \)

\section*{6.6 Energy in simple harmonic motion}

The total energy \( (E) \) of an oscillating particle is equal to the sum of its kinetic energy and potential energy if conservative force acts on it.

The velocity of a particle executing SHM at a position where its displacement is \( y \) from its mean position is \( v = \omega \sqrt{a^2 - y^2} \)

\textbf{Kinetic energy}

Kinetic energy of the particle of mass \( m \) is

\[ K = \frac{1}{2} m \left( \omega \sqrt{a^2 - y^2} \right)^2 \]

\[ K = \frac{1}{2} m \omega^2 (a^2 - y^2) \quad \text{...(1)} \]
Potential energy

From definition of SHM \( F = -ky \) the work done by the force during the small displacement \( dy \) is \( dW = -F \, dy = -(ky) \, dy = ky \, dy \)

\[ \therefore \text{Total work done for the displacement } y \text{ is,} \]
\[ W = \int_0^y ky \, dy \]
\[ W = \int_0^y m\omega^2 y \, dy \quad \therefore (k = m\omega^2) \]
\[ \therefore W = \frac{1}{2} m \omega^2 y^2 \]

This work done is stored in the body as potential energy
\[ U = \frac{1}{2} m \omega^2 y^2 \quad \ldots(2) \]

Total energy \( E = K + U \)
\[ = \frac{1}{2} m\omega^2 (a^2 - y^2) + \frac{1}{2} m \omega^2 y^2 \]
\[ = \frac{1}{2} m \omega^2 a^2 \]

Thus we find that the total energy of a particle executing simple harmonic motion is \( \frac{1}{2} m \omega^2 a^2 \).

Special cases

(i) When the particle is at the mean position \( y = 0 \), from eqn (1) it is known that kinetic energy is maximum and from eqn. (2) it is known that potential energy is zero. Hence the total energy is wholly kinetic
\[ E = K_{\text{max}} = \frac{1}{2} m\omega^2 a^2 \]

(ii) When the particle is at the extreme position \( y = \pm a \), from eqn. (1) it is known that kinetic energy is zero and from eqn. (2) it is known that Potential energy is maximum. Hence the total energy is wholly potential.
\[ E = U_{\text{max}} = \frac{1}{2} m \omega^2 a^2 \]

(iii) when \( y = \frac{a}{2} \),

\[ K = \frac{1}{2} m \omega^2 \left[ a^2 - \frac{a^2}{4} \right] \]

\[ \therefore K = \frac{3}{4} \left( \frac{1}{2} m \omega^2 a^2 \right) \]

\[ K = \frac{3}{4} E \]

\[ U = \frac{1}{2} m \omega^2 \left( \frac{a}{2} \right)^2 = \frac{1}{4} \left( \frac{1}{2} m \omega^2 a^2 \right) \]

\[ \therefore U = \frac{1}{4} E \]

If the displacement is half of the amplitude, \( K = \frac{3}{4} E \) and \( U = \frac{1}{4} E \). \( K \) and \( U \) are in the ratio 3 : 1.

\[ E = K + U = \frac{1}{2} m \omega^2 a^2 \]

At any other position the energy is partly kinetic and partly potential.

This shows that the particle executing SHM obeys the law of conservation of energy.

**Graphical representation of energy**

The values of \( K \) and \( U \) in terms of \( E \) for different values of \( y \) are given in the Table 6.2. The variation of energy of an oscillating particle with the displacement can be represented in a graph as shown in the Fig. 6.16.
Table 6.2 Energy of SHM

<table>
<thead>
<tr>
<th></th>
<th>y</th>
<th>$a/2$</th>
<th>$a$</th>
<th>$-a/2$</th>
<th>$-a$</th>
</tr>
</thead>
<tbody>
<tr>
<td>Kinetic energy</td>
<td>E</td>
<td>$3/4$E</td>
<td>0</td>
<td>$3/4$E</td>
<td>0</td>
</tr>
<tr>
<td>Potential energy</td>
<td>0</td>
<td>$1/4$E</td>
<td>E</td>
<td>$1/4$E</td>
<td>E</td>
</tr>
</tbody>
</table>

6.7 Types of oscillations

There are three main types of oscillations.

(i) Free oscillations

When a body vibrates with its own natural frequency, it is said to execute free oscillations. The frequency of oscillations depends on the inertial factor and spring factor, which is given by,

$$n = \frac{1}{2\pi} \sqrt{\frac{k}{m}}$$

Examples

(i) Vibrations of tuning fork
(ii) Vibrations in a stretched string
(iii) Oscillations of simple pendulum
(iv) Air blown gently across the mouth of a bottle.

(ii) Damped oscillations

Most of the oscillations in air or in any medium are damped. When an oscillation occurs, some kind of damping force may arise due to friction or air resistance offered by the medium. So, a part of the energy is dissipated in overcoming the resistive force. Consequently, the amplitude of oscillation decreases with time and finally becomes zero. Such oscillations are called damped oscillations (Fig. 6.17).
Examples:

(i) The oscillations of a pendulum
(ii) Electromagnetic damping in galvanometer (oscillations of a coil in galvanometer)
(iii) Electromagnetic oscillations in tank circuit

(iii) Maintained oscillations

The amplitude of an oscillating system can be made constant by feeding some energy to the system. If an energy is fed to the system to compensate the energy it has lost, the amplitude will be a constant. Such oscillations are called maintained oscillations (Fig. 6.18).

Example:

A swing to which energy is fed continuously to maintain amplitude of oscillation.

(iv) Forced oscillations

When a vibrating body is maintained in the state of vibration by a periodic force of frequency \(n\) other than its natural frequency of the body, the vibrations are called forced vibrations. The external force is driver and body is driven.

The body is forced to vibrate with an external periodic force. The amplitude of forced vibration is determined by the difference between the frequencies of the driver and the driven. The larger the frequency difference, smaller will be the amplitude of the forced oscillations.

Examples:

(i) Sound boards of stringed instruments execute forced vibration,
(ii) Press the stem of vibrating tuning fork, against tabla. The tabla suffers forced vibration.

(v) Resonance

In the case of forced vibration, if the frequency difference is small,
the amplitude will be large (Fig. 6.19). Ultimately when the two frequencies are same, amplitude becomes maximum. This is a special case of forced vibration.

If the frequency of the external periodic force is equal to the natural frequency of oscillation of the system, then the amplitude of oscillation will be large and this is known as resonance.

**Advantages**

(i) Using resonance, frequency of a given tuning fork is determined with a sonometer.

(ii) In radio and television, using tank circuit, required frequency can be obtained.

**Disadvantages**

(i) Resonance can cause disaster in an earthquake, if the natural frequency of the building matches the frequency of the periodic oscillations present in the Earth. The building begins to oscillate with large amplitude thus leading to a collapse.

(ii) A singer maintaining a note at a resonant frequency of a glass, can cause it to shatter into pieces.
Solved problems

6.1 Obtain an equation for the SHM of a particle whose amplitude is 0.05 m and frequency 25 Hz. The initial phase is $\pi/3$.

**Data** : $a = 0.05$ m, $n = 25$ Hz, $\phi_0 = \pi/3$.

**Solution** : $\omega = 2\pi n = 2\pi \times 25 = 50\pi$

The equation of SHM is $y = a \sin (\omega t + \phi_0)$

The displacement equation of SHM is : $y = 0.05 \sin (50\pi t + \pi/3)$

6.2 The equation of a particle executing SHM is $y = 5 \sin \left(\frac{\pi t}{3}\right)$.

Calculate (i) amplitude (ii) period (iii) maximum velocity and (iv) velocity after 1 second (y is in metre).

**Data** : $y = 5 \sin \left(\frac{\pi t}{3}\right)$

**Solution** : The equation of SHM is $y = a \sin (\omega t + \phi_0)$

Comparing the equations

(i) Amplitude $a = 5$ m

(ii) Period, $T = \frac{2\pi}{\omega} = \frac{2\pi}{\pi} = 2$ s

(iii) $v_{\text{max}} = a\omega = 5 \times \pi = 15.7$ m s$^{-1}$

(iv) Velocity after 1 s = $a\omega \cos (\omega t + \phi_0)$

$= 15.7 \cos \left(\frac{\pi t \times 1 + \pi}{3}\right)$

$= 15.7 \times \frac{1}{2} = 7.85$ m s$^{-1}$

$\therefore v = 7.85$ m s$^{-1}$

6.3 A particle executes a simple harmonic motion of time period T. Find the time taken by the particle to have a displacement from mean position equal to one half of the amplitude.
Solution: The displacement is given by \( y = a \sin \omega t \)

When the displacement \( y = \frac{a}{2} \),

we get \( \frac{a}{2} = a \sin \omega t \)

or \( \sin \omega t = \frac{1}{2} \)

\( \omega t = \frac{\pi}{6} \)

\( t = \frac{\pi}{6\omega} = \frac{\pi}{6 \cdot \frac{2\pi}{T}} \)

The time taken is \( t = \frac{T}{12} \) s

6.4 The velocities of a particle executing SHM are 4 cm s\(^{-1}\) and 3 cm s\(^{-1}\), when its distance from the mean position is 2 cm and 3 cm respectively. Calculate its amplitude and time period.

Data: \( v_1 = 4 \text{ cm s}^{-1} = 4 \times 10^{-2} \text{ m s}^{-1} \); \( v_2 = 3 \text{ cm s}^{-1} = 3 \times 10^{-2} \text{ m s}^{-1} \);
\( y_1 = 2 \text{ cm} = 2 \times 10^{-2} \text{ m} \); \( y_2 = 3 \text{ cm} = 3 \times 10^{-2} \text{ m} \)

Solution:

\( v_1 = \omega \sqrt{a^2 - y_1^2} \) \( \cdots (1) \)
\( v_2 = \omega \sqrt{a^2 - y_2^2} \) \( \cdots (2) \)

Squaring and dividing the equations

\[
\left( \frac{v_1}{v_2} \right)^2 = \frac{a^2 - y_1^2}{a^2 - y_2^2}
\]

\[
\left( \frac{4 \times 10^{-2}}{3 \times 10^{-2}} \right)^2 = \frac{a^2 - 4 \times 10^{-4}}{a^2 - 9 \times 10^{-4}}
\]

\[
9a^2 - 36 \times 10^{-4} = 16a^2 - 144 \times 10^{-4}
\]

\[
7a^2 = 108 \times 10^{-4}
\]

\( \therefore a = \sqrt{15.42 \times 10^{-2}} = 0.03928 \text{ m} \)
Substituting the value of $a^2$ in equation (1) we have

$$4 \times 10^{-2} = \omega \sqrt{\frac{108 \times 10^4}{7} - 4 \times 10^4}$$

$$\therefore \omega = \frac{\sqrt{7}}{5} \text{ rad s}^{-1}$$

$$\therefore \text{Time period } T = \frac{2\pi}{\omega} = 2\pi \sqrt{\frac{5}{7}}$$

$$T = 5.31 \text{ s}$$

6.5 A circular disc of mass 10 kg is suspended by a wire attached to its centre. The wire is twisted by rotating the disc and released. The period of torsional oscillation is found to be 1.5 s. The radius of the disc is 15 cm. Calculate the torsional spring constant.

**Data**: $m = 10 \text{ kg}$, $T = 1.5 \text{ s}$, $r = 15 \text{ cm} = 15 \times 10^{-2} \text{ m}$ \quad $C = ?$

**Solution**: MI of the disc about an axis through the centre is

$$I = \frac{1}{2} MR^2$$

The time period of angular SHM is

$$T = 2\pi \sqrt{\frac{I}{C}}$$

Squaring the equation, $T^2 = 4\pi^2 \frac{I}{C}$

$$\therefore C = \frac{4\pi^2 I}{T^2}$$

$$C = \frac{4\pi^2 \times \frac{1}{2} MR^2}{T^2}$$

$$= \frac{2 \times (3.14)^2 \times 10 \times 0.15^2}{(1.5)^2}$$

$$C = 2.0 \text{ N m rad}^{-1}$$

6.6 A body of mass 2 kg executing SHM has a displacement $y = 3 \sin \left(100t + \frac{\pi}{4}\right) \text{ cm}$. Calculate the maximum kinetic energy of the body.
Solution: Comparing with equation of SHM

\[ y = a \sin (\omega t + \phi) \]
\[ a = 3 \text{ cm} = 3 \times 10^{-2} \text{ m}, \ \omega = 100 \text{ rad s}^{-1}, \ m = 2 \text{ kg} \]

\[ y = 3 \sin \left(100 t + \frac{\pi}{4}\right) \]

Maximum kinetic energy = \( \frac{1}{2} m a^2 \omega^2 \)
\[ = \frac{1}{2} \times 2 \times (0.03^2 \times 100^2) \]
Maximum kinetic energy = 9 joule

6.7 A block of mass 15 kg executes SHM under the restoring force of a spring. The amplitude and the time period of the motion are 0.1 m and 3.14 s respectively. Find the maximum force exerted by the spring on the block.

Data: \( m = 15 \text{ kg}, \ a = 0.1 \text{ m and } T = 3.14 \text{ s} \)

Solution: The maximum force exerted on the block is \( ka \), when the block is at the extreme position, where \( k \) is the spring constant.

The angular frequency \( \omega = \frac{2\pi}{T} = 2 \text{ s}^{-1} \)

The spring constant \( k = m \omega^2 \)
\[ = 15 \times 4 = 60 \text{ N m}^{-1} \]

The maximum force exerted on the block is \( ka = 60 \times 0.1 = 6 \text{ N} \)

6.8 A block of mass 680 g is attached to a horizontal spring whose spring constant is 65 N m\(^{-1}\). The block is pulled to a distance of 11 cm from the mean position and released from rest. Calculate: (i) angular frequency, frequency and time period (ii) displacement of the system (iii) maximum speed and acceleration of the system

Data: \( m = 680 \text{ g} = 0.68 \text{ kg}, \ k = 65 \text{ N m}^{-1}, \ a = 11 \text{ cm} = 0.11 \text{ m} \)

Solution: The angular frequency \( \omega = \sqrt{\frac{k}{m}} \)
\[ \omega = \sqrt{\frac{65}{0.68}} = 9.78 \text{ rad s}^{-1} \]
The frequency $n = \frac{\omega}{2\pi} = \frac{9.78}{2\pi} = 1.56 \text{ Hz}$

The time period $T = \frac{1}{n} = \frac{1}{1.56} = 0.64 \text{ s}$

maximum speed $= a \omega$

$= 0.11 \times 9.78$

$= 1.075 \text{ m s}^{-1}$

Acceleration of the block $= a \omega^2 = a\omega \times \omega$

$= 1.075 \times (9.78)$

$= 10.52 \text{ m s}^{-2}$

Displacement $y(t) = a \sin \omega t$

$\therefore y(t) = 0.11 \sin 9.78 t \text{ metre}$

6.9 A mass of 10 kg is suspended by a spring of length 60 cm and force constant $4 \times 10^3 \text{ N m}^{-1}$. If it is set into vertical oscillations, calculate the (i) frequency of oscillation of the spring and (ii) the length of the stretched string.

**Data** : $k = 4 \times 10^3 \text{ N m}^{-1}$, $F = 10 \times 9.8 \text{ N}$, $l = 60 \times 10^{-2} \text{ m}$, $m = 10 \text{ kg}$

**Solution** : (i) $n = \frac{1}{2\pi} \sqrt{\frac{k}{m}}$

$$= \frac{1}{2\pi} \sqrt{\frac{4 \times 10^3}{10}} = \frac{20}{2\pi}$$

Frequency $= 3.184 \text{ Hz}$

(ii) $T = 2\pi \sqrt{\frac{dl}{g}}$ or $T^2 = 4\pi^2 \frac{dl}{g}$

length $(dl) = \frac{T^2 g}{4\pi^2} = \frac{1}{n^2} \times \frac{g}{4\pi^2}$

$\therefore dl = \frac{9.8}{(3.184)^2 \times 4 \times (3.14)^2}$

$dl = 0.0245 \text{ m}$

$\therefore$ The length of the stretched string $= 0.6 + 0.0245 = 0.6245 \text{ m}$
6.10 A mass $m$ attached to a spring oscillates every 4 seconds. If the mass is increased by 4 kg, the period increases by 1 s. Find its initial mass $m$.

**Data:** Mass $m$ oscillates with a period of 4 s

When the mass is increased by 4 kg period is 5 s

**Solution:** Period of oscillation $T = 2\pi \sqrt{\frac{m}{k}}$

$$4 = 2\pi \sqrt{\frac{m}{k}} \quad \ldots \ (1)$$

$$5 = 2\pi \sqrt{\frac{m+4}{k}} \quad \ldots \ (2)$$

Squaring and dividing the equations

$$\frac{25}{16} = \frac{m+4}{m}$$

$$25m = 16m + 64$$

$$9m = 64$$

$$m = \frac{64}{9} = 7.1 \text{ kg}$$

6.11 The acceleration due to gravity on the surface of moon is 1.7 m s$^{-2}$. What is the time period of a simple pendulum on the surface of the moon, if its period on the Earth is 3.5 s?

**Data:**

$g$ on moon = 1.7 m s$^{-2}$

$g$ on the Earth = 9.8 ms$^{-2}$

Time period on the Earth = 3.5 s

**Solution:**

$$T = 2\pi \sqrt{\frac{l}{g}}$$

Let $T_m$ represent the time period on moon

$$T_m = 2\pi \sqrt{\frac{l}{1.7}} \quad \ldots \ (1)$$

On the Earth,

$$3.5 = 2\pi \sqrt{\frac{l}{9.8}} \quad \ldots \ (2)$$

Dividing the equation (2) by (1) and squaring
\[
\left( \frac{3.5}{T_m} \right)^2 = \frac{1.7}{9.8}
\]
\[
T_m^2 \times 1.7 = (3.5)^2 \times 9.8
\]
\[
T_m^2 = \frac{(3.5)^2 \times 9.8}{1.7} = \frac{12.25 \times 9.8}{1.7}
\]
\[
\therefore T_m = \sqrt{\frac{120.05}{1.7}} = 8.40 \text{ s}
\]

6.12 A simple pendulum has a period 4.2 s. When the pendulum is shortened by 1 m the period is 3.7 s. Calculate its (i) acceleration due to gravity (ii) original length of the pendulum.

**Data:** \( T = 4.2 \text{ s} \); when length is shortened by 1 m the period is 3.7 s.

**Solution:**

\[
T = 2\pi \sqrt{\frac{l}{g}}
\]

**Squaring and rearranging**

\[
g = 4\pi^2 \frac{l}{T^2}
\]

\[
\therefore g = 4\pi^2 \frac{l}{(4.2)^2} \quad \text{...(1)}
\]

When the length is shortened by 1 m

\[
g = \frac{4\pi^2(l-1)}{(3.7)^2} \quad \text{...(2)}
\]

From the above equations

\[
\frac{l}{(4.2)^2} = \frac{l-1}{(3.7)^2}
\]

\[
(7.9 \times 0.5) l = 17.64
\]

\[
l = \frac{17.64}{7.9 \times 0.5} = 4.46 \text{ m}
\]

Substituting in equation (1)

\[
g = \frac{4.46}{(4.2)^2} \times \frac{175.89}{17.64}
\]

\[
g = 9.97 \text{ m s}^{-2}
\]
Self evaluation

(The questions and problems given in this self evaluation are only samples. In the same way any question and problem could be framed from the text matter. Students must be prepared to answer any question and problem from the text matter, not only from the self evaluation.)

6.1 Which of the following is the necessary condition for SHM?
(a) constant period
(b) constant acceleration
(c) displacement and acceleration are proportional
(d) displacement and torque are proportional

6.2 The displacement of a particle executing SHM is given by
\[ x = 0.01 \sin (100 \pi t + 0.05) \]. Its time period is
(a) 0.01 s  (b) 0.02 s
(c) 0.1 s  (d) 0.2 s

6.3 If the displacement of a particle executing SHM is given by
\[ y = 0.05 \sin (100 t + \frac{\pi}{2}) \text{ cm}. \] The maximum velocity of the particle is
(a) 0.5 cm s\(^{-1}\)  (b) 0.05 m s\(^{-1}\)
(c) 100 m s\(^{-1}\)  (d) 50 m s\(^{-1}\)

6.4 If the magnitude of displacement is equal to acceleration, then the time period is.
(a) 1 s  (b) \(\pi\) s
(c) 2\(\pi\) s  (d) 4\(\pi\) s

6.5 A body of mass 2 g is executing SHM about a mean position with an amplitude 10 cm. If the maximum velocity is 100 cm s\(^{-1}\) its velocity is 50 cm s\(^{-1}\) at a distance of (in cm).
(a) \(5\sqrt{2}\)  (b) \(50\sqrt{3}\)
(c) \(5\sqrt{3}\)  (d) \(10\sqrt{3}\)

6.6 A linear harmonic oscillator has a total energy of 160 J. Its
(a) maximum potential energy is 100 J
(b) maximum kinetic energy is 160 J
(c) minimum potential energy is 100 J
(d) maximum kinetic energy is 100 J
6.7 A force of 6.4 N stretches a vertical spring by 0.1 m. The mass that must be suspended from the spring so that it oscillates with a period of \( \frac{\pi}{4} \) s is

(a) \( \frac{\pi}{4} \) kg \quad (b) 1 kg \quad (c) \( \frac{1}{4} \) kg \quad (d) 10 kg

6.8 The length of seconds pendulum at a place where \( g = 9.8 \text{ m s}^{-2} \) is

(a) 0.25 m \quad (b) 1 m \quad (c) 0.99 m \quad (d) 0.50 m

6.9 A particle executes SHM with an amplitude 4 cm. At what displacement from the mean position its energy is half kinetic and half potential?

(a) \( 2\sqrt{2} \) cm \quad (b) \( \sqrt{2} \) cm \quad (c) 2 cm \quad (d) 1 cm

6.10 A particle executes SHM along a straight line with an amplitude 'a'. PE is maximum when the displacement is

(a) \( \pm a \) \quad (b) zero \quad (c) \( \frac{a}{2} \) \quad (d) \( \frac{a}{\sqrt{2}} \)

6.11 Define simple harmonic motion. What are the conditions of SHM?

6.12 Every SHM is periodic motion but every periodic motion need not be SHM. Why? Support your answer with an example.

6.13 Show that the projection of uniform circular motion on the diameter of a circle is simple harmonic motion.

6.14 Explain: (i) displacement (ii) velocity and (iii) acceleration in SHM using component method.

6.15 Show graphically the variation of displacement, velocity and acceleration of a particle executing SHM.

6.16 What is the phase difference between (i) velocity and acceleration (ii) acceleration and displacement of a particle executing SHM?
6.17 Derive the differential formula for SHM.
6.18 Define the terms (i) time period (ii) frequency and (iii) angular frequency.
6.19 Define force constant. Give its unit and dimensional formula.
6.20 What is phase of SHM? Explain the term phase difference.
6.21 Derive an expression for the time period of a body when it executes angular SHM.
6.22 What is an epoch? Give its unit.
6.23 Explain the oscillations of a mass attached to a horizontal spring. Hence deduce an expression for its time period.
6.24 Obtain an expression for the frequency of vertical oscillations of a loaded spring.
6.25 Distinguish between linear and angular harmonic oscillator?
6.26 What is a spring factor?
6.27 Show that the oscillations of a simple pendulum are simple harmonic. Hence deduce the expression for the time period.
6.28 The bob of a simple pendulum is a hollow sphere filled with water. How does the period of oscillation change if the water begins to drain out of the sphere?
6.29 Why does the oscillation of a simple pendulum eventually stop?
6.30 What will happen to the time period of a simple pendulum if its length is doubled?
6.31 Derive an expression for the total energy of a particle executing SHM.
6.32 On what factors the natural frequency of a body depend on?
6.33 What is forced vibration? Give an example.
6.34 What forces keep the simple pendulum in SHM?
6.35 Illustrate an example to show that resonance is disastrous sometimes.
6.36 If two springs are connected in parallel, what is its equivalent spring constant?
6.37 If two springs are connected in series, what is its equivalent spring constant?
Problems

6.38 Obtain an equation for the SHM of a particle of amplitude 0.5 m, frequency 50 Hz. The initial phase is \( \frac{\pi}{2} \). Find the displacement at \( t = 0 \).

6.39 The equation of SHM is represented by \( y = 0.25 \sin (3014 t + 0.35) \), where \( y \) and \( t \) are in mm and s respectively. Deduce (i) amplitude (ii) frequency (iii) angular frequency (iv) period and (v) initial phase.

6.40 A particle executing SHM is represented by \( y = 2 \sin \left( \frac{2\pi}{T} t + \phi_0 \right) \).
At \( t = 0 \), the displacement is \( \sqrt{3} \) cm. Find the initial phase.

6.41 A particle executing SHM has angular frequency of \( \pi \) rad s\(^{-1} \) and amplitude of 5 m. Deduce (i) time period (ii) maximum velocity (iii) maximum acceleration (iv) velocity when the displacement is 3 m.

6.42 A body executes SHM with an amplitude 10 cm and period 2 s. Calculate the velocity and acceleration of the body when the displacement is (i) zero and (ii) 6 cm.

6.43 A disc suspended by a wire makes angular oscillations. When it is displaced through 30° from the mean position, it produces a restoring torque of 4.6 N m. If the moment of inertia of the disc is 0.082 kg m\(^2\), calculate the frequency of angular oscillations.

6.44 A spring of force constant 1200 N m\(^{-1}\) is mounted on a horizontal table as shown in figure. A mass of 3 kg is attached to its free end and pulled side ways to a distance of 2 cm and released. Calculate (i) the frequency of oscillation (ii) the maximum velocity and (iii) maximum acceleration of the mass.

6.45 A mass of 0.2 kg attached to one end of a spring produces an extension of 15 mm. The mass is pulled 10 mm downwards and set into vertical oscillations of amplitude 10 mm. Calculate (i) the period of oscillation (ii) maximum kinetic energy.
6.46  A 5 kg mass is suspended by a system of two identical springs of spring constant 250 N m\(^{-1}\) as shown in figure. Determine the period of oscillation the system.

![Diagram of two identical springs suspending a mass](image)

6.47  A trolley of mass 2 kg is connected between two identical springs of spring constant 400 N m\(^{-1}\). If the trolley is displaced from its mean position by 3 cm and released, calculate its (i) time period (ii) maximum velocity (iii) maximum kinetic energy.

![Diagram of a trolley connected by two springs](image)

6.48  A vertical U tube of uniform cross section contains water to a height of 0.3 m. Show that, if water in one of the limbs is depressed and then released, the oscillations of the water column in the tube are SHM. Calculate its time period also.

6.49  A bob of a simple pendulum oscillates with an amplitude of 4 cm and time period 1 s. Find (i) length of the pendulum and (ii) velocity of the bob in the mean position.

6.50  Compare the acceleration due to gravity at two places if the time for 100 oscillations of a simple pendulum are 8 minutes 2 seconds and 8 minutes 20 seconds respectively of the two places.

6.51  A particle of mass 0.2 kg executes SHM of amplitude 2 cm and time period 6 s. Calculate (i) the total energy (ii) kinetic and potential energy when displacement is 1 cm from the mean position.

6.52  The length of a seconds pendulum in a clock is increased by 2%. How many seconds will it lose or gain in a day?
Answers

6.1 (c) 6.2 (b) 6.3 (b) 6.4 (c)
6.5 (c) 6.6 (b) 6.7 (b) 6.8 (c)
6.9 (a) 6.10 (a)

6.38 0.5 m
6.39 $0.25 \times 10^{-3}$ m, 480 Hz, 3014 rad s$^{-1}$, 0.0021 s, 0.35 rad
6.40 60°
6.41 2 s, 15.7 m s$^{-1}$, 49.3 m s$^{-2}$, 12.56 m s$^{-1}$
6.42 0.314 m s$^{-1}$, zero; 0.2512 m s$^{-1}$, 0.5915 m s$^{-2}$
6.43 1.64 Hz
6.44 3.2 Hz, 0.40 m s$^{-1}$, 8.07 m s$^{-2}$
6.45 0.25 s, 6.533 $\times 10^{-3}$ J
6.46 0.628 s
6.47 0.314 s, 0.6 m s$^{-1}$, 0.36 J
6.48 1.0098 s
6.49 0.25 m, 0.2512 m s$^{-1}$
6.50 1.076
6.51 $4.386 \times 10^{-5}$ J, $3.286 \times 10^{-5}$ J, $1.1 \times 10^{-5}$ J
6.52 loss of time is 864 s
7. Wave Motion

Wave motion is a mode of transmission of energy through a medium in the form of a disturbance. It is due to the repeated periodic motion of the particles of the medium about an equilibrium position transferring the energy from one particle to another.

The waves are of three types - mechanical, electromagnetic and matter waves. Mechanical waves can be produced only in media which possess elasticity and inertia. Water waves, sound waves and seismic waves are common examples of this type. Electromagnetic waves do not require any material medium for propagation. Radio waves, microwaves, infrared rays, visible light, the ultraviolet rays, X rays and γ rays are electromagnetic waves. The waves associated with particles like electrons, protons and fundamental particles in motion are matter waves.

Waves on surface of water

In order to understand the concept of wave motion, let us drop a stone in a trough of water. We find that small circular waves seem to originate from the point where the stone touches the surface of water. These waves spread out in all directions. It appears as if water moves away from that point. If a piece of paper is placed on the water surface, it will be observed that the piece of paper moves up and down, when the waves pass through it. This shows that the waves are formed due to the vibratory motion of the water particles, about their mean position.

Wave motion is a form of disturbance which travels through a medium due to the repeated periodic motion of the particles of the medium about their mean position. The motion is transferred continuously from one particle to its neighbouring particle.

7.1 Characteristics of wave motion

(i) Wave motion is a form of disturbance travelling in the medium due to the periodic motion of the particles about their mean position.
(ii) It is necessary that the medium should possess elasticity and inertia.

(iii) All the particles of the medium do not receive the disturbance at the same instant (i.e.) each particle begins to vibrate a little later than its predecessor.

(iv) The wave velocity is different from the particle velocity. The velocity of a wave is constant for a given medium, whereas the velocity of the particles goes on changing and it becomes maximum in their mean position and zero in their extreme positions.

(v) During the propagation of wave motion, there is transfer of energy from one particle to another without any actual transfer of the particles of the medium.

(vi) The waves undergo reflection, refraction, diffraction and interference.

7.1.1 Mechanical wave motion

The two types of mechanical wave motion are (i) transverse wave motion and (ii) longitudinal wave motion

(i) Transverse wave motion

Transverse wave motion is that wave motion in which particles of the medium execute SHM about their mean positions in a direction perpendicular to the direction of propagation of the wave. Such waves are called transverse waves. Examples of transverse waves are waves produced by plucked strings of veena, sitar or violin and electromagnetic waves. Transverse waves travel in the form of crests and troughs. The maximum displacement of the particle in the positive direction i.e. above its mean position is called crest and maximum displacement of the particle in the negative direction i.e. below its mean position is called trough.

Thus if ABCDEFG is a transverse wave, the points B and F are crests while D is trough (Fig. 7.1).

For the propagation of transverse waves, the medium must possess force of cohesion and volume elasticity. Since gases and liquids do not have rigidity (cohesion), transverse waves
cannot be produced in gases and liquids. Transverse waves can be produced in solids and surfaces of liquids only.

(ii) Longitudinal wave motion

‘Longitudinal wave motion is that wave motion in which each particle of the medium executes simple harmonic motion about its mean position along the direction of propagation of the wave.’

Sound waves in fluids (liquids and gases) are examples of longitudinal wave. When a longitudinal wave travels through a medium, it produces compressions and rarefactions.

In the case of a spiral spring, whose one end is tied to a hook of a wall and the other end is moved forward and backward, the coils of the spring vibrate about their original position along the length of the spring and longitudinal waves propagate through the spring (Fig.7.2).

The regions where the coils are closer are said to be in the state of compression, while the regions where the coils are farther are said to be in the state of rarefaction.

When we strike a tuning fork on a rubber pad, the prongs of the tuning fork begin to vibrate to and fro about their mean positions. When the prong A moves outwards to A₁, it compresses the layer of air in its neighbourhood. As the compressed layer moves forward it compresses the next layer and a wave of compression passes through air. But when the prong moves inwards to A₂, the particles of the medium which moved to the right, now move backward to the left due to elasticity of air. This gives rise to rarefaction.

Thus a longitudinal wave is characterised by the formation of compressions and rarefactions following each other.

Longitudinal waves can be produced in all types of material medium, solids, liquids and gases. The density and pressure of the
medium in the region of compression are more than that in the region of rarefaction.

**7.1.2 Important terms used in wave motion**

(i) **Wavelength (λ)**

The distance travelled by a wave during which a particle of the medium completes one vibration is called wavelength. It is also defined as the distance between any two nearest particles on the wave having same phase.

Wavelength may also be defined as the distance between two successive crests or troughs in transverse waves, or the distance between two successive compressions or rarefactions in longitudinal waves.

(ii) **Time period (T)**

The time period of a wave is the time taken by the wave to travel a distance equal to its wavelength.

(iii) **Frequency (n)**

This is defined as the number of waves produced in one second. If \( T \) represents the time required by a particle to complete one vibration, then it makes \( \frac{1}{T} \) waves in one second.

Therefore frequency is the reciprocal of the time period (i.e) \( n = \frac{1}{T} \).

**Relationship between velocity, frequency and wavelength of a wave**

The distance travelled by a wave in a medium in one second is called the velocity of propagation of the wave in that medium. If \( v \) represents the velocity of propagation of the wave, it is given by

\[
\text{Velocity} = \frac{\text{Distance travelled}}{\text{Time taken}}
\]

\[
v = \frac{\lambda}{T} = n\lambda \quad \left[ \therefore n = \frac{1}{T} \right]
\]

The velocity of a wave \( (v) \) is given by the product of the frequency and wavelength.
7.2 Velocity of wave in different media

The velocity of mechanical wave depends on elasticity and inertia of the medium.

7.2.1 Velocity of a transverse wave along a stretched string

Let us consider a string fixed at one of its ends and tension be applied at the other end. When the string is plucked at a point, it begins to vibrate.

Consider a transverse wave proceeding from left to right in the form of a pulse when the string is plucked at a point as shown in Fig. 7.4. EF is the displaced position of the string at an instant of time. It forms an arc of a circle with O as centre and R as radius. The arc EF subtends an angle $2\theta$ at O.

If $m$ is the mass per unit length of the string and $dx$ is the length of the arc EF, then the mass of the portion of the string is $m\,dx$.

\[ \therefore \text{Centripetal force} = \frac{m\,dx\,v^2}{R} \tag{1} \]

This force is along CO. To find the resultant of the tension T at the points E and F, we resolve T into two components $T\cos\theta$ and $T\sin\theta$.

$T\cos\theta$ components acting perpendicular to CO are of equal in magnitude but opposite in direction, they cancel each other.

$T\sin\theta$ components act parallel to CO. Therefore the resultant of the tensions acting at E and F is $2T\sin\theta$. It is directed along CO. If $\theta$ is small, $\sin\theta = \theta$ and the resultant force due to tension is $2T\theta$.

resultant force = $2T\theta$
\[
\frac{dR}{dx} = 2T \cdot \frac{dx}{2R} \quad \left( \therefore 2\rho = \frac{dx}{R} \right)
\]

\[
= T \cdot \frac{dx}{R}
\]

For the arc EF to be in equilibrium,

\[
\frac{m dx v^2}{R} = \frac{T dx}{R}
\]

\[
v^2 = \frac{T}{m}
\]

or

\[
v = \frac{T}{\sqrt{m}}
\] ... (3)

### 7.2.2 Velocity of longitudinal waves in an elastic medium

Velocity of longitudinal waves in an elastic medium is

\[
v = \sqrt{\frac{E}{\rho}}
\] ... (1)

where \(E\) is the modulus of elasticity, \(\rho\) is the density of the medium.

(i) In the case of a solid rod

\[
v = \sqrt{\frac{q}{\rho}}
\] ... (2)

where \(q\) is the Young’s modulus of the material of the rod and \(\rho\) is the density of the rod.

(ii) In liquids, \(v = \sqrt{\frac{k}{\rho}}\)

... (3)

where \(k\) is the Bulk modulus and \(\rho\) is the density of the liquid.

### 7.2.3 Newton’s formula for the velocity of sound waves in air

Newton assumed that sound waves travel through air under isothermal conditions (i.e) temperature of the medium remains constant.

The change in pressure and volume obeys Boyle’s law.

\[
\therefore PV = \text{constant}
\]

Differentiating, \(P \cdot dV + V \cdot dP = 0\)

\[
P \cdot dV = -V \cdot dP
\]
\[ P = \frac{-\text{change in pressure}}{\text{volume strain}} \]

\[ P = k \text{ (Volume Elasticity)} \]

Therefore under isothermal condition, \( P = k \)

\[ v = \sqrt{\frac{k}{\rho}} = \sqrt{\frac{P}{\rho}} \]

where \( P \) is the pressure of air and \( \rho \) is the density of air. The above equation is known as Newton’s formula for the velocity of sound waves in a gas.

At NTP, \( P = 76 \) cm of mercury

\[ = (0.76 \times 13.6 \times 10^3 \times 9.8) \text{ N m}^{-2} \]

\[ \rho = 1.293 \text{ kg m}^{-3}. \]

\[ \therefore \text{Velocity of sound in air at NTP is} \]

\[ v = \sqrt{\frac{0.76 \times 13.6 \times 10^3 \times 9.8}{1.293}} = 280 \text{ m s}^{-1} \]

The experimental value for the velocity of sound in air is 332 m s\(^{-1}\). But the theoretical value of 280 m s\(^{-1}\) is 15% less than the experimental value. This discrepancy could not be explained by Newton’s formula.

7.2.4 Laplace’s correction

The above discrepancy between the observed and calculated values was explained by Laplace in 1816. Sound travels in air as a longitudinal wave. The wave motion is therefore, accompanied by compressions and rarefactions. At compressions the temperature of air rises and at rarefactions, due to expansion, the temperature decreases.

Air is a very poor conductor of heat. Hence at a compression, air cannot lose heat due to radiation and conduction. At a rarefaction it cannot gain heat, during the small interval of time. As a result, the temperature throughout the medium does not remain constant.

Laplace suggested that sound waves travel in air under adiabatic condition and not under isothermal condition.
For an adiabatic change, the relation between pressure and volume is given by

\[ P \, V^\gamma = \text{constant} \]

where \( \gamma = \left( \frac{C_p}{C_v} \right) \) is the ratio of two specific heat capacities of the gas.

Differentiating

\[ P \gamma \, V^{\gamma-1} \, dV + V^\gamma \, dP = 0 \]

\[ P \gamma = \frac{-V \, dP}{dV} \]

\[ P \gamma = \frac{dP}{dV} = k \]

\[ \therefore P \gamma = k \text{ (Volume elasticity)} \]

Therefore under adiabatic condition

velocity of sound \( v = \sqrt{\frac{k}{\rho}} = \sqrt{\frac{\gamma P}{\rho}} \)

This is Laplace’s corrected formula.

For air at NTP
\( \gamma = 1.41, \rho = 1.293 \text{ kg m}^{-3} \)

\[ \therefore v = \sqrt{\frac{\gamma P}{\rho}} = \sqrt{1.41 \times 280} = 331.3 \text{ ms}^{-1} \]

This result agrees with the experimental value of 332 ms\(^{-1}\).

7.2.5 Factors affecting velocity of sound in gases

(i) Effect of pressure

If the temperature of the gas remains constant, then by Boyle’s law \( PV = \text{constant} \)

i.e \( P \cdot \frac{m}{\rho} = \text{constant} \)

\( \frac{P}{\rho} \) is a constant, when mass \( m \) of a gas is constant. If the pressure changes from \( P \) to \( P' \) then the corresponding density also will change from \( \rho \) to \( \rho' \) such that \( \frac{P}{\rho} \) is a constant.
In Laplace’s formula \( \sqrt{\frac{\gamma P}{\rho}} \) is also a constant. Therefore the **velocity of sound in a gas is independent of the change in pressure provided the temperature remains constant.**

**(ii) Effect of temperature**

For a gas, \( PV = RT \)
\[
\frac{m}{\rho} = RT
\]
or
\[
\frac{P}{\rho} = \frac{RT}{m}
\]

where \( m \) is the mass of the gas, \( T \) is the absolute temperature and \( R \) is the gas constant.

Therefore \( v = \sqrt{\frac{\gamma RT}{m}} \)

It is clear that the velocity of sound in a gas is directly proportional to the square root of its absolute temperature.

Let \( v_0 \) and \( v_t \) be the velocity of sound at \( 0^\circ C \) and \( t^\circ C \) respectively. Then, from the above equation,
\[
v_0 = \sqrt{\frac{\gamma R}{m} \times \sqrt{273}}
\]
\[
v_t = \sqrt{\frac{\gamma R}{m} \times \sqrt{273+t}}
\]
\[
\therefore \frac{v_t}{v_0} = \sqrt{\frac{273+t}{273}}
\]
\[
\therefore v_t = v_0 \left( 1 + \frac{t}{273} \right)^{1/2}
\]

Using binomial expansion and neglecting higher powers we get,
\[
v_t = v_0 \left( 1 + \frac{t}{2} \cdot \frac{1}{273} \right)
\]
\[
v_t = v_0 \left( 1 + \frac{t}{546} \right)
\]

Since \( v_0 = 331 \text{ m s}^{-1} \) at \( 0^\circ C \)
\[
v_t = 331 + 0.61 \text{ m s}^{-1}
\]
Thus the velocity of sound in air increases by 0.61 m s\(^{-1}\) per degree centigrade rise in temperature.

(iii) **Effect of density**

Consider two different gases at the same temperature and pressure with different densities. The velocity of sound in two gases are given by

\[
v_1 = \sqrt{\frac{\gamma_1 P}{\rho_1}} \quad \text{and} \quad v_2 = \sqrt{\frac{\gamma_2 P}{\rho_2}}
\]

\[
\therefore \frac{v_1}{v_2} = \sqrt{\frac{\rho_1}{\rho_2}} \cdot \frac{\gamma_2}{\gamma_1}
\]

For gases having same value of \(\gamma\), \(\frac{v_1}{v_2} = \sqrt{\frac{\rho_1}{\rho_2}}\)

The velocity of sound in a gas is inversely proportional to the square root of the density of the gas.

(iv) **Effect of humidity**

When the humidity of air increases, the amount of water vapour present in it also increases and hence its density decreases, because the density of water vapour is less than that of dry air. Since velocity of sound is inversely proportional to the square root of density, the sound travels faster in moist air than in dry air. Due to this reason it can be observed that on a rainy day sound travels faster.

(v) **Effect of wind**

The velocity of sound in air is affected by wind. If the wind blows with the velocity \(w\) along the direction of sound, then the velocity of sound increases to \(v + w\). If the wind blows in the opposite direction to the direction of sound, then the velocity of sound decreases to \(v - w\). If the wind blows at an angle \(\theta\) with the direction of sound, the effective velocity of sound will be \((v + w \cos \theta)\).

**Note:** In a medium, sound waves of different frequencies or wavelengths travel with the same velocity. Hence there is no effect of frequency on the velocity of sound.
7.3 Progressive wave

A progressive wave is defined as the onward transmission of the vibratory motion of a body in an elastic medium from one particle to the successive particle.

7.3.1 Equation of a plane progressive wave

An equation can be formed to represent generally the displacement of a vibrating particle in a medium through which a wave passes. Thus each particle of a progressive wave executes simple harmonic motion of the same period and amplitude differing in phase from each other.

Let us assume that a progressive wave travels from the origin O along the positive direction of X axis, from left to right (Fig. 7.6). The displacement of a particle at a given instant is

\[ y = a \sin \omega t \]  \hspace{1cm} \ldots (1)

where \( a \) is the amplitude of the vibration of the particle and \( \omega = 2\pi n \).

---

**Table 7.1 Velocity of sound in various media**

<table>
<thead>
<tr>
<th>Medium</th>
<th>Velocity (ms(^{-1}))</th>
</tr>
</thead>
<tbody>
<tr>
<td>Gases</td>
<td></td>
</tr>
<tr>
<td>Air 0(^{\circ}) C</td>
<td>331</td>
</tr>
<tr>
<td>Air 20(^{\circ}) C</td>
<td>343</td>
</tr>
<tr>
<td>Helium</td>
<td>965</td>
</tr>
<tr>
<td>Hydrogen</td>
<td>1284</td>
</tr>
<tr>
<td>Water 0(^{\circ}) C</td>
<td>1402</td>
</tr>
<tr>
<td>Water at 20(^{\circ}) C</td>
<td>1482</td>
</tr>
<tr>
<td>Sea water</td>
<td>1522</td>
</tr>
<tr>
<td>Solids</td>
<td></td>
</tr>
<tr>
<td>Aluminum</td>
<td>6420</td>
</tr>
<tr>
<td>Steel</td>
<td>5921</td>
</tr>
<tr>
<td>Granite</td>
<td>6000</td>
</tr>
</tbody>
</table>

(Not for examination)
The displacement of the particle P at a distance $x$ from O at a given instant is given by,

$$y = a \sin (\omega t - \phi) \quad \ldots (2)$$

If the two particles are separated by a distance $\lambda$, they will differ by a phase of $2\pi$. Therefore, the phase $\phi$ of the particle P at a distance $x$ is

$$\phi = \frac{2\pi}{\lambda} \cdot x$$

$$y = a \sin \left(\omega t - \frac{2\pi x}{\lambda}\right) \quad \ldots (3)$$

Since $\omega = 2\pi n = 2\pi \frac{v}{\lambda}$, the equation is given by

$$y = a \sin \left(\frac{2\pi vt}{\lambda} - \frac{2\pi x}{\lambda}\right)$$

$$y = a \sin \frac{2\pi}{\lambda} (vt - x) \quad \ldots (4)$$

Since $\omega = \frac{2\pi}{T}$, the eqn. (3) can also be written as

$$y = a \sin 2\pi \left(\frac{t}{T} - \frac{x}{\lambda}\right) \quad \ldots (5)$$

If the wave travels in opposite direction, the equation becomes.

$$y = a \sin 2\pi \left(\frac{t}{T} + \frac{x}{\lambda}\right) \quad \ldots (6)$$

(i) Variation of phase with time

The phase changes continuously with time at a constant distance.

At a given distance $x$ from O let $\phi_1$ and $\phi_2$ be the phase of a particle at time $t_1$ and $t_2$ respectively.

$$\phi_1 = 2\pi \left(\frac{t_1}{T} - \frac{x}{\lambda}\right)$$

$$\phi_2 = 2\pi \left(\frac{t_2}{T} - \frac{x}{\lambda}\right)$$
∴ \phi_2 - \phi_1 = 2\pi \left( \frac{t_2}{T} - \frac{t_1}{T} \right) = \frac{2\pi}{T} (t_2 - t_1)

\Delta \phi = \frac{2\pi}{T} \Delta t

This is the phase change \Delta \phi of a particle in time interval \Delta t. If \Delta t = T, \Delta \phi = 2\pi. This shows that after a time period T, the phase of a particle becomes the same.

(ii) Variation of phase with distance

At a given time t phase changes periodically with distance x. Let \phi_1 and \phi_2 be the phase of two particles at distance x_1 and x_2 respectively from the origin at a time t.

Then \phi_1 = 2\pi \left( \frac{t}{T} - \frac{x_1}{\lambda} \right)

\phi_2 = 2\pi \left( \frac{t}{T} - \frac{x_2}{\lambda} \right)

∴ \phi_2 - \phi_1 = -\frac{2\pi}{\lambda} (x_2 - x_1)

∴ \Delta \phi = -\frac{2\pi}{\lambda} \Delta x

The negative sign indicates that the forward points lag in phase when the wave travels from left to right.

When \Delta x = \lambda, \Delta \phi = 2\pi, the phase difference between two particles having a path difference \lambda is 2\pi.

7.3.2 Characteristics of progressive wave

1. Each particle of the medium executes vibration about its mean position. The disturbance progresses onward from one particle to another.

2. The particles of the medium vibrate with same amplitude about their mean positions.

3. Each successive particle of the medium performs a motion similar to that of its predecessor along the propagation of the wave, but later in time.

4. The phase of every particle changes from 0 to 2\pi.

5. No particle remains permanently at rest. Twice during each
vibration, the particles are momentarily at rest at extreme positions, different particles attain the position at different time.

6. Transverse progressive waves are characterised by crests and troughs. Longitudinal waves are characterised by compressions and rarefactions.

7. There is a transfer of energy across the medium in the direction of propagation of progressive wave.

8. All the particles have the same maximum velocity when they pass through the mean position.

9. The displacement, velocity and acceleration of the particle separated by $m\lambda$ are the same, where $m$ is an integer.

7.3.3 Intensity and sound level

If we hear the sound produced by violin, flute or harmonium, we get a pleasing sensation in the ear, whereas the sound produced by a gun, horn of a motor car etc. produce unpleasant sensation in the ear. The loudness of a sound depends on intensity of sound wave and sensitivity of the ear.

The intensity is defined as the amount of energy crossing per unit area per unit time perpendicular to the direction of propagation of the wave.

Intensity is measured in W m$^{-2}$.

The intensity of sound depends on (i) Amplitude of the source ($I \propto a^2$), (ii) Surface area of the source ($I \propto A$), (iii) Density of the medium ($I \propto \rho$), (iv) Frequency of the source ($I \propto n^2$) and (v) Distance of the observer from the source ($I \propto \frac{1}{r^2}$).

The lowest intensity of sound that can be perceived by the human ear is called threshold of hearing. It is denoted by $I_0$.

For sound of frequency 1 KHz, $I_0 = 10^{-12}$ W m$^{-2}$. The level of sound intensity is measured in decibel. According to Weber-Fechner law,

$\text{decibel level (β)} = 10 \log_{10} \left( \frac{I}{I_0} \right)$

where $I_0$ is taken as $10^{-12}$ W m$^{-2}$ which corresponds to the lowest sound intensity that can be heard. Its level is 0 dB. I is the maximum intensity that an ear can tolerate which is 1 W m$^{-2}$ equal to 120 dB.
\[ \beta = 10 \log_{10} \left( \frac{1}{10^{-12}} \right) \]

\[ \beta = 10 \log_{10} (10^{12}) \]

\[ \beta = 120 \ \text{dB}. \]

Table 7.2 gives the decibel value and power density (intensity) for various sources.

**Table 7.2 Intensity of sound sources**

*(NOT FOR EXAMINATION)*

<table>
<thead>
<tr>
<th>Source of sound</th>
<th>Sound intensity (dB)</th>
<th>Intensity (W m(^{-2}))</th>
</tr>
</thead>
<tbody>
<tr>
<td>Threshold of pain</td>
<td>120</td>
<td>1</td>
</tr>
<tr>
<td>Busy traffic</td>
<td>70</td>
<td>10(^{-5})</td>
</tr>
<tr>
<td>Conversation</td>
<td>65</td>
<td>3.2 \times 10(^{-6})</td>
</tr>
<tr>
<td>Quiet car</td>
<td>50</td>
<td>10(^{-7})</td>
</tr>
<tr>
<td>Quiet Radio</td>
<td>40</td>
<td>10(^{-8})</td>
</tr>
<tr>
<td>Whisper</td>
<td>20</td>
<td>10(^{-10})</td>
</tr>
<tr>
<td>Rustle of leaves</td>
<td>10</td>
<td>10(^{-11})</td>
</tr>
<tr>
<td>Threshold of hearing</td>
<td>0</td>
<td>10(^{-12})</td>
</tr>
</tbody>
</table>

7.4. Reflection of sound

Take two metal tubes A and B. Keep one end of each tube on a metal plate as shown in Fig. 7.7. Place a wrist watch at the open end of the tube A and interpose a cardboard between A and B. Now at a particular inclination of the tube B with the cardboard, ticking of the watch is clearly heard. The angle of reflection made by the tube B with the cardboard is equal to the angle of incidence made by the tube A with the cardboard.

7.4.1 Applications of reflection of sound waves

(i) **Whispering gallery**: The famous whispering gallery at
St. Paul’s Cathedral is a circular shaped chamber whose walls repeatedly reflect sound waves round the gallery, so that a person talking quietly at one end can be heard distinctly at the other end. This is due to multiple reflections of sound waves from the curved walls (Fig. 7.8).

(ii) Stethoscope: Stethoscope is an instrument used by physicians to listen to the sounds produced by various parts of the body. It consists of a long tube made of rubber or metal. When sound pulses pass through one end of the tube, the pulses get concentrated to the other end due to several reflections on the inner surface of the tube. Using this doctors hear the patients’ heart beat as concentrated rays.

(iii) Echo: Echoes are sound waves reflected from a reflecting surface at a distance from the listener. Due to persistence of hearing, we keep hearing the sound for \( \frac{1}{10} \)th of a second, even after the sounding source has stopped vibrating. Assuming the velocity of sound as 340 ms\(^{-1}\), if the sound reaches the obstacle and returns after 0.1 second, the total distance covered is 34 m. No echo is heard if the reflecting obstacle is less than 17 m away from the source.

7.5 Refraction of sound

This is explained with a rubber bag filled with carbon-di-oxide as shown in Fig. 7.9. The velocity of sound in carbon-di-oxide is less than that in air and hence the bag acts as a lens. If a whistle is used as a source S, the sound passes through the lens and converges at O which is located with the help of flame. The flame will be disturbed only at the point O.

When sound travels from one medium to another, it undergoes refraction.

7.5.1 Applications of refraction of sound

It is easier to hear the sound during night than during day-time.
During day time, the upper layers of air are cooler than the layers of air near the surface of the Earth. During night, the layers of air near the Earth are cooler than the upper layers of air. As sound travels faster in hot air, during day-time, the sound waves will be refracted upwards and travel a short distance on the surface of the Earth. On the other hand, during night the sound waves are refracted downwards to the Earth and will travel a long distance.

7.6 Superposition principle

When two waves travel in a medium simultaneously in such a way that each wave represents its separate motion, then the resultant displacement at any point at any time is equal to the vector sum of the individual displacements of the waves.

This principle is illustrated by means of a slinky in the Fig. 7.10(a).

1. In the figure, (i) shows that the two pulses pass each other.
2. In the figure, (ii) shows that they are at some distance apart.
3. In the figure, (iii) shows that they overlap partly.
4. In the figure, (iv) shows that resultant is maximum.

Fig. 7.10 b illustrates the same events but with pulses that are equal and opposite.

\[ \vec{Y}_1 \text{ and } \vec{Y}_2 \text{ are the displacements at a point, then the resultant displacement is given by } \vec{Y} = \vec{Y}_1 + \vec{Y}_2. \]

If \( |\vec{Y}_1| = |\vec{Y}_2| = a \), and if the two waves have their displacements in the same direction, then \( |\vec{Y}| = a + a = 2a \)

If the two waves have their displacements in the opposite direction, then \( |\vec{Y}| = a + (-a) = 0 \)

The principle of superposition of waves is applied in wave phenomena such as interference, beats and stationary waves.
7.6.1 Interference of waves

When two waves of same frequency travelling in the same direction in a medium superpose with each other, their resultant intensity is maximum at some points and minimum at some other points. This phenomenon of superposition is called interference.

Let us consider two simple harmonic waves of same frequency travelling in the same direction. If \( a_1 \) and \( a_2 \) are the amplitudes of the waves and \( \phi \) is the phase difference between them, then their instantaneous displacements are

\[
y_1 = a_1 \sin \omega t \quad \ldots (1)
\]
\[
y_2 = a_2 \sin (\omega t + \phi) \quad \ldots (2)
\]

According to the principle of superposition, the resultant displacement is represented by

\[
y = y_1 + y_2
\]
\[
= a_1 \sin \omega t + a_2 \sin (\omega t + \phi)
\]
\[
= a_1 \sin \omega t + a_2 (\sin \omega t \cdot \cos \phi + \cos \omega t \cdot \sin \phi)
\]
\[
= (a_1 + a_2 \cos \phi) \sin \omega t + a_2 \sin \phi \cos \omega t \quad \ldots (3)
\]

Put \( a_1 + a_2 \cos \phi = A \cos \theta \quad \ldots (4) \)
\[
a_2 \sin \phi = A \sin \theta \quad \ldots (5)
\]

where \( A \) and \( \theta \) are constants, then

\[
y = A \sin \omega t \cdot \cos \theta + A \cos \omega t \cdot \sin \theta
\]
or \( y = A \sin (\omega t + \theta) \quad \ldots (6) \)

This equation gives the resultant displacement with amplitude \( A \).

From eqn. (4) and (5)

\[
A^2 \cos^2 \theta + A^2 \sin^2 \theta
\]

\[
= (a_1 + a_2 \cos \phi)^2 + (a_2 \sin \phi)^2
\]

\[ : A^2 = a_1^2 + a_2^2 + 2a_1a_2 \cos \phi \]

\[ : A = \sqrt{a_1^2 + a_2^2 + 2a_1a_2 \cos \phi} \quad \ldots (7) \]

Also \( \tan \theta = \frac{a_2 \sin \phi}{a_1 + a_2 \cos \phi} \quad \ldots (8) \)
We know that intensity is directly proportional to the square of the amplitude

\( I \alpha A^2 \)

\( \therefore I \alpha \left( a_1^2 + a_2^2 + 2a_1a_2 \cos \phi \right) \) ... (9)

**Special cases**

The resultant amplitude \( A \) is maximum, when \( \cos \phi = 1 \) or \( \phi = 2m\pi \) where \( m \) is an integer (i.e) \( I_{\text{max}} \alpha (a_1 + a_2)^2 \)

The resultant amplitude \( A \) is minimum when

\[ \cos \phi = -1 \text{ or } \phi = (2m + 1)\pi \]

\[ I_{\text{min}} \alpha (a_1 - a_2)^2 \]

The points at which interfering waves meet in the same phase \( \phi = 2m\pi \) i.e. 0, 2\( \pi \), 4\( \pi \), ... are points of maximum intensity, where constructive interference takes place. The points at which two interfering waves meet out of phase \( \phi = (2m + 1)\pi \) i.e. \( \pi \), 3\( \pi \), ... are called points of minimum intensity, where destructive interference takes place.

**7.6.2 Experimental demonstration of interference of sound**

The phenomenon of interference between two longitudinal waves in air can be demonstrated by Quincke’s tube shown in Fig. 7.11.

Quincke’s tube consists of U shaped glass tubes \( A \) and \( B \). The tube \( SAR \) has two openings at \( S \) and \( R \). The other tube \( B \) can slide over the tube \( A \). A sound wave from \( S \) travels along both the paths \( SAR \) and \( SBR \) in opposite directions and meet at \( R \).

If the path difference between the two waves (i.e) \( SAR \sim SBR \) is an integral multiple of wavelength, intensity of sound will be maximum due to constructive interference.

i.e \( SAR \sim SBR = m\lambda \)

The corresponding phase difference \( \phi \) between the two waves is even multiples of \( \pi \). (i.e) \( \phi = m \cdot 2\pi \) where \( m = 0, 1, 2, 3 \ldots \)
If the tube B is gradually slided over A, a stage is reached when the intensity of sound is zero at \( R \) due to destructive interference. Then no sound will be heard at \( R \).

If the path difference between the waves is odd multiples of \( \frac{\lambda}{2} \), intensity of sound will be minimum.

\[ \text{i.e. } \text{SAR} \sim \text{SBR} = (2m + 1)\frac{\lambda}{2} \]

The corresponding phase difference \( \phi \) between the two waves is odd multiples of \( \pi \). (i.e) \( \phi = (2m + 1)\pi \) where \( m = 0, 1, 2, 3 \ldots \).

### 7.6.3 Beats

When two waves of nearly equal frequencies travelling in a medium along the same direction superimpose upon each other, beats are produced. The amplitude of the resultant sound at a point rises and falls regularly.

The intensity of the resultant sound at a point rises and falls regularly with time. When the intensity rises to maximum we call it as waxing of sound, when it falls to minimum we call it as waning of sound.

The phenomenon of waxing and waning of sound due to interference of two sound waves of nearly equal frequencies are called beats. The number of beats produced per second is called beat frequency, which is equal to the difference in frequencies of two waves.

#### Analytical method

Let us consider two waves of slightly different frequencies \( n_1 \) and \( n_2 \) (\( n_1 \sim n_2 < 10 \)) having equal amplitude travelling in a medium in the same direction.

At time \( t = 0 \), both waves travel in same phase.

The equations of the two waves are

\[ y_1 = a \sin \omega_1 t \]
\[ y_1 = a \sin (2\pi n_1) t \] ...(1)

\[ y_2 = a \sin \omega_2 t \]
\[ = a \sin (2\pi n_2) t \] ...(2)

When the two waves superimpose, the resultant displacement is given by

\[ y = y_1 + y_2 \]
\[ = a \sin (2\pi n_1) t + a \sin (2\pi n_2) t \] ...(3)

Therefore

\[ y = 2a \sin 2\pi \left( \frac{n_1 + n_2}{2} \right) t \cos 2\pi \left( \frac{n_1 - n_2}{2} \right) t \] ...(4)

Substitute \( A = 2a \cos 2\pi \left( \frac{n_1 - n_2}{2} \right) t \) and \( n = \frac{n_1 + n_2}{2} \) in equation (4)

\[ \therefore \ y = A \sin 2\pi n t \]

This represents a simple harmonic wave of frequency \( n = \frac{n_1 + n_2}{2} \) and amplitude \( A \) which changes with time.

(i) The resultant amplitude is maximum (i.e) \( \pm 2a \), if

\[ \cos 2\pi \left( \frac{n_1 - n_2}{2} \right) t = \pm 1 \]

\[ \therefore 2\pi \left( \frac{n_1 - n_2}{2} \right) t = m\pi \] (where \( m = 0, 1, 2 \ldots \)) or \( (n_1 - n_2) t = m \)

The first maximum is obtained at \( t_1 = 0 \)

The second maximum is obtained at

\[ t_2 = \frac{1}{n_1 - n_2} \]

The third maximum at \( t_3 = \frac{2}{n_1 - n_2} \) and so on.

The time interval between two successive maxima is

\[ t_2 - t_1 = t_3 - t_2 = \frac{1}{n_1 - n_2} \]

Hence the number of beats produced per second is equal to the reciprocal of the time interval between two successive maxima.
(ii) The resultant amplitude is minimum (i.e) equal to zero, if
\[ \cos 2\pi \left( \frac{n_1 - n_2}{2} \right) t = 0 \]

(i.e) \[ 2\pi \left( \frac{n_1 - n_2}{2} \right) t = \frac{\pi}{2} + m\pi = (2m + 1) \frac{\pi}{2} \] or \( (n_1 - n_2)t = \frac{(2m+1)}{2} \)

where \( m = 0, 1, 2 \ldots \)

The first minimum is obtained at
\[ t_1' = \frac{1}{2(n_1 - n_2)} \]

The second minimum is obtained at
\[ t_2' = \frac{3}{2(n_1 - n_2)} \]

The third minimum is obtained at
\[ t_3' = \frac{5}{2(n_1 - n_2)} \] and so on

Time interval between two successive minima is
\[ t_2' - t_1' = t_3' - t_2' = \frac{1}{n_1 - n_2} \]

Hence, the number of beats produced per second is equal to the reciprocal of time interval between two successive minima.

7.6.4 Uses of beats

(i) The phenomenon of beats is useful in tuning two vibrating bodies in unison. For example, a sonometer wire can be tuned in unison with a tuning fork by observing the beats. When an excited tuning fork is kept on the sonometer and if the sonometer wire is also excited, beats are heard, when the frequencies are nearly equal. If the length of the wire is adjusted carefully so that the number of beats gradually decreases to zero, then the two are said to be in unison. Most of the musical instruments are made to be in unison based on this method.

(ii) The frequency of a tuning fork can be found using beats. A standard tuning fork of frequency \( N \) is excited along with the experimental fork. If the number of beats per second is \( n \), then the frequency of experimental tuning fork is \( N + n \). The experimental tuning
fork is then loaded with a little bees’ wax, thereby decreasing its frequency. Now the observations are repeated. If the number of beats increases, then the frequency of the experimental tuning fork is \( N-n \), and if the number of beats decreases its frequency is \( N+n \).

### 7.6.5 Stationary waves

When two progressive waves of same amplitude and wavelength travelling along a straight line in opposite directions superimpose on each other, stationary waves are formed.

**Analytical method**

Let us consider a progressive wave of amplitude \( a \) and wavelength \( \lambda \) travelling in the direction of X axis.

\[
y_1 = a \sin 2\pi \left( \frac{t}{T} - \frac{x}{\lambda} \right)
\]

This wave is reflected from a free end and it travels in the negative direction of X axis, then

\[
y_2 = a \sin 2\pi \left( \frac{t}{T} + \frac{x}{\lambda} \right)
\]

According to principle of superposition, the resultant displacement is

\[
y = y_1 + y_2
\]

\[
= a \left[ \sin 2\pi \left( \frac{t}{T} - \frac{x}{\lambda} \right) + \sin 2\pi \left( \frac{t}{T} + \frac{x}{\lambda} \right) \right]
\]

\[
= a \left[ 2\sin \frac{2\pi}{T} \cos \frac{2\pi x}{\lambda} \right]
\]

\[
\therefore \quad y = 2a \cos \frac{2\pi x}{\lambda} \sin \frac{2\pi t}{T}
\]

This is the equation of a stationary wave.

(i) At points where \( x = 0, \frac{\lambda}{2}, \frac{3\lambda}{2}, \frac{5\lambda}{2} \), the values of \( \cos \frac{2\pi x}{\lambda} = \pm 1 \)

\[
\therefore A = \pm 2a. \text{ At these points the resultant amplitude is maximum. They are called antinodes (Fig. 7.13).}
\]

(ii) At points where \( x = \frac{\lambda}{4}, \frac{3\lambda}{4}, \frac{5\lambda}{4} \), the values of \( \cos \frac{2\pi x}{\lambda} = 0 \).

\[
\therefore A = 0. \text{ The resultant amplitude is zero at these points. They are} \]
called nodes (Fig. 7.16).

The distance between any two successive antinodes or nodes is equal to \( \frac{\lambda}{2} \) and the distance between an antinode and a node is \( 2\lambda \).

(iii) When \( t = 0, \frac{T}{2}, T, \frac{3T}{2}, 2T, \ldots \) then \( \sin \frac{2\pi t}{T} = 0 \), the displacement is zero.

(iv) When \( t = \frac{T}{4}, \frac{3T}{4}, \frac{5T}{4}, \ldots \) etc, \( \sin \frac{2\pi t}{T} = \pm 1 \), the displacement is maximum.

### 7.6.6 Characteristics of stationary waves

1. The waveform remains stationary.
2. Nodes and antinodes are formed alternately.
3. The points where displacement is zero are called nodes and the points where the displacement is maximum are called antinodes.
4. Pressure changes are maximum at nodes and minimum at antinodes.
5. All the particles except those at the nodes, execute simple harmonic motions of same period.
6. Amplitude of each particle is not the same, it is maximum at antinodes decreases gradually and is zero at the nodes.
7. The velocity of the particles at the nodes is zero. It increases gradually and is maximum at the antinodes.
8. Distance between any two consecutive nodes or antinodes is equal to \( \frac{\lambda}{2} \), whereas the distance between a node and its adjacent antinode is equal to \( \frac{\lambda}{4} \).
9. There is no transfer of energy. All the particles of the medium pass through their mean position simultaneously twice during each vibration.
10. Particles in the same segment vibrate in the same phase and...
between the neighbouring segments, the particles vibrate in opposite phase.

7.7 Standing waves in strings

In musical instruments like sitar, violin, etc. sound is produced due to the vibrations of the stretched strings. Here, we shall discuss the different modes of vibrations of a string which is rigidly fixed at both ends.

When a string under tension is set into vibration, a transverse progressive wave moves towards the end of the wire and gets reflected. Thus stationary waves are formed.

7.7.1 Sonometer

The sonometer consists of a hollow sounding box about a metre long. One end of a thin metallic wire of uniform cross-section is fixed to a hook and the other end is passed over a pulley and attached to a weight hanger as shown in Fig. 7.14. The wire is stretched over two knife edges P and Q by adding sufficient weights on the hanger. The distance between the two knife edges can be adjusted to change the vibrating length of the wire.

A transverse stationary wave is set up in the wire. Since the ends are fixed, nodes are formed at P and Q and antinode is formed in the middle.

The length of the vibrating segment is \( l = \lambda/2 \)

\[
\therefore \lambda = 2l. \text{ If } n \text{ is the frequency of vibrating segment, then}
\]

\[
n = \frac{v}{\lambda} = \frac{v}{2l} \quad \text{...(1)}
\]

We know that \( v = \frac{T}{\sqrt{m}} \) where T is the tension and m is the mass per unit length of the wire.

\[
\therefore n = \frac{1}{2l} \sqrt{\frac{T}{m}} \quad \text{...(2)}
\]
Modes of vibration of stretched string

(i) Fundamental frequency

If a wire is stretched between two points, a transverse wave travels along the wire and is reflected at the fixed end. A transverse stationary wave is thus formed as shown in Fig. 7.15.

When a wire AB of length \( l \) is made to vibrate in one segment then

\[
 l = \frac{\lambda_1}{2}. 
\]

\[ \therefore \lambda_1 = 2l. \] This gives the lowest frequency called fundamental frequency \( n_1 = \frac{v}{\lambda_1} \)

\[ \therefore n_1 = \frac{1}{2l} \sqrt{\frac{T}{m}} \quad \text{(3)} \]

(ii) Overtones in stretched string

If the wire AB is made to vibrate in two segments then

\[ l = \frac{\lambda_2}{2} + \frac{\lambda_2}{2} \]

\[ \therefore \lambda_2 = l. \]

But, \( n_2 = \frac{v}{\lambda_2} \)

\[ \therefore n_2 = \frac{l}{2l} \sqrt{\frac{T}{m}} = 2n_1 \quad \text{(4)} \]

\( n_2 \) is the frequency of the first overtone.

Since the frequency is equal to twice the fundamental, it is also known as second harmonic.

Similarly, higher overtones are produced, if the wire vibrates with more segments. If there are \( P \) segments, the length of each segment is

\[ \frac{l}{P} = \frac{\lambda_P}{2} \quad \text{or} \quad \lambda_P = \frac{2l}{P} \]

\[ \therefore \text{Frequency } n_P = \frac{P}{2l} \sqrt{\frac{T}{m}} = P n_1 \quad \text{(5)} \]

(i.e) \( P^{th} \) harmonic corresponds to \( (P-1)^{th} \) overtone.
7.7.2 Laws of transverse vibrations of stretched strings

The laws of transverse vibrations of stretched strings are (i) the law of length (ii) law of tension and (iii) the law of mass.

(i) For a given wire \((m \text{ is constant})\), when \(T\) is constant, the fundamental frequency of vibration is inversely proportional to the vibrating length (i.e)

\[ n \propto \frac{1}{l} \text{ or } nl = \text{constant}. \]

(ii) For constant \(l \text{ and } m\), the fundamental frequency is directly proportional to the square root of the tension (i.e) \(n \propto \sqrt{T}\).

(iii) For constant \(l \text{ and } T\), the fundamental frequency varies inversely as the square root of the mass per unit length of the wire (i.e) \(n \propto \frac{1}{\sqrt{m}}\).

7.8 Vibrations of air column in pipes

Musical wind instruments like flute, clarinet etc. are based on the principle of vibrations of air columns. Due to the superposition of the incident wave and the reflected wave, longitudinal stationary waves are formed in the pipe.

7.8.1 Organ pipes

Organ pipes are musical instruments which are used to produce musical sound by blowing air into the pipe. Organ pipes are two types (i) closed organ pipes, closed at one end (ii) open organ pipe, open at both ends.

(i) Closed organ pipe: If the air is blown lightly at the open end of the closed organ pipe, then the air column vibrates (Fig. 7.16a) in the fundamental mode. There is a node at the closed end and an antinode at the open end. If \(l\) is the length of the tube,

\[ l = \frac{\lambda}{4} \text{ or } \lambda_1 = 4l \quad \ldots \text{(1)} \]

If \(n_1\) is the fundamental frequency of the waves in a closed pipe (Fundamental mode)
vibrations and \( v \) is the velocity of sound in air, then

\[
\nu_1 = \frac{\nu}{\lambda_1} = \frac{\nu}{4l}
\]

If air is blown strongly at the open end, frequencies higher than fundamental frequency can be produced. They are called overtones. Fig. 7.16b & Fig. 7.16c shows the mode of vibration with two or more nodes and antinodes.

\[
l = \frac{3\lambda_3}{4} \quad \text{or} \quad \lambda_3 = \frac{4l}{3}
\]

\[
\therefore \ n_3 = \frac{v}{\lambda_3} = \frac{3v}{4l} = 3n_1
\]

This is the first overtone or third harmonic.

Similarly \( n_5 = \frac{5v}{4l} = 5n_1 \) \( \ldots(5) \)

This is called as second overtone or fifth harmonic.

Therefore the frequency of \( p \)th overtone is \( (2p + 1) n_1 \) where \( n_1 \) is the fundamental frequency. In a closed pipe only odd harmonics are produced. The frequencies of harmonics are in the ratio of 1 : 3 : 5....

(ii) Open organ pipe - When air is blown into the open organ pipe, the air column vibrates in the fundamental mode Fig. 7.17a. Antinodes are formed at the ends and a node is formed in the middle of the pipe. If \( l \) is the length of the pipe, then

\[
l = \frac{\lambda_1}{2} \quad \text{or} \quad \lambda_1 = 2l
\]

\[
v = n_1 \lambda_1 = n_1 2l
\]

The fundamental frequency

\[
n_1 = \frac{\nu}{2l}
\]

In the next mode of vibration additional nodes and antinodes are formed as shown in Fig. 7.17a Stationary waves in an open pipe (Fundamental mode)
Fig. 7.17b and Fig. 7.17c.

\[ l = \lambda_2 \text{ or } v = n_2\lambda_2 = n_2 \cdot l \]

\[ \therefore n_2 = \left( \frac{v}{T} \right) = 2n_1 \quad \text{(3)} \]

This is the first overtone or second harmonic.

Similarly, \( n_3 = \frac{v}{\lambda_3} = \frac{3v}{2l} = 3n_1 \) \quad \text{(4)}

This is the second overtone or third harmonic.

Therefore the frequency of \( P^{th} \) overtone is \( (P + 1) n_1 \) where \( n_1 \) is the fundamental frequency.

The frequencies of harmonics are in the ratio of 1 : 2 : 3 ....

7.9 Resonance air column apparatus

The resonance air column apparatus consists of a glass tube \( G \) about one metre in length (Fig. 7.18) whose lower end is connected to a reservoir \( R \) by a rubber tube.

The glass tube is mounted on a vertical stand with a scale attached to it. The glass tube is partly filled with water. The level of water in the tube can be adjusted by raising or lowering the reservoir.

Fig. 7.18 Resonance air column apparatus
A vibrating tuning fork of frequency \( n \) is held near the open end of the tube. The length of the air column is adjusted by changing the water level. The air column of the tube acts like a closed organ pipe. When this air column resonates with the frequency of the fork the intensity of sound is maximum.

Here longitudinal stationary wave is formed with node at the water surface and an antinode near the open end. If \( l_1 \) is the length of the resonating air column

\[
\frac{\lambda}{4} = l_1 + e
\] ..(1)

where \( e \) is the end correction.

The length of air column is increased until it resonates again with the tuning fork. If \( l_2 \) is the length of the air column.

\[
\frac{3\lambda}{4} = l_2 + e
\] ...(2)

From equations (1) and (2)

\[
\frac{\lambda}{2} = (l_2 - l_1)
\] ...(3)

The velocity of sound in air at room temperature

\[
v = n\lambda = 2n (l_2 - l_1)
\] ...(4)

**End correction**

The antinode is not exactly formed at the open end, but at a small distance above the open end. This is called the end correction.

As \( l_1 + e = \frac{\lambda}{4} \) and \( l_2 + e = \frac{3\lambda}{4} \)

\[
e = \frac{(l_2 - 3l_1)}{2}
\]

It is found that \( e = 0.61r \), where \( r \) is the radius of the glass tube.

**7.10 Doppler effect**

The whistle of a fast moving train appears to increase in pitch as it approaches a stationary observer and it appears to decrease as the train moves away from the observer. This apparent change in frequency was first observed and explained by Doppler in 1845.
The phenomenon of the apparent change in the frequency of sound due to the relative motion between the source of sound and the observer is called Doppler effect.

The apparent frequency due to Doppler effect for different cases can be deduced as follows.

(i) Both source and observer at rest

Suppose S and O are the positions of the source and the observer respectively. Let \( n \) be the frequency of the sound and \( v \) be the velocity of sound. In one second, \( n \) waves produced by the source travel a distance \( SO = v \) (Fig. 7.19a).

The wavelength is \( \lambda = \frac{v}{n} \).

(ii) When the source moves towards the stationary observer

If the source moves with a velocity \( v_s \) towards the stationary observer, then after one second, the source will reach \( S' \), such that \( SS' = v_s \). Now \( n \) waves emitted by the source will occupy a distance of \( (v-v_s) \) only as shown in Fig. 7.19b.

Therefore the apparent wavelength of the sound is
\[
\lambda' = \frac{v - v_s}{n}
\]

The apparent frequency
\[
n' = \frac{v}{\lambda'} = \left(\frac{v}{v - v_s}\right) n \quad \text{...(1)}
\]

As \( n' > n \), the pitch of the sound appears to increase.

When the source moves away from the stationary observer

If the source moves away from the stationary observer with velocity \( v_s \), the apparent frequency will be given by
\[
n' = \left(\frac{v}{v - (-v_s)}\right)n = \left(\frac{v}{v + v_s}\right)n \quad \text{...(2)}
\]
As \( n' < n \), the pitch of the sound appears to decrease.

**(iii) Source is at rest and observer in motion**

S and O represent the positions of source and observer respectively. The source S emits \( n \) waves per second having a wavelength \( \lambda = \frac{v}{n} \).

Consider a point A such that OA contains \( n \) waves, which crosses the ear of the observer in one second (Fig. 7.20a). (i.e) when the first wave is at the point A, the \( n \)th wave will be at O, where the observer is situated.

When the observer moves towards the stationary source

Suppose the observer is moving towards the stationary source with velocity \( v_o \). After one second the observer will reach the point \( O' \) such that \( OO' = v_o \). The number of waves crossing the observer will be \( n \) waves in the distance OA in addition to the number of waves in the distance \( OO' \) which is equal to \( \frac{v_o}{\lambda} \) as shown in Fig. 7.20b.

Therefore, the apparent frequency of sound is

\[
n' = n + \frac{v_o}{\lambda} = n + \left( \frac{v_o}{v} \right) n
\]

\[
\therefore n' = \left( \frac{v + v_o}{v} \right) n
\]

As \( n' > n \), the pitch of the sound appears to increase.

When the observer moves away from the stationary source

\[
n' = \left[ \frac{v + (v_o) f}{v} \right] n
\]
\[ n' = \left( \frac{v - v_o}{v} \right) n \]  

...(4)

As \( n' < n \), the pitch of sound appears to decrease.

**Note** : If the source and the observer move along the same direction, the equation for apparent frequency is

\[ n' = \left( \frac{v - v_o}{v - v_s} \right) n \]  

...(5)

Suppose the wind is moving with a velocity \( W \) in the direction of propagation of sound, the apparent frequency is

\[ n' = \left( \frac{v + W - v_o}{v + W - v_s} \right) n \]  

...(6)

**Applications of Doppler effect**

**(i) To measure the speed of an automobile**

An electromagnetic wave is emitted by a source attached to a police car. The wave is reflected by a moving vehicle, which acts as a moving source. There is a shift in the frequency of the reflected wave. From the frequency shift using beats, the speeding vehicles are trapped by the police.

**(ii) Tracking a satellite**

The frequency of radio waves emitted by a satellite decreases as the satellite passes away from the Earth. The frequency received by the Earth station, combined with a constant frequency generated in the station gives the beat frequency. Using this, a satellite is tracked.

**(iii) RADAR (RADIO DETECTION AND RANGING)**

A RADAR sends high frequency radiowaves towards an aeroplane. The reflected waves are detected by the receiver of the radar station. The difference in frequency is used to determine the speed of an aeroplane.

**(iv) SONAR (SOUND NAVIGATION AND RANGING)**

Sound waves generated from a ship fitted with SONAR are transmitted in water towards an approaching submarine. The frequency of the reflected waves is measured and hence the speed of the submarine is calculated.
Solved Problems

7.1 What is the distance travelled by sound in air when a tuning fork of frequency 256 Hz completes 25 vibrations? The speed of sound in air is 343 m s\(^{-1}\).

**Data**: \( v = 343 \text{ m s}^{-1} \), \( n = 256 \text{ Hz} \), \( d = ? \)

**Solution**: \( v = n\lambda \)

\[
\therefore \lambda = \frac{343}{256} = 1.3398 \text{ m}
\]

Wavelength is the distance travelled by the wave in one complete vibration of the tuning fork.

\( \therefore \) Distance travelled by sound wave in 25 vibrations = \( 25 \times 1.3398 \)

Distance travelled by sound wave is = 33.49 m

7.2 Ultrasonic sound of frequency 100 kHz emitted by a bat is incident on a water surface. Calculate the wavelength of reflected sound and transmitted sound? (speed of sound in air 340 m s\(^{-1}\) and in water 1486 m s\(^{-1}\))

**Data**: \( n = 100 \text{ kHz} = 10^5 \text{ Hz} \), \( v_a = 340 \text{ m s}^{-1} \), \( v_w = 1486 \text{ m s}^{-1} \);

\( \lambda_a = ? \), \( \lambda_w = ? \)

Wavelength of reflected sound \( \lambda_a = \frac{v_a}{n} \)

\[
\lambda_a = \frac{340}{10^5} = 3.4 \times 10^{-3} \text{ m}
\]

Wavelength of transmitted sound \( \lambda_w = \frac{v_w}{n} \)

\[
\lambda_w = \frac{1486}{10^5} = 1.486 \times 10^{-2} \text{ m}
\]

7.3 A string of mass 0.5 kg and length 50 m is stretched under a tension of 400 N. A transverse wave of frequency 10 Hz travels through the wire. (i) Calculate the wave velocity and wavelength. (ii) How long does the disturbance take to reach the other end?

**Data**: \( m = 0.5 \text{ kg} \), length of the wire = 50 m; \( T = 400 \text{ N} \); \( n=10 \text{ Hz} \)

\( v = ? \); \( \lambda = ? \); \( t = ? \)

**Solution**: mass per unit length \( m = \frac{\text{mass of the wire}}{\text{length of the wire}} \)

\[
m = \frac{0.5}{50} = 0.01 \text{ kg m}^{-1}
\]
Velocity in the stretched string \( v = \sqrt{\frac{T}{m}} \)

\[
v = \sqrt{\frac{400}{0.01}} = 200 \text{ m s}^{-1}
\]

\( v = n\lambda \\
200 = 10\lambda \\
\therefore \lambda = 20 \text{ m}
\]

The length of the wire = 50 m

\( \therefore \) Time taken for the transverse wave to travel a distance 50 m = \( \frac{50}{200} = 0.25 \text{ s} \)

7.4 Determine the velocity and wavelength of sound of frequency 256 Hz travelling in water of Bulk modulus 0.022 \( \times \) 10\(^{11} \) Pa

**Data :** \( k = 0.022 \times 10^{11} \) Pa, \( \rho = 1000 \) kg m\(^{-3} \), \( n = 256 \) Hz

**Solution :** Velocity of sound in water \( v = \sqrt{\frac{k}{\rho}} \)

\[
v = \sqrt{\frac{0.022\times10^{11}}{1000}} = 1483 \text{ ms}^{-1}
\]

\( \therefore \lambda = \frac{v}{n} = \frac{1483}{256} = 5.79 \text{ m} \)

7.5 Calculate the speed of longitudinal wave in air at 27\(^{\circ}\) C (The molecular mass of air is 28.8 g mol\(^{-1} \), \( \gamma \) for air is 1.4, \( R = 8.314 \) J mol\(^{-1} \) K\(^{-1} \))

**Data :** \( m = 28.8 \times 10^{-3} \) kg mol\(^{-1} \), \( \gamma = 1.4, \)

\[
R = 8.314 \text{ J mol}^{-1} \text{ K}^{-1}, \quad T = 27^{\circ} \text{C} = 300 \text{ K}
\]

**Solution :** \( v = \sqrt{\frac{\gamma RT}{m}} = \sqrt{\frac{1.4\times8.314\times300}{28.8\times10^{-3}}} \)

\( v = 348.2 \text{ m s}^{-1} \)

7.6 For air at NTP, the density is 0.001293 g cm\(^{-3} \). Calculate the velocity of longitudinal wave (i) using Newton’s formula (ii) Laplace’s correction

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Data: \( \gamma = 1.4, P = 1.013 \times 10^5 \text{ N m}^{-2}, \)
\( \rho = 0.001293 \times 10^3 \text{ kg m}^{-3}. \)

**Solution**: By Newton’s formula the velocity of longitudinal wave

\[
v = \sqrt{\frac{P}{\rho}} = \sqrt{\frac{1.013 \times 10^5}{0.001293 \times 10^3}} \approx 531.18 \text{ m s}^{-1}
\]

By Laplace’s formula

\[
v = \sqrt{\frac{\gamma P}{\rho}} = \sqrt{\frac{1.4 \times 1.013 \times 10^5}{0.001293 \times 10^3}} \approx 331.18 \text{ m s}^{-1}
\]

7.7 The velocity of sound at 27°C is 347 m s\(^{-1}\). Calculate the velocity of sound in air at 627°C.

**Data**: \( v_{27} = 347 \text{ m s}^{-1}, v_{627} = ? \)

**Solution**: \(v \propto \sqrt{T}\)

\[
\frac{v_{27}}{v_{627}} = \sqrt{\frac{273 + 27}{273 + 627}} = \sqrt{\frac{300}{900}}
\]

\[
\frac{v_{27}}{v_{627}} = \frac{1}{\sqrt{3}}
\]

\[
\therefore v_{627} = v_{27} \times \sqrt{3} = 347 \times \sqrt{3} = 601 \text{ m s}^{-1}
\]

Velocity of sound in air at 627°C is 601 m s\(^{-1}\)

7.8 The equation of a progressive wave is \(y = 0.50 \sin (500t - 0.025x)\), where \(y\), \(t\) and \(x\) are in cm, second and metre. Calculate (i) amplitude (ii) angular frequency (iii) period (iv) wavelength and (v) speed of propagation of wave.

**Solution**: The general equation of a progressive wave is given by

\[
y = \alpha \sin \left( \omega t - \frac{2\pi}{\lambda} x \right)
\]

given \(y = 0.50 \sin (500t - 0.025x)\)

comparing the two equations,
(i) amplitude \( a = 0.50 \times 10^{-2} \text{ m} \)

(ii) angular frequency \( \omega = 500 \text{ rad s}^{-1} \)

(iii) time period \( T = \frac{2\pi}{\omega} = \frac{2\pi}{500} = \frac{\pi}{250} \text{ s} \)

(iv) wavelength \( \lambda = \frac{2\pi}{0.025} \text{ m} \)

\[ \lambda = 80\pi = 251.2 \text{ m} \]

(v) wave velocity \( v = n\lambda \)

\[ v = \frac{250}{\pi} \times 80\pi = 2 \times 10^4 \text{ m s}^{-1} \]

7.9 A source of sound radiates energy uniformly in all directions at a rate of 2 watt. Find the intensity (i) in \( W \text{ m}^2 \) and (ii) in decibels, at a point 20 m from the source.

Data: Power = 2 watt, \( r = 20 \text{ m} \)

Solution: Intensity of sound \( I = \frac{\text{Power}}{\text{area}} \)

\[ I = \frac{2}{4\pi(20)^2} \]

(A spherical surface of radius 20 m with source of sound as centre is imagined)

\[ I = 4 \times 10^4 \text{ W m}^2 \]

Intensity = \( 10 \log_{10} \left( \frac{I}{I_0} \right) \)

\[ = 10 \log_{10} \left( \frac{4 \times 10^4}{10^{-12}} \right) \quad (\because \ I_0 = 10^{-12}) \]

\[ = 10 \log_{10} (4 \times 10^8) \]

Intensity = 86 dB

7.10 Two tuning forks A and B when sounded together produce 4 beats. If A is in unison with the 0.96 m length of a sonometer wire under a tension, B is in unison with 0.97 m length of the same wire under same tension. Calculate the frequencies of the forks.

Data: \( l_1 = 0.96 \text{ m}; l_2 = 0.97 \text{ m}; n_1 = ?; n_2 = ? \)

\( l_1 < l_2 \quad \therefore \quad n_1 > n_2 \)
Solution: Let \( n_1 = n \) and \( n_2 = n - 4 \)

According to 1st law of transverse vibrations

\[
\begin{align*}
  n_1 l_1 &= n_2 l_2 \\
  n \times 0.96 &= (n-4) \times 0.97 \\
  n(0.97 - 0.96) &= 3.88
\end{align*}
\]

\[
\therefore n = \frac{3.88}{0.01} = 388 \text{ Hz}
\]

\[
\therefore n_2 = 388 - 4 = 384 \text{ Hz}
\]

The frequency of the fork A is \( n_1 = 388 \text{ Hz} \).
The frequency of the fork B is \( n_2 = 384 \text{ Hz} \).

7.11 A string of length 1 m and mass \( 5 \times 10^{-4} \) kg fixed at both ends is under a tension of 20 N. If it vibrates in two segments, determine the frequency of vibration of the string.

Data: The string vibrates with 2 segments,
\( P = 2 \) loops, \( l = 1 \) m, \( m = 5 \times 10^{-4} \) kg m\(^{-1}\), \( T = 20 \) N

Solution: Frequency of vibration \( n = \frac{P}{2l} \sqrt{\frac{T}{m}} \)

\[
\therefore n = \frac{2}{2 \times 1} \sqrt{\frac{20}{5 \times 10^{-4}}} = 200 \text{ Hz}
\]

7.12 A stretched string made of aluminium is vibrating at its fundamental frequency of 512 Hz. What is the fundamental frequency of a second string made from the same material which has a diameter and length twice that of the original and which is subjected to three times the force of the original?

Data: \( n = 512 \) Hz. In the second case, tension = 3T, length = 2l, radius = 2r

Solution: Let \( l \) be the length, \( T \) be the tension and \( r \) be the radius of the wire, then

\[
n = \frac{1}{2l} \sqrt{\frac{T}{m}}
\]

Mass per unit length can be written as the product of cross-sectional area of the wire and density (i.e) \( m = \pi r^2 d \)
\[ 512 = \frac{1}{2\pi} \sqrt{\frac{T}{\pi^2 d}} \quad \text{...(1)} \]

In the second case
\[ n = \frac{1}{2 \times 2\pi} \sqrt{\frac{3T}{\pi(2r)^2 d}} \quad \text{...(2)} \]

Dividing the second equation by first equation
\[ \frac{n}{512} = \frac{1}{2} \sqrt{\frac{3}{(2r)^2}} \quad \text{(i.e) } n = \frac{512}{4} \sqrt{3} = 222 \text{ Hz} \]

7.13 The third overtone of a closed pipe is found to be in unison with the first overtone of an open pipe. Determine the ratio of the lengths of the pipes.

**Solution:** Let \( l_1 \) and \( l_2 \) be the lengths of the closed pipe and open pipe respectively, \( n_1 \) and \( n_2 \) are their fundamental frequencies.

For closed pipe \( n_1 = \frac{v}{4l_1} \)

For open pipe \( n_2 = \frac{v}{2l_2} \)

Third overtone of closed pipe = \((2P + 1) \) \( n_1 = (2 \times 3 + 1) \) \( n_1 = 7n_1 \)

First overtone of open pipe = \((P + 1) \) \( n_2 = (1 + 1) \) \( n_2 = 2n_2 \)

\[ \therefore 7n_1 = 2n_2 \]

\[ 7 \times \frac{v}{4l_1} = 2 \times \frac{v}{2l_2} \]

\[ \therefore \frac{l_1}{l_2} = \frac{7}{4} \]

7.14 The shortest length of air in a resonance tube which resonates with a tuning fork of frequency 256 Hz is 32 cm. The corresponding length for the fork of frequency 384 Hz is 20.8 cm. Calculate the end correction and velocity of sound in air.

**Data:** \( n_1 = 256 \text{ Hz, } l_1 = 32 \times 10^{-2} \text{ m} \)

\( n_2 = 384 \text{ Hz, } l_2 = 20.8 \times 10^{-2} \text{ m} \)

**Solution:** In a closed pipe \( n = \frac{v}{4(l + e)} \)

For the first tuning fork, \( 256 = \frac{v}{4(32 + e) \times 10^{-2}} \) and
for the second tuning fork, \( 384 = \frac{v}{4(20.8 + e) \times 10^2} \)

Dividing the first equation by second equation,

\[
\frac{256}{384} = \frac{20.8 + e}{32 + e}
\]

\(\therefore e = 1.6 \text{ cm.}\)

\(v = 256 \times 4 (32 + 1.6) \times 10^{-2}\)

Velocity of sound in air \(v = 344 \text{ m s}^{-1}\)

7.15 A railway engine and a car are moving parallel but in opposite direction with velocities 144 km/hr and 72 km/hr respectively. The frequency of engine’s whistle is 500 Hz and the velocity of sound is 340 m s\(^{-1}\). Calculate the frequency of sound heard in the car when (i) the car and engine are approaching each other (ii) both are moving away from each other.

**Data :** The velocity of source \(v_S = 144 \text{ km/hr}\) and the velocity of observer \(v_o = 72 \text{ km/hr}\)

\(v = 340 \text{ m s}^{-1}, n = 500 \text{ Hz}\)

**Solution :** (i) When the car and engine approaches each other

\[n' = \left(\frac{v + v_o}{v - v_S}\right)n\]

\[v_S = \frac{144 \times 10^3}{60 \times 60} = 40 \text{ m s}^{-1}\]

\[v_o = \frac{72 \times 10^3}{60 \times 60} = 20 \text{ m s}^{-1}\]

\(\therefore n' = \frac{340 + 20}{340 - 40} \times 500\)

The frequency of sound heard is = 600 Hz

(ii) When the car and engine are moving away from each other

\[n'' = \left(\frac{v - v_o}{v + v_S}\right)n\]

\[= \frac{340 - 20}{340 + 40} \times 500\]

The frequency of sound heard is = 421 Hz
Self evaluation

(The questions and problems given in this self evaluation are only samples. In the same way any question and problem could be framed from the text matter. Students must be prepared to answer any question and problem from the text matter, not only from the self evaluation.)

7.1 In a longitudinal wave there is state of maximum compression at a point at an instant. The frequency of wave is 50 Hz. After what time will the same point be in the state of maximum rarefaction.
(a) 0.01 s     (b) 0.002 s
(c) 25 s       (d) 50 s

7.2 Sound of frequency 256 Hz passes through a medium. The maximum displacement is 0.1 m. The maximum velocity is equal to
(a) $60\pi$ m s$^{-1}$  (b) $51.2\pi$ m s$^{-1}$
(c) 256 m s$^{-1}$       (d) 512 m s$^{-1}$

7.3 Which of the following does not affect the velocity of sound?
(a) temperature of the gas  (b) pressure of the gas
(c) mass of the gas          (d) specific heat capacities of the gas

7.4 When a wave passes from one medium to another, there is change of
(a) frequency and velocity  (b) frequency and wavelength
(c) wavelength and velocity (d) frequency, wavelength and velocity

7.5 Sound waves from a point source are propagating in all directions. What will be the ratio of amplitude at a distance 9 m and 25 m from the source?
(a) 25:9     (b) 9: 25
(c) 3 : 5    (d) 81 : 625

7.6 The intensity level of two sounds are 100 dB and 50 dB. Their ratio of intensities are
(a) $10^1$     (b) $10^5$
(c) $10^3$     (d) $10^{10}$
7.7 Number of beats produced by two waves of \( y_1 = a \sin 2000 \pi t \), \( y_2 = a \sin 2008 \pi t \) is
(a) 0 (b) 1
(c) 4 (d) 8

7.8 In order to increase the fundamental frequency of a stretched string from 100 Hz to 400 Hz, the tension must be increased by
(a) 2 times (b) 4 times
(c) 8 times (d) 16 times

7.9 The second overtone of an open pipe has the same frequency as the first overtone of a closed pipe of 2 m long. The length of the open pipe is
(a) 2 m (b) 4 m
(c) 0.5 m (d) 0.75 m

7.10 A source of sound of frequency 150 Hz is moving in a direction towards an observer with a velocity 110 m s\(^{-1}\). If the velocity of sound is 330 m s\(^{-1}\), the frequency of sound heard by the person is
(a) 225 Hz (b) 200 Hz
(c) 150 Hz (d) 100 Hz

7.11 Define wave motion. Mention the properties of the medium in which a wave propagates.

7.12 What are the important characteristics of wave motion?

7.13 Distinguish between transverse and longitudinal waves.

7.14 In solids both longitudinal and transverse waves are possible, but transverse waves are not produced in gases. Why?

7.15 Define the terms wavelength and frequency in wave motion. Prove that \( v = n \lambda \).

7.16 Obtain an expression for the velocity of transverse wave in a stretched string, when it is vibrating in fundamental mode.

7.17 Derive Newton - Laplace formula for the velocity of sound in gases.

7.18 Show that the velocity of sound increases by 0.61 m s\(^{-1}\) for every degree rise of temperature.
7.19 Sound travels faster on rainy days. Why?
7.20 Obtain the equation for plane progressive wave.
7.21 Distinguish between intensity and loudness of sound.
7.22 What do you understand by decibel?
7.23 On what factors does the intensity of sound depend?
7.24 What is an echo? Why an echo cannot be heard in a small room?
7.25 Write a short note on whispering gallery.
7.26 State the principle of superposition.
7.27 What are the essential conditions for the formation of beats?
7.28 What are beats? Show that the number of beats produced per second is equal to the difference in frequencies.
7.29 What is interference of sound waves? Describe an experiment to explain the phenomenon of interference of waves.
7.30 How are stationary waves formed?
7.31 Derive the equation of stationary wave and deduce the condition for nodes and antinodes.
7.32 What are the properties of stationary waves?
7.33 State the laws of transverse vibrations in stretched strings.
7.34 List out the differences between a progressive wave and a stationary wave.
7.35 What are overtones and harmonics?
7.36 Why open organ pipes are preferred for making flute?
7.37 Prove that in a pipe closed at one end, frequency of harmonics are in the ratio 1:3:5.
7.38 Explain how overtones are produced in an open pipe. Show that all harmonics are present in the open pipe.
7.39 What is meant by end correction?
7.40 What is doppler effect? Derive the formula for the change in frequency (f) when the source is approaching and receding from
the observer and (ii) when the source is stationary and observer is moving towards and away from the source.

**Problems**

7.41 A wave of length 0.60 cm is produced in air and travels with a velocity of 340 m s\(^{-1}\). Will it be audiable to human ear?

7.42 The velocity of sound in water is 1480 m s\(^{-1}\). Find the frequency of sound wave such that its wavelength in water is the same as the wavelength in air of a sound wave of frequency 1000 Hz. (The velocity of sound in air is 340 m s\(^{-1}\)).

7.43 Calculate the ratio of velocity of sound in hydrogen \(\gamma = \frac{7}{5}\) to that in helium \(\gamma = \frac{5}{3}\) at the same temperature.

7.44 The equation of a progressive wave travelling along the x axis is given by \(y = 10 \sin \pi (2t - 0.01x)\) where \(y\) and \(x\) are in m and \(t\) in s. Calculate (i) amplitude (ii) frequency and wavelength (iii) wave velocity.

7.45 If the intensity is increased by a factor 60, by how many decibels the sound level is increased.

7.46 Two sound waves, originating from the same source, travel along different paths in air and then meet at a point. If the source vibrates at a frequency of 1.0 kHz and one path is 83 cm longer than the other, what will be the nature of interference? The speed of sound in air is 332 m s\(^{-1}\).

7.47 In an experiment, the tuning fork and sonometer give 5 beats per second, when their lengths are 1m and 1.05m respectively. Calculate the frequency of the fork.

7.48 A steel wire of length 1.2 m with a tension of 9.8 N is found to resonate in five segments at a frequency of 240 Hz. Find the mass of the string.

7.49 How can a stretched string of length 114 cm be divided into three segments so that the fundamental frequency of the three segments be in the ratio of 1 : 3 : 4.
7.50 An open organ pipe has a fundamental frequency of 240 Hz. The first overtone of a closed organ pipe has the same frequency as the first overtone of the open pipe. How long is each pipe? Velocity of sound at room temperature is 350 ms\(^{-1}\).

7.51 A tuning fork of frequency 800 Hz produces resonance in a resonance column apparatus. If successive resonances are produced at lengths 9.75 cm and 31.25 cm, calculate the velocity of sound in air.

7.52 A train standing at a signal of a railway station blows a whistle of frequency 256 Hz in air. Calculate the frequency of the sound as heard by a person standing on the platform when the train (i) approaches the platform with a speed of 40 m s\(^{-1}\) (ii) recedes from the platform with the same speed.

7.53 A whistle of frequency 480 Hz rotates in a circle of radius 1.25 m at an angular speed of 16.0 rad s\(^{-1}\). What is the lowest and highest frequency heard by a listener a long distance away at rest with respect to the centre of the circle. The velocity of sound is 340 m s\(^{-1}\).

7.54 Two tuning forks A and B when sounded together give 4 beats per second. The fork A is in resonance with a closed column of air of length 15 cm, while the second is in resonance with an open column of length 30.5 cm. Calculate their frequencies.
Answers

7.1 (a) 7.2 (b) 7.3 (b)
7.4 (c) 7.5 (a) 7.6 (b)
7.7 (c) 7.8 (d) 7.9 (b)
7.10 (a)

7.41 $5.666 \times 10^4$ Hz, not audible
7.42 4353 Hz
7.43 1.833
7.44 10 m, 1 Hz, 200 m, 200 ms$^{-1}$
7.45 18 dB
7.46 destructive interference, as this is odd multiple of $\pi$
7.47 205 Hz
7.48 $7.38 \times 10^{-4}$ kg
7.49 $72 \times 10^{-2}$ m, $24 \times 10^{-2}$ m, $18 \times 10^{-2}$ m
7.50 $54.7 \times 10^{-2}$ m, $72.9 \times 10^{-2}$ m
7.51 344 m s$^{-1}$
7.52 290 Hz, 229 Hz
7.53 510 Hz, 453 Hz
7.54 240 Hz, 244 Hz
8. Heat and Thermodynamics

In early days, according to caloric theory of heat, heat was regarded as an invisible and weightless fluid called “caloric”. The two bodies at different temperatures placed in contact attain thermal equilibrium by the exchange of caloric. The caloric flows from the hot body to the cold body, till their temperature becomes equal. However, this theory failed to explain the production of heat due to friction in the experiments conducted by Court Rumford. Rubbing our hands against each other produces heat. Joule’s paddle wheel experiment led to the production of heat by friction. These observations led to the dynamic theory of heat, according to which heat is a form of energy called thermal energy.

Every body is made up of molecules. Depending on its nature and temperature, the molecules may possess translatory motion, vibratory motion and rotatory motion about its axis. Each type of motion provides some kinetic energy to the molecules. Heat possessed by a body is the total thermal energy of the body, which is the sum of kinetic energies of all the individual molecules of the body.

Temperature of a body is the degree of hotness or coldness of the body. Heat flows from a body at high temperature to a body at low temperature when they are in contact with each other. Modern concept of temperature follows from zeroth law of thermodynamics. Temperature is the thermal state of the body, that decides the direction of flow of heat. Temperature is now regarded as one of the fundamental quantities.

8.1 Kinetic theory of gases

The founder of modern kinetic theory of heat by common consent is Daniel Bernoulli. But the credit for having established it on a firm mathematical basis is due to Clausius and Maxwell in whose hands it attained the present form.

8.1.1 Postulates of Kinetic theory of gases

(1) A gas consists of a very large number of molecules. Each one is a perfectly identical elastic sphere.
(2) The molecules of a gas are in a state of continuous and random motion. They move in all directions with all possible velocities.

(3) The size of each molecule is very small as compared to the distance between them. Hence, the volume occupied by the molecule is negligible in comparison to the volume of the gas.

(4) There is no force of attraction or repulsion between the molecules and the walls of the container.

(5) The collisions of the molecules among themselves and with the walls of the container are perfectly elastic. Therefore, momentum and kinetic energy of the molecules are conserved during collisions.

(6) A molecule moves along a straight line between two successive collisions and the average distance travelled between two successive collisions is called the mean free path of the molecules.

(7) The collisions are almost instantaneous (i.e.) the time of collision of two molecules is negligible as compared to the time interval between two successive collisions.

**Avogadro number**

*Avogadro number is defined as the number of molecules present in one mole of a substance.* It is constant for all the substances. Its value is $6.023 \times 10^{23}$.

**8.1.2 Pressure exerted by a gas**

The molecules of a gas are in a state of random motion. They continuously collide against the walls of the container. During each collision, momentum is transferred to the walls of the container. The pressure exerted by the gas is due to the continuous collision of the molecules against the walls of the container. Due to this continuous collision, the walls experience a continuous force which is equal to the total momentum imparted to the walls per second. The force experienced per unit area of the walls of the container determines the pressure exerted by the gas.

Consider a cubic container of side $l$ containing $n$ molecules of perfect gas

![Fig. 8.1 Pressure exerted by a gas](image-url)
moving with velocities \( C_1, C_2, C_3 \ldots C_n \) (Fig. 8.1). A molecule moving with a velocity \( C_1 \), will have velocities \( u_1, v_1 \) and \( w_1 \) as components along the \( x, y \) and \( z \) axes respectively. Similarly \( u_2, v_2 \) and \( w_2 \) are the velocity components of the second molecule and so on. Let a molecule \( P \) (Fig. 8.2) having velocity \( C_1 \) collide against the wall marked I (BCFG) perpendicular to the \( x \)-axis. Only the \( x \)-component of the velocity of the molecule is relevant for the wall I. Hence momentum of the molecule before collision is \( \mu u_1 \) where \( \mu \) is the mass of the molecule. Since the collision is elastic, the molecule will rebound with the velocity \( u_1 \) in the opposite direction. Hence momentum of the molecule after collision is \(-\mu u_1\).

Change in the momentum of the molecule

\[
\text{Change in the momentum} = \text{Final momentum} - \text{Initial momentum} = -\mu u_1 - \mu u_1 = -2\mu u_1
\]

During each successive collision on face I the molecule must travel a distance \( 2l \) from face I to face II and back to face I.

Time taken between two successive collisions is \( \frac{2l}{u_1} \)

\[\therefore \text{Rate of change of momentum} = \frac{-2\mu u_1}{2l} = \frac{-2\mu u_1^2}{l} = \frac{-\mu u_1^2}{l} \]

(i.e) Force exerted on the molecule \( = \frac{-\mu u_1^2}{l} \)

\[\therefore \text{According to Newton’s third law of motion, the force exerted by the molecule} \]

\[= -\frac{(-\mu u_1^2)}{l} = \frac{\mu u_1^2}{l} \]
Force exerted by all the \( n \) molecules is
\[
F_x = \frac{m u_1^2}{l} + \frac{m u_2^2}{l} + \ldots + \frac{m u_n^2}{l}
\]

Pressure exerted by the molecules
\[
P_x = \frac{F_x}{A}
\]
\[
= \frac{1}{l^2} \left( \frac{m u_1^2}{l} + \frac{m u_2^2}{l} + \ldots + \frac{m u_n^2}{l} \right)
\]
\[
= \frac{m}{l^3} \left( u_1^2 + u_2^2 + \ldots + u_n^2 \right)
\]

Similarly, pressure exerted by the molecules along Y and Z axes are
\[
P_y = \frac{m}{l^3} \left( v_1^2 + v_2^2 + \ldots + v_n^2 \right)
\]
\[
P_z = \frac{m}{l^3} \left( w_1^2 + w_2^2 + \ldots + w_n^2 \right)
\]

Since the gas exerts the same pressure on all the walls of the container
\[
P_x = P_y = P_z = P
\]
\[
P = \frac{P_x + P_y + P_z}{3}
\]
\[
P = \frac{1}{3} \frac{m}{l^3} \left[ (u_1^2 + u_2^2 + \ldots + u_n^2) + (v_1^2 + v_2^2 + \ldots + v_n^2) + (w_1^2 + w_2^2 + \ldots + w_n^2) \right]
\]
\[
P = \frac{1}{3} \frac{m}{l^3} \left[ (u_1^2 + v_2^2 + w_n^2) + (u_2^2 + v_2^2 + w_n^2) + \ldots + (u_n^2 + v_n^2 + w_n^2) \right]
\]
\[
P = \frac{1}{3} \frac{m}{l^3} \left[ C_1^2 + C_2^2 + \ldots + C_n^2 \right]
\]

where \( C_1^2 = (u_1^2 + u_2^2 + \ldots + u_n^2) \)
\[
P = \frac{1}{3} \frac{mn}{l^3} \left[ \frac{C_1^2 + C_2^2 + \ldots + C_n^2}{n} \right]
\]
\[
P = \frac{1}{3} \frac{mn}{V} \cdot C^2
\]

where \( C \) is called the root mean square (RMS) velocity, which is defined as the square root of the mean value of the squares of velocities of individual molecules.
\[ C = \sqrt{\frac{C_1^2 + C_2^2 + \ldots + C_n^2}{n}} \]

8.1.3 **Relation between the pressure exerted by a gas and the mean kinetic energy of translation per unit volume of the gas**

Pressure exerted by unit volume of a gas, \( P = \frac{1}{3} \rho C^2 \)

\( P = \frac{1}{3} \rho C^2 \) (\( \because \) \( \rho \) = mass per unit volume of the gas; \( \rho \) = density of the gas)

Mean kinetic energy of translation per unit volume of the gas

\[ E = \frac{1}{2} \rho C^2 \]

\[ \frac{P}{E} = \frac{1}{3} \rho C^2 \]

\[ \frac{1}{2} \rho C^2 = \frac{2}{3} \rho C^2 \]

\[ P = \frac{2}{3} E \]

8.1.4 **Average kinetic energy per molecule of the gas**

Let us consider one mole of gas of mass \( M \) and volume \( V \).

\[ P = \frac{1}{3} \rho C^2 \]

\[ P = \frac{1}{3} \frac{M}{V} C^2 \]

\[ PV = \frac{1}{3} MC^2 \]

From gas equation

\[ PV = RT \]

\[ \therefore RT = \frac{1}{3} MC^2 \]

\[ \frac{3}{2} RT = \frac{1}{2} MC^2 \]

(i.e) Average kinetic energy of one mole of the gas is equal to \( \frac{3}{2} \) RT

Since one mole of the gas contains \( N \) number of atoms where \( N \) is the Avogadro number
we have \( M = Nm \)

\[ \therefore \frac{1}{2} mNC^2 = \frac{3}{2} RT \]

\[ \frac{1}{2} mC^2 = \frac{3}{2} \frac{R}{N} T \]

\[ = \frac{3}{2} kT \quad \text{where} \quad k = \frac{R}{N}, \text{is the Boltzmann constant} \]

Its value is \( 1.38 \times 10^{-23} \text{ J K}^{-1} \)

\[ \therefore \text{Average kinetic energy per molecule of the gas is equal to} \quad \frac{3}{2} kT \]

Hence, it is clear that the temperature of a gas is the measure of the mean translational kinetic energy per molecule of the gas.

### 8.2 Degrees of freedom

The number of degrees of freedom of a dynamical system is defined as the total number of co-ordinates or independent variables required to describe the position and configuration of the system.

**For translatory motion**

(i) A particle moving in a straight line along any one of the axes has one degree of freedom (e.g) Bob of an oscillating simple pendulum.

(ii) A particle moving in a plane (X and Y axes) has two degrees of freedom. (eg) An ant that moves on a floor.

(iii) A particle moving in space (X, Y and Z axes) has three degrees of freedom. (eg) a bird that flies.

A point mass cannot undergo rotation, but only translatory motion. A rigid body with finite mass has both rotatory and translatory motion. The rotatory motion also can have three co-ordinates in space, like translatory motion ; Therefore a rigid body will have six degrees of freedom ; three due to translatory motion and three due to rotatory motion.

#### 8.2.1 Monoatomic molecule

Since a monoatomic molecule consists of only a single atom of point mass it has three degrees of freedom of translatory motion along the three co-ordinate axes as shown in Fig. 8.3.
Examples: molecules of rare gases like helium, argon, etc.

8.2.2 Diatomic molecule

The diatomic molecule can rotate about any axis at right angles to its own axis. Hence it has two degrees of freedom of rotational motion in addition to three degrees of freedom of translational motion along the three axes. So, a diatomic molecule has five degrees of freedom (Fig. 8.4). Examples: molecules of \( \text{O}_2, \text{N}_2, \text{CO}, \text{Cl}_2 \), etc.

8.2.3 Triatomic molecule (Linear type)

In the case of triatomic molecule of linear type, the centre of mass lies at the central atom. It, therefore, behaves like a diatomic molecule with three degrees of freedom of translation and two degrees of freedom of rotation, totally it has five degrees of freedom (Fig. 8.5). Examples: molecules of \( \text{CO}_2, \text{CS}_2 \), etc.

8.2.4 Triatomic molecule (Non-linear type)

A triatomic non-linear molecule may rotate, about the three mutually perpendicular axes, as shown in Fig. 8.6. Therefore, it possesses three degrees of freedom of rotation in addition to three degrees of freedom of translation along the three co-ordinate axes. Hence it has six degrees of freedom. Examples: molecules of \( \text{H}_2\text{O}, \text{SO}_2 \), etc.

In all the above cases, only the translatory and rotatory motion of the molecules have been considered. The vibratory motion of the molecules has not been taken into consideration.
8.3 Law of equipartition of energy

Law of equipartition of energy states that for a dynamical system in thermal equilibrium the total energy of the system is shared equally by all the degrees of freedom. The energy associated with each degree of freedom per molecule is \( \frac{1}{2} kT \), where \( k \) is the Boltzmann’s constant.

Let us consider one mole of a monoatomic gas in thermal equilibrium at temperature \( T \). Each molecule has 3 degrees of freedom due to translatory motion. According to kinetic theory of gases, the mean kinetic energy of a molecule is \( \frac{3}{2} kT \).

Let \( C_x \), \( C_y \) and \( C_z \) be the components of RMS velocity of a molecule along the three axes. Then the average energy of a gas molecule is given by

\[
\frac{1}{2} mC^2 = \frac{1}{2} mC^2_x + \frac{1}{2} mC^2_y + \frac{1}{2} mC^2_z
\]

\[
\therefore \frac{1}{2} mC^2_x + \frac{1}{2} mC^2_y + \frac{1}{2} mC^2_z = \frac{3}{2} kT
\]

Since molecules move at random, the average kinetic energy corresponding to each degree of freedom is the same.

\[
\therefore \frac{1}{2} mC^2_x = \frac{1}{2} mC^2_y = \frac{1}{2} mC^2_z
\]

(i.e) \[
\frac{1}{2} mC^2_x = \frac{1}{2} mC^2_y = \frac{1}{2} mC^2_z = \frac{1}{2} kT
\]

\[
\therefore \text{Mean kinetic energy per molecule per degree of freedom is } \frac{1}{2} kT.
\]

8.4. Thermal equilibrium

Let us consider a system requiring a pair of independent co-ordinates \( X \) and \( Y \) for their complete description. If the values of \( X \) and \( Y \) remain unchanged so long as the external factors like temperature also remains the same, then the system is said to be in a state of thermal equilibrium.

Two systems \( A \) and \( B \) having their thermodynamic co-ordinates \( X \) and \( Y \) and \( X_1 \) and \( Y_1 \) respectively separated from each other, for example, by a wall, will have new and common co-ordinates \( X \) and \( Y \).
spontaneously, if the wall is removed. Now the two systems are said to be in thermal equilibrium with each other.

8.4.1 Zeroth law of thermodynamics

If two systems A and B are separately in thermal equilibrium with a third system C, then the three systems are in thermal equilibrium with each other. Zeroth law of thermodynamics states that two systems which are individually in thermal equilibrium with a third one, are also in thermal equilibrium with each other.

This Zeroth law was stated by Flower much later than both first and second laws of thermodynamics.

This law helps us to define temperature in a more rigorous manner.

8.4.2 Temperature

If we have a number of gaseous systems, whose different states are represented by their volumes and pressures $V_1, V_2, V_3 \ldots$ and $P_1, P_2, P_3 \ldots$ etc., in thermal equilibrium with one another, we will have $\phi_1 (P_1, V_1) = \phi_2 (P_2, V_2) = \phi_3 (P_3, V_3)$ and so on, where $\phi$ is a function of $P$ and $V$. Hence, despite their different parameters of $P$ and $V$, the numerical value of the these functions or the temperature of these systems is same.

Temperature may be defined as the particular property which determines whether a system is in thermal equilibrium or not with its neighbouring system when they are brought into contact.

8.5 Specific heat capacity

Specific heat capacity of a substance is defined as the quantity of heat required to raise the temperature of 1 kg of the substance through 1K. Its unit is $J \, kg^{-1} \, K^{-1}$.

Molar specific heat capacity of a gas

Molar specific heat capacity of a gas is defined as the quantity of heat required to raise the temperature of 1 mole of the gas through 1K. Its unit is $J \, mol^{-1} \, K^{-1}$.

Specific heat capacity of a gas may have any value between $-\infty$ and $+\infty$ depending upon the way in which heat energy is given.
Let $m$ be the mass of a gas and $C$ its specific heat capacity. Then
\[ \Delta Q = m \times C \times \Delta T \] where $\Delta Q$ is the amount of heat absorbed and $\Delta T$ is the corresponding rise in temperature.

(i.e) $C = \frac{\Delta Q}{m \Delta T}$

**Case (i)**

If the gas is insulated from its surroundings and is suddenly compressed, it will be heated up and there is rise in temperature, even though no heat is supplied from outside

(i.e) $\Delta Q = 0$

$\therefore C = 0$

**Case (ii)**

If the gas is allowed to expand slowly, in order to keep the temperature constant, an amount of heat $\Delta Q$ is supplied from outside,

then $C = \frac{\Delta Q}{m \times \Delta T} = \frac{\Delta Q}{0} = +\infty$

($\therefore$ $\Delta Q$ is +ve as heat is supplied from outside)

**Case (iii)**

If the gas is compressed gradually and the heat generated $\Delta Q$ is conducted away so that temperature remains constant, then

$C = \frac{\Delta Q}{m \times \Delta T} = \frac{-\Delta Q}{0} = -\infty$

($\therefore$ $\Delta Q$ is -ve as heat is supplied by the system)

Thus we find that if the external conditions are not controlled, the value of the specific heat capacity of a gas may vary from $+\infty$ to $-\infty$.

Hence, in order to find the value of specific heat capacity of a gas, either the pressure or the volume of the gas should be kept constant. Consequently a gas has two specific heat capacities (i) Specific heat capacity at constant volume (ii) Specific heat capacity at constant pressure.

**Molar specific heat capacity of a gas at constant volume**

Molar specific heat capacity of a gas at constant volume $C_v$ is defined as the quantity of heat required to raise the temperature of one mole of a gas through 1 K, keeping its volume constant.
**Molar specific heat capacity of a gas at constant pressure**

Molar specific heat capacity of a gas at constant pressure $C_p$ is defined as the quantity of heat to raise the temperature of one mole of a gas through 1 K keeping its pressure constant.

**Specific heat capacity of monoatomic, diatomic and triatomic gases**

Monoatomic gases like argon, helium etc. have three degrees of freedom.

We know, kinetic energy per molecule, per degree of freedom is $\frac{1}{2} kT$.

∴ Kinetic energy per molecule with three degrees of freedom is $\frac{3}{2} kT$.

Total kinetic energy of one mole of the monoatomic gas is given by $E = \frac{3}{2} kT \times N = \frac{3}{2} RT$, where $N$ is the Avogadro number.

∴ $\frac{dE}{dT} = \frac{3}{2} R$

If $dE$ is a small amount of heat required to raise the temperature of 1 mole of the gas at constant volume, through a temperature $dT$,

$dE = 1 \times C_V \times dT$

$C_V = \frac{dE}{dT} = \frac{3}{2} R$

As $R = 8.31$ J mol$^{-1}$ K$^{-1}$

$C_V = \frac{3}{2} \times 8.31 = 12.465$ J mol$^{-1}$ K$^{-1}$

Then $C_p - C_V = R$

$C_p = C_V + R$

$= \frac{3}{2} R + R = \frac{5}{2} R = \frac{5}{2} \times 8.31$

∴ $C_p = 20.775$ J mol$^{-1}$ K$^{-1}$
In diatomic gases like hydrogen, oxygen, nitrogen etc., a molecule has five degrees of freedom. Hence the total energy associated with one mole of diatomic gas is

\[ E = 5 \times \frac{1}{2} kT \times N = \frac{5}{2} RT \]

Also,

\[ C_v = \frac{dE}{dT} = \frac{d}{dT} \left( \frac{5}{2} RT \right) = \frac{5}{2} R \]

\[ C_v = \frac{5}{2} \times 8.31 = 20.775 \text{ J mol}^{-1} \text{ K}^{-1} \]

But \[ C_p = C_v + R \]

\[ = \frac{5}{2} R + R = \frac{7}{2} R \]

\[ C_p = \frac{7}{2} \times 8.31 \]

\[ = 29.085 \text{ J mol}^{-1} \text{ K}^{-1} \]

Similarly, \( C_p \) and \( C_v \) can be calculated for triatomic gases.

**Internal energy**

Internal energy \( U \) of a system is the energy possessed by the system due to molecular motion and molecular configuration. The internal kinetic energy \( U_K \) of the molecules is due to the molecular motion and the internal potential energy \( U_P \) is due to molecular configuration. Thus

\[ U = U_K + U_P \]

It depends only on the initial and final states of the system. In case of an ideal gas, it is assumed that the intermolecular forces are zero. Therefore, no work is done, although there is change in the intermolecular distance. Thus \( U_P = 0 \). Hence, internal energy of an ideal gas has only internal kinetic energy, which depends only on the temperature of the gas.

In a real gas, intermolecular forces are not zero. Therefore, a definite amount of work has to be done in changing the distance between the molecules. Thus the internal energy of a real gas is the sum of internal kinetic energy and internal potential energy. Hence, it would depend upon both the temperature and the volume of the gas.
8.6 First law of thermodynamics

Let us consider a gas inside a cylinder fitted with a movable frictionless piston. The walls of the cylinder are made up of non-conducting material and the bottom is made up of conducting material (Fig. 8.7).

Let the bottom of the cylinder be brought in contact with a hot body like burner. The entire heat energy given to the gas is not converted into work. A part of the heat energy is used up in increasing the temperature of the gas (i.e.) in increasing its internal energy and the remaining energy is used up in pushing the piston upwards (i.e.) in doing work.

If $\Delta Q$ is the heat energy supplied to the gas, $U_1$ and $U_2$ are initial and final internal energies and $\Delta W$ is the work done by the system, then

\[ \Delta Q = \Delta W + (U_2 - U_1) \]
\[ \Delta Q = \Delta W + \Delta U \]

where $\Delta U$ is the change in the internal energy of the system.

Hence, the first law of thermodynamics states that the amount of heat energy supplied to a system is equal to the sum of the change in internal energy of the system and the work done by the system. This law is in accordance with the law of conservation of energy.

8.7 Relation between $C_p$ and $C_v$ (Meyer’s relation)

Let us consider one mole of an ideal gas enclosed in a cylinder provided with a frictionless piston of area $A$. Let $P$, $V$ and $T$ be the pressure, volume and absolute temperature of gas respectively (Fig. 8.8).

A quantity of heat $dQ$ is supplied to the gas. To keep the volume of the gas constant, a small weight is placed over the piston. The pressure and the temperature of the gas increase to $P + dP$ and $T + dT$ respectively. This heat energy $dQ$ is used to increase the internal energy $dU$ of the gas. But the gas does not do any work ($dW = 0$).

\[ \therefore dQ = dU = 1 \times C_v \times dT \quad \ldots \quad (1) \]

The additional weight is now removed from the piston. The piston now moves upwards through a distance $dx$, such that the pressure of
the enclosed gas is equal to the atmospheric pressure \( P \). The temperature of the gas decreases due to the expansion of the gas.

Now a quantity of heat \( dq' \) is supplied to the gas till its temperature becomes \( T + \Delta T \). This heat energy is not only used to increase the internal energy \( dU \) of the gas but also to do external work \( dW \) in moving the piston upwards.

\[
\therefore dq' = dU + dW
\]

Since the expansion takes place at constant pressure,

\[
dq' = C_p dT
\]

\[
\therefore C_p dT = C_v dT + dW \quad \ldots (2)
\]

Work done, \( dW = \text{force} \times \text{distance} \)

\[
dW = P \times A \times dx
\]

\[
dW = P \ dV \ (\text{since } A \times dx = dV, \text{change in volume})
\]

\[
\therefore C_p dT = C_v dT + P \ dV \quad \ldots (3)
\]

The equation of state of an ideal gas is

\[
PV = RT
\]

Differentiating both the sides

\[
PdV = RdT \quad \ldots (4)
\]

Substituting equation (4) in (3).

\[
C_p dT = C_v dT + RdT
\]

\[
C_p = C_v + R
\]
This equation is known as Meyer’s relation.

8.8 Indicator diagram (P-V diagram)

A curve showing variation of volume of a substance taken along the X-axis and the variation of pressure taken along Y-axis is called an indicator diagram or P-V diagram. The shape of the indicator diagram shall depend on the nature of the thermodynamical process the system undergoes.

Let us consider one mole of an ideal gas enclosed in a cylinder fitted with a perfectly frictionless piston. Let $P_1$, $V_1$ and $T$ be the initial state of the gas. If $dV$ is an infinitesimally small increase in volume of the gas during which the pressure $P$ is assumed to be constant, then small amount of work done by the gas is $dW = PdV$

In the indicator diagram $dW = \text{area } a_1b_1c_1d_1$

∴ The total work done by the gas during expansion from $V_1$ to $V_2$ is

$$W = \int_{V_1}^{V_2} PdV = \text{Area ABCD, in the indicator diagram.}$$

Hence, in an indicator diagram the area under the curve represents the work done (Fig. 8.9).

8.8.1 Isothermal process

When a gas undergoes expansion or compression at constant temperature, the process is called isothermal process.

Let us consider a gas in a cylinder provided with a frictionless piston. The cylinder and the piston are made up of conducting material. If the piston is pushed down slowly, the heat energy produced will be quickly transmitted to the surroundings. Hence, the temperature remains constant but the pressure of the gas increases and its volume decreases.

The equation for an isothermal process is $PV = \text{constant}$.

If a graph is drawn between $P$ and $V$, keeping temperature constant, we get a curve called an isothermal curve. Isotherms for three different
temperatures $T_1$, $T_2$ and $T_3$ are shown in the Fig. 8.10. The curve moves away from the origin at higher temperatures.

During an isothermal change, the specific heat capacity of the gas is infinite.

\[
C = \frac{\Delta Q}{m\Delta T} = \infty \quad \therefore \Delta T = 0
\]

(e.g) Melting of ice at its melting point and vapourisation of water at its boiling point.

### 8.8.2 Work done in an isothermal expansion

Consider one mole of an ideal gas enclosed in a cylinder with perfectly conducting walls and fitted with a perfectly frictionless and conducting piston. Let $P_1$, $V_1$ and $T$ be the initial pressure, volume and temperature of the gas. Let the gas expand to a volume $V_2$ when pressure reduces to $P_2$, at constant temperature $T$. At any instant during expansion let the pressure of the gas be $P$. If $A$ is the area of cross section of the piston, then force

\[
F = P \times A
\]

Let us assume that the pressure of the gas remains constant during an infinitesimally small outward displacement $dx$ of the piston.

Work done

\[
dW = Fdx = PAdx = PdV
\]

Total work done by the gas in expansion from initial volume $V_1$ to final volume $V_2$ is

\[
W = \int_{V_1}^{V_2} P \, dV
\]

We know, $PV = RT$, $P = \frac{RT}{V}$

\[
\therefore W = \int_{V_1}^{V_2} \frac{RT}{V} \, dV = RT \int_{V_1}^{V_2} \frac{1}{V} \, dV
\]

\[
W = RT \left[ \log_e V \right]_{V_1}^{V_2}
\]
\[ W = RT \left[ \log_e V_2 - \log_e V_1 \right] \]
\[ = RT \log_e \frac{V_2}{V_1} \]
\[ W = 2.3026 RT \log_{10} \frac{V_2}{V_1} \]

This is the equation for the work done during an isothermal process.

8.8.3 Adiabatic process

In Greek, adiabatic means "nothing passes through". The process in which pressure, volume and temperature of a system change in such a manner that during the change no heat enters or leaves the system is called adiabatic process. Thus in adiabatic process, the total heat of the system remains constant.

Let us consider a gas in a perfectly thermally insulated cylinder fitted with a piston. If the gas is compressed suddenly by moving the piston downward, heat is produced and hence the temperature of the gas will increase. Such a process is adiabatic compression.

If the gas is suddenly expanded by moving the piston outward, energy required to drive the piston is drawn from the internal energy of the gas, causing fall in temperature. This fall in temperature is not compensated by drawing heat from the surroundings. This is adiabatic expansion.

Both the compression and expansion should be sudden, so that there is no time for the exchange of heat. Hence, in an adiabatic process always there is change in temperature.

Expansion of steam in the cylinder of a steam engine, expansion of hot gases in internal combustion engine, bursting of a cycle tube or car tube, propagation of sound waves in a gas are adiabatic processes.

The adiabatic relation between \(P\) and \(V\) for a gas, is

\[ PV^\gamma = k, \text{ a constant} \quad \text{... (1)} \]

where \(\gamma = \frac{\text{specific heat capacity of the gas at constant pressure}}{\text{specific heat capacity of the gas at constant volume}}\)

From standard gas equation,

\[ PV = RT \]
\[ P = \frac{RT}{V} \]

substituting the value \( P \) in \( (1) \)

\[ \frac{RT}{V} V^\gamma = \text{constant} \]

\[ T V^{\gamma-1} = \text{constant} \]

In an adiabatic process \( Q = \text{constant} \)

\[ \therefore \Delta Q = 0 \]

\[ \therefore \text{specific heat capacity} \ C = \frac{\Delta Q}{m \Delta T} \]

\[ \therefore C = 0 \]

### 8.8.4 Work done in an adiabatic expansion

Consider one mole of an ideal gas enclosed in a cylinder with perfectly non conducting walls and fitted with a perfectly frictionless, non conducting piston.

Let \( P_1, V_1 \) and \( T_1 \) be the initial pressure, volume and temperature of the gas. If \( A \) is the area of cross section of the piston, then force exerted by the gas on the piston is

\[ F = P \times A, \text{ where } P \text{ is pressure of the gas at any instant during expansion.} \]

If we assume that pressure of the gas remains constant during an infinitesimally small outward displacement \( dx \) of the piston,

then work done \( dW = F \times dx = P \times A \times dx \)

\[ dW = P \, dV \]

Total work done by the gas in adiabatic expansion from volume \( V_1 \) to \( V_2 \) is

\[ W = \int_{V_1}^{V_2} P \, dV \]

But \( PV^{\gamma} = \text{constant (k) for adiabatic process} \)

where \( \gamma = \frac{C_p}{C_v} \)

\[ \therefore W = \int_{V_1}^{V_2} k V^{\gamma} \, dV = k \left[ \frac{V^{\gamma+1}}{1+\gamma} \right]_{V_1}^{V_2} \left( \because P = \frac{k}{V^\gamma} \right) \]
\[
W = \frac{k}{1 - \gamma} \left[ V_2^{1/\gamma} - V_1^{1/\gamma} \right]
\]
\[
W = \frac{1}{1 - \gamma} \left[ kV_2^{1/\gamma} - kV_1^{1/\gamma} \right] \quad \ldots (1)
\]

but, \( P_2 V_2^{\gamma} = P_1 V_1^{\gamma} = k \) \quad \ldots (2)

Substituting the value of \( k \) in (1)
\[
\therefore W = \frac{1}{1 - \gamma} \left[ P_2 V_2^{\gamma} - P_1 V_1^{\gamma} \right]
\]
\[
W = \frac{1}{1 - \gamma} \left[ P_2 V_2 - P_1 V_1 \right] \quad \ldots (3)
\]

If \( T_2 \) is the final temperature of the gas in adiabatic expansion, then
\[
P_1 V_1 = RT_1, \ P_2 V_2 = RT_2
\]
Substituting in (3)
\[
W = \frac{1}{1 - \gamma} [RT_2 - RT_1]
\]
\[
W = \frac{R}{1 - \gamma} [T_2 - T_1] \quad \ldots (4)
\]

This is the equation for the work done during adiabatic process.

### 8.9 Reversible and irreversible processes

#### 8.9.1 Reversible process

A thermodynamic process is said to be reversible when (i) the various stages of an operation to which it is subjected can be reversed in the opposite direction and in the reverse order and (ii) in every part of the process, the amount of energy transferred in the form of heat or work is the same in magnitude in either direction. At every stage of the process there is no loss of energy due to friction, inelasticity, resistance, viscosity etc. The heat losses to the surroundings by conduction, convection or radiation are also zero.

**Condition for reversible process**

(i) The process must be infinitely slow.
(ii) The system should remain in thermal equilibrium (i.e) system and surrounding should remain at the same temperature.

**Examples**

(a) Let a gas be compressed isothermally so that heat generated is conducted away to the surrounding. When it is allowed to expand in the same small equal steps, the temperature falls but the system takes up the heat from the surrounding and maintains its temperature.

(b) Electrolysis can be regarded as reversible process, provided there is no internal resistance.

**8.9.2 Irreversible process**

An irreversible process is one which cannot be reversed back. Examples: diffusion of gases and liquids, passage of electric current through a wire, and heat energy lost due to friction. As an irreversible process is generally a very rapid one, temperature adjustments are not possible. Most of the chemical reactions are irreversible.

**8.10 Second law of thermodynamics**

The first law of thermodynamics is a general statement of equivalence between work and heat. The second law of thermodynamics enables us to know whether a process which is allowed by first law of thermodynamics can actually occur or not. The second law of thermodynamics tells about the extent and direction of energy transformation.

Different scientists have stated this law in different ways to bring out its salient features.

**(i) Kelvin’s statement**

Kelvin’s statement of second law is based on his experience about the performance of heat engine.

*It is impossible to obtain a continuous supply of work from a body by cooling it to a temperature below the coldest of its surroundings.*

**(ii) Clausius statement**

*It is impossible for a self acting machine unaided by any external...*
agency to transfer heat from a body at a lower temperature to another body at a higher temperature.

(iii) Kelvin - Planck’s statement

It is impossible to construct a heat engine operating in a cycle, that will extract heat from a reservoir and perform an equivalent amount of work.

8.11 Carnot engine

Heat engine is a device which converts heat energy into mechanical energy.

In the year 1824, Carnot devised an ideal cycle of operation for a heat engine. The machine used for realising this ideal cycle of operation is called an ideal heat engine or carnot heat engine.

The essential parts of a Carnot engine are shown in Fig. 8.11

(i) Source

It is a hot body which is kept at a constant temperature $T_1$. It has infinite thermal capacity. Any amount of heat can be drawn from it at a constant temperature $T_1$ (i.e) its temperature will remain the same even after drawing any amount of heat from it.

(ii) Sink

It is a cold body which is kept at a constant lower temperature $T_2$. Its thermal capacity is also infinite that any amount of heat added to it will not increase its temperature.

(iii) Cylinder

Cylinder is made up of non-conducting walls and conducting bottom. A perfect gas is used as a working substance. The cylinder is fitted with a perfectly non-conducting and frictionless piston.
(iv) Insulating stand

It is made up of non conducting material so as to perform adiabatic operations.

Working: The Carnot engine has the following four stages of operations.


Isothermal expansion

Let us consider one mole of an ideal gas enclosed in the cylinder. Let \( V_1, P_1 \) be the initial volume and pressure of the gas respectively. The initial state of the gas is represented by the point \( A \) on the \( P-V \) diagram. The cylinder is placed over the source which is at the temperature \( T_1 \).

The piston is allowed to move slowly outwards, so that the gas expands. Heat is gained from the source and the process is isothermal at constant temperature \( T_1 \). In this process the volume of the gas changes from \( V_1 \) to \( V_2 \) and the pressure changes from \( P_1 \) to \( P_2 \). This process is represented by \( AB \) in the indicator diagram (Fig. 8.12). During this process, the quantity of heat absorbed from the source is \( Q_1 \) and \( W_1 \) is the corresponding amount of work done by the gas.

\[
Q_1 = W_1 = \int_{V_1}^{V_2} \frac{PdV}{R} = \int_{V_1}^{V_2} \frac{RT_1 \log \left( \frac{V_2}{V_1} \right)}{V} = \text{area ABGEA} \quad \text{...(1)}
\]

Adiabatic expansion

The cylinder is taken from the source and is placed on the insulating stand and the piston is moved further so that the volume of the gas changes from \( V_2 \) to \( V_3 \) and the pressure changes from \( P_2 \) to \( P_3 \). This adiabatic expansion is represented by BC. Since the gas is thermally insulated from all sides no heat can be gained from the surroundings. The temperature of the gas falls from \( T_1 \) to \( T_2' \).
Let $W_2$ be the work done by the gas in expanding adiabatically.

\[ W_2 = \int_{v_2}^{v_1} PdV = \frac{R}{\gamma - 1} (T_1 - T_2) = \text{Area BCHGB} \quad \ldots(2) \]

**Isothermal compression**

The cylinder is now placed on the sink at a temperature $T_2$. The piston is moved slowly downward to compress the gas isothermally. This is represented by CD. Let $(V_4, P_4)$ be the volume and pressure corresponding to the point D. Since the base of the cylinder is conducting the heat produced during compression will pass to the sink so that, the temperature of the gas remains constant at $T_2$. Let $Q_2$ be the amount of heat rejected to the sink and $W_3$ be the amount of work done on the gas in compressing it isothermally.

\[ Q_2 = W_3 = \int_{V_3}^{V_4} -P \, dV = \frac{RT_2}{\gamma - 1} \log \left( \frac{V_4}{V_3} \right) = \text{area CDFHC} \quad \ldots(3) \]

The negative sign indicates that work is done on the working substance.

\[ \therefore Q_2 = RT_2 \log \left( \frac{V_2}{V_4} \right) \]

**Adiabatic compression**

The cylinder is now placed on the insulating stand and the piston is further moved down in such a way that the gas is compressed adiabatically to its initial volume $V_1$ and pressure $P_1$. As the gas is insulated from all sides heat produced raises the temperature of the gas to $T_1$. This change is adiabatic and is represented by DA. Let $W_4$ be the work done on the gas in compressing it adiabatically from a state $D (V_4, P_4)$ to the initial state $A (V_1, P_1)$.

\[ \therefore W_4 = \int_{V_4}^{V_1} -P \, dV = \frac{R}{\gamma - 1} (T_2 - T_1) \]

The negative sign indicates that work is done on the working substance.
∴ \( W_4 = \frac{R}{\gamma - 1} (T_1 - T_2) = \text{Area DAEFD} \) \( \cdots \)(4)

**Work done by the engine per cycle**

Total work done by the gas during one cycle of operation is \( (W_1 + W_2) \).

Total work done on the gas during one cycle of operation is \( (W_3 + W_4) \).

∴ Net work done by the gas in a complete cycle

\[
W = W_1 + W_2 - (W_3 + W_4)
\]

But \( W_2 = W_4 \)

∴ \( W = W_1 - W_3 \)

\[
W = \dot{Q}_1 - \dot{Q}_2
\]

Also, \( W = \text{Area ABGEA} + \text{Area BCHGB} - \text{Area CDFHC} - \text{Area DAEFD} \)

(i.e) \( W = \text{Area ABCDA} \)

Hence in Carnot heat engine, net work done by the gas per cycle is numerically equal to the area of the loop representing the cycle.

**Efficiency of Carnot’s engine**

\[
\eta = \frac{\text{Heat converted into work}}{\text{Heat drawn from the source}} = \frac{\dot{Q}_1 - \dot{Q}_2}{\dot{Q}_1}
\]

\[
\eta = 1 - \frac{\dot{Q}_2}{\dot{Q}_1}
\]

But

\[
\frac{\dot{Q}_1}{\dot{Q}_2} \cdot \frac{W_1}{W_3} = \frac{RT_1 \log \left( \frac{V_2}{V_1} \right)}{RT_2 \log \left( \frac{V_3}{V_4} \right)}
\]

\[
= \frac{T_1 \log \left( \frac{V_2}{V_1} \right)}{T_2 \log \left( \frac{V_3}{V_4} \right)} \cdots \text{(5)}
\]

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Since B and C lie on the same adiabatic curve BC

\[ T_1 V_2^{\gamma - 1} = T_2 V_3^{\gamma - 1} \quad (\therefore \ TV^{\gamma - 1} = \text{constant}) \text{ where } \gamma = \frac{C_p}{C_v} \]

\[ \therefore \frac{T_1}{T_2} = \frac{V_3^{\gamma - 1}}{V_2^{\gamma - 1}} \] \hspace{1cm} ...(6)

Similarly D & A lie on the same adiabatic curve DA

\[ \therefore T_1 V_1^{\gamma - 1} = T_2 V_4^{\gamma - 1} \]

\[ \frac{T_1}{T_2} = \frac{V_4^{\gamma - 1}}{V_1^{\gamma - 1}} \] \hspace{1cm} ...(7)

From (6) & (7)

\[ \frac{V_3}{V_2} = \frac{V_4}{V_1} \quad \text{(or)} \quad \frac{V_2}{V_1} = \frac{V_3}{V_4} \] \hspace{1cm} ...(8)

Substituting equation (8) in equation (5)

\[ \frac{Q_1}{Q_2} = \frac{T_1}{T_2} \ \log \left( \frac{V_3}{V_4} \right) \]

\[ \left( \text{i.e} \right) \frac{Q_2}{Q_1} = \frac{T_2}{T_1} \]

\[ \therefore \text{We have} \quad \eta = 1 - \frac{Q_2}{Q_1} = 1 - \frac{T_2}{T_1} \]

or \[ \eta = \frac{T_1 - T_2}{T_1} \] \hspace{1cm} ...(9)

Inferences

Efficiency of Carnot’s cycle is independent of the working substance, but depends upon the temperatures of heat source and sink.

Efficiency of Carnot’s cycle will be 100% if \( T_1 = \infty \) or \( T_2 = 0 \) K. As neither the temperature of heat source can be made infinite, nor the temperature of the sink can be made 0 K, the inference is that the
Carnot heat engine working on the reversible cycle cannot have 100% efficiency.

8.12 Refrigerator

A refrigerator is a cooling device. An ideal refrigerator can be regarded as Carnot’s heat engine working in the reverse direction. Therefore, it is also called a heat pump. In a refrigerator the working substance would absorb certain quantity of heat from the sink at lower temperature and reject a large amount of heat to the source at a higher temperature with the help of an external agency like an electric motor (Fig. 8.13).

In an actual refrigerator vapours of freon (dichloro difluoro methane $CCl_2F_2$) act as the working substance. Things kept inside the refrigerator act as a sink at a lower temperature $T_2$. A certain amount of work $W$ is performed by the compressor (operated by an electric motor) on the working substance. Therefore, it absorbs heat energy $Q_2$ from the sink and rejects $Q_1$ amount of heat energy to the source (atmosphere) at a temperature $T_1$.

Since this is a reversible cyclic process, the change in the internal energy of the working substance is zero (i.e) $dU = 0$

According to the first law of thermodynamics,

$$dQ = dU + dW$$

But

$$dQ = Q_2 - Q_1$$
$$dW = -W$$

and

$$dU = 0$$

:. $dQ = Q_2 - Q_1 = -W$

Negative sign with $W$ represents work is done on the system

(i.e) $W = Q_1 - Q_2$

Coefficient of performance

Coefficient of performance (COP) is defined as the ratio of quantity
of heat \( Q_2 \) removed per cycle from the contents of the refrigerator to the energy spent per cycle \( W \) to remove this heat.

\[
\text{(i.e) } \text{COP} = \frac{Q_2}{W} = \frac{Q_2}{Q_1 - Q_2} 
\]

\[
\text{(i.e) } \text{COP} = \frac{T_2}{T_1 - T_2} 
\]

The efficiency of the heat engine is

\[
\eta = 1 - \frac{T_2}{T_1}; \quad 1 - \eta = \frac{T_2}{T_1}
\]

\[
\frac{1 - \eta}{\eta} = \frac{T_2}{T_1} \times \frac{T_1}{T_1 - T_2}
\]

\[
\text{(i.e) } \frac{1 - \eta}{\eta} = \frac{T_2}{T_1 - T_2} 
\]

From equations (1) and (2)

\[
\text{COP} = \frac{1 - \eta}{\eta} 
\]

**Inferences**

(i) Equation (1) shows that smaller the value of \((T_1 - T_2)\) greater is the value of COP. (i.e.) smaller is the difference in temperature between atmosphere and the things to be cooled, higher is the COP.

(ii) As the refrigerator works, \( T_2 \) goes on decreasing due to the formation of ice. \( T_1 \) is almost steady. Hence COP decreases. When the refrigerator is defrosted, \( T_2 \) increases.

Therefore defrosting is essential for better working of the refrigerator.

**8.13 Transfer of heat**

There are three ways in which heat energy may get transferred from one place to another. These are conduction, convection and radiation.

**8.13.1 Conduction**

Heat is transmitted through the solids by the process of conduction
only. When one end of the solid is heated, the atoms or molecules of the solid at the hotter end becomes more strongly agitated and start vibrating with greater amplitude. The disturbance is transferred to the neighbouring molecules.

**Applications**

(i) The houses of Eskimos are made up of double walled blocks of ice. Air enclosed in between the double walls prevents transmission of heat from the house to the coldest surroundings.

(ii) Birds often swell their feathers in winter to enclose air between their body and the feathers. Air prevents the loss of heat from the body of the bird to the cold surroundings.

(iii) Ice is packed in gunny bags or sawdust because, air trapped in the saw dust prevents the transfer of heat from the surroundings to the ice. Hence ice does not melt.

**Coefficient of thermal conductivity**

Let us consider a metallic bar of uniform cross section A whose one end is heated. After sometime each section of the bar attains constant temperature but it is different at different sections. This is called steady state. In this state there is no further absorption of heat.

If \( \Delta x \) is the distance between the two sections with a difference in temperature of \( \Delta T \) and \( \Delta Q \) is the amount of heat conducted in a time \( \Delta t \), then it is found that the rate of conduction of heat \( \frac{\Delta Q}{\Delta t} \) is

(i) directly proportional to the area of cross section \( (A) \)

(ii) directly proportional to the temperature difference between the two sections \( (\Delta T) \)

(iii) inversely proportional to the distance between the two sections \( (\Delta x) \).

(i.e) \[
\frac{\Delta Q}{\Delta t} \propto A \frac{\Delta T}{\Delta x}
\]

\[
\frac{\Delta Q}{\Delta t} = KA \frac{\Delta T}{\Delta x}
\]
where $K$ is a constant of proportionality called co-efficient of thermal conductivity of the metal.

\[
\frac{\Delta T}{\Delta x} \text{ is called temperature gradient}
\]

If $A = 1 \ m^2$, and \( \frac{\Delta T}{\Delta x} \) = unit temperature gradient

then, \( \frac{\Delta Q}{\Delta t} = K \times 1 \times 1 \)

or \( K = \frac{\Delta Q}{\Delta t} \)

Coefficient of thermal conductivity of the material of a solid is equal to the rate of flow of heat per unit area per unit temperature gradient across the solid. Its unit is $W \ m^{-1} K^{-1}$.

8.13.2 Convection

It is a phenomenon of transfer of heat in a fluid with the actual movement of the particles of the fluid.

When a fluid is heated, the hot part expands and becomes less dense. It rises and upper colder part replaces it. This again gets heated, rises up replaced by the colder part of the fluid. This process goes on. This mode of heat transfer is different from conduction where energy transfer takes place without the actual movement of the molecules.

Application

It plays an important role in ventilation and in heating and cooling system of the houses.

8.13.3 Radiation

It is the phenomenon of transfer of heat without any material medium. Such a process of heat transfer in which no material medium takes part is known as radiation.

Thermal radiation

The energy emitted by a body in the form of radiation on account of its temperature is called thermal radiation.

It depends on,

(i) temperature of the body,
(ii) nature of the radiating body
The wavelength of thermal radiation ranges from $8 \times 10^{-7}$ m to $4 \times 10^{-4}$ m. They belong to infra-red region of the electromagnetic spectrum.

**Properties of thermal radiations**

1. Thermal radiations can travel through vacuum.
2. They travel along straight lines with the speed of light.
3. They can be reflected and refracted. They exhibit the phenomenon of interference and diffraction.
4. They do not heat the intervening medium through which they pass.
5. They obey inverse square law.

**Absorptive and Emissive power**

**Absorptive power**

Absorptive power of a body for a given wavelength and temperature is defined as the ratio of the radiant energy absorbed per unit area per unit time to the total energy incident on it per unit area per unit time.

It is denoted by $a_\lambda$.

**Emissive power**

Emissive power of a body at a given temperature is the amount of energy emitted per unit time per unit area of the surface for a given wavelength. It is denoted by $e_\lambda$. Its unit is $W \, m^{-2}$.

**8.14 Perfect black body**

A perfect black body is the one which absorbs completely heat radiations of all wavelengths which fall on it and emits heat radiations of all wavelengths when heated. Since a perfect black body neither reflects nor transmits any radiation, the absorptive power of a perfectly black body is unity.

**8.14.1 Fery's black body**

Fery's black body consists of a double walled hollow sphere having a small opening O on one side and a conical projection P just opposite
to it (Fig. 8.14). Its inner surface is coated with lamp black. Any radiation entering the body through the opening O suffers multiple reflections at its innerwall and about 97% of it is absorbed by lamp black at each reflection. Therefore, after a few reflections, almost entire radiation is absorbed. The projection helps in avoiding any direct reflections which even otherwise is not possible because of the small opening O. When this body is placed in a bath at fixed temperature, the heat radiations come out of the hole. The opening O thus acts as a black body radiator.

8.14.2 Prevost’s theory of heat exchanges

Prevost applied the idea of thermal equilibrium to radiation. According to him the rate at which a body radiates or absorbs heat depends on the nature of its surface, its temperature and the temperature of the surroundings. The total amount of heat radiated by a body increases as its temperature rises. A body at a higher temperature radiates more heat energy to the surroundings than it receives from the surroundings. That is why we feel warm when we stand before the furnace.

Similarly a body at a lower temperature receives more heat energy than it loses to the surroundings. That is why we feel cold when we stand before an ice block.

Thus the rise or fall of temperature is due to the exchange of heat radiation. When the temperature of the body is the same as that of surroundings, the exchanges of heat do not stop. In such a case, the amount of heat energy radiated by the body is equal to the amount of heat energy absorbed by it.

A body will stop emitting radiation only when it is at absolute zero. (i.e) 0 K or −273° C. At this temperature the kinetic energy of the molecule is zero.

Therefore, Prevost theory states that all bodies emit thermal radiation at all temperatures above absolute zero, irrespective of the nature of the surroundings.
8.14.3 Kirchoff’s Law

According to this law, the ratio of emissive power to the absorptive power corresponding to a particular wavelength and at a given temperature is always a constant for all bodies. This constant is equal to the emissive power of a perfectly black body at the same temperature and the same wavelength. Thus, if $e_\lambda$ is the emissive power of a body corresponding to a wavelength $\lambda$ at any given temperature, $a_\lambda$ is the absorptive power of the body corresponding to the same wavelength at the same temperature and $E_\lambda$ is the emissive power of a perfectly black body corresponding to the same wavelength and the same temperature, then according to Kirchoff’s law

$$\frac{e_\lambda}{a_\lambda} = \text{constant} = E_\lambda$$

From the above equation it is evident that if $a_\lambda$ is large, then $e_\lambda$ will also be large (i.e) if a body absorbs radiation of certain wavelength strongly then it will also strongly emit the radiation of same wavelength. In other words, good absorbers of heat are good emitters also.

**Applications of Kirchoff’s law**

(i) The silvered surface of a thermos flask is a bad absorber as well as a bad radiator. Hence, ice inside the flask does not melt quickly and hot liquids inside the flask do not cool quickly.

(ii) Sodium vapours on heating, emit two bright yellow lines. These are called $D_1$ and $D_2$ lines of sodium. When continuous white light from carbon arc passes through sodium vapour at low temperature, the continuous spectrum is absorbed at two places corresponding to the wavelengths of $D_1$ and $D_2$ lines and appear as dark lines. This is in accordance with Kirchoff’s law.

8.14.4 Wien’s displacement law

Wien’s displacement law states that the wavelength of the radiation corresponding to the maximum energy ($\lambda_m$) decreases as the temperature $T$ of the body increases.

\[(i.e) \lambda_m T = b \quad \text{where } b \text{ is called Wien’s constant.}\]

Its value is $2.898 \times 10^{-3} \text{ m K}$
8.14.5 Stefan’s law

Stefan’s law states that the total amount of heat energy radiated per second per unit area of a perfect black body is directly proportional to the fourth power of its absolute temperature.

\( E \propto T^4 \) or \( E = \sigma T^4 \)

where \( \sigma \) is called the Stefan’s constant. Its value is \( 5.67 \times 10^{-8} \) W m\(^{-2}\) K\(^{-4}\).

It is also called Stefan - Boltzmann law, as Boltzmann gave a theoretical proof of the result given by Stefan.

8.14.6 Newton’s law of cooling

Newton’s law of cooling states that the rate of cooling of a body is directly proportional to the temperature difference between the body and the surroundings.

The law holds good only for a small difference of temperature. Loss of heat by radiation depends on the nature of the surface and the area of the exposed surface.

Experimental verification of Newton’s law of cooling

Let us consider a spherical calorimeter of mass \( m \) whose outer surface is blackened. It is filled with hot water of mass \( m_1 \). The calorimeter with a thermometer is suspended from a stand (Fig. 8.15).

The calorimeter and the hot water radiate heat energy to the surroundings. Using a stop clock, the temperature is noted for every 30 seconds interval of time till the temperature falls by about 20° C. The readings are entered in a tabular column.

If the temperature of the calorimeter and the water falls from \( T_1 \) to \( T_2 \) in \( t \) seconds, the quantity of heat energy lost by radiation \( Q = (ms + m_1 s_1) (T_1 - T_2) \), where \( s \) is the specific heat capacity of the material of the calorimeter and \( s_1 \) is the specific heat capacity of water.
Rate of cooling = \frac{\text{Heat energy lost}}{\text{time taken}}

\therefore \frac{Q}{t} = \frac{(m_s + m_i)(T_i - T_f)}{t}

If the room temperature is \( T_o \), the average excess temperature of the calorimeter over that of the surroundings is

\( \frac{T_o + T + T_o - T_o}{2} \)

According to Newton's Law of cooling,

\( \frac{Q}{t} \propto \left( \frac{T_i + T_f}{2} \right) \cdot T_o \)

\therefore \frac{(m_s + m_i)(T_i - T_f)}{t} \cdot \left( \frac{T_i + T_f}{2} \cdot T_o \right) = \text{constant}

The time for every 4º fall in temperature is noted. The last column in the tabular column is found to be the same. This proves Newton’s Law of cooling.

**Table 8.1 Newton’s law of cooling**

<table>
<thead>
<tr>
<th>Temperature range</th>
<th>Time ( t ) for every 4º fall of temperature</th>
<th>Average excess of temperature ( \left( \frac{T_i + T_f}{2} \cdot T_o \right) )</th>
<th>( \left( \frac{T_i + T_f}{2} \cdot T_o \right) \cdot t )</th>
</tr>
</thead>
</table>

A cooling curve is drawn by taking time along X-axis and temperature along Y-axis (Fig. 8.16).

From the cooling curve, the rate of fall of temperature at \( T \) is

\( \frac{dT}{dt} = \frac{AB}{BC} \)

![Fig. 8.16 Cooling curve](image-url)
The rate of cooling \( \frac{dT}{dt} \) is found to be directly proportional to \( (T - T_o) \). Hence Newton’s law of cooling is verified.

### 8.15 Solar constant

The solar constant is the amount of radiant energy received per second per unit area by a perfect black body on the Earth with its surface perpendicular to the direction of radiation from the sun in the absence of atmosphere. It is denoted by \( S \) and its value is \( 1.388 \times 10^3 \text{ W m}^2 \). Surface temperature of the Sun can be calculated from solar constant.

#### Surface temperature of the Sun

The Sun is a perfect black body of radius \( r \) and surface temperature \( T \). According to Stefan’s law, the energy radiated by the Sun per second per unit area is equal to \( \sigma T^4 \).

Where \( \sigma \) is Stefan’s Constant.

Hence, the total energy radiated per second by the Sun will be given by

\[
E = \text{surface area of the Sun} \times \sigma T^4
\]

\[
E = 4\pi r^2 \sigma T^4 \quad ... (1)
\]

Let us imagine a sphere with Sun at the centre and the distance between the Sun and Earth \( R \) as radius (Fig. 8.17). The heat energy from the Sun will necessarily pass through this surface of the sphere.

If \( S \) is the solar constant, the amount of heat energy that falls on this sphere per unit time is

\[
E = 4\pi R^2 S \quad ... (2)
\]

By definition, equations (1) & (2) are equal.

\[
4\pi^2 \sigma T^4 = 4\pi R^2 S
\]

\[
T^4 = \frac{R^2 S}{r^2 \sigma}
\]

\[
T = \left( \frac{R^2 S}{r^2 \sigma} \right)^{\frac{1}{4}} \quad (i.e.) \quad T = \left( \frac{R}{r} \right)^{\frac{1}{2}} \left( \frac{S}{\sigma} \right)^{\frac{1}{2}}
\]
Knowing the values of \( R, \ r, \ S \) and \( \sigma \) the surface temperature of the Sun can be calculated.

### 8.15.1 Angstrom pyrheliometer

Pyrheliometer is an instrument used to measure the quantity of heat radiation and solar constant.

Pyrheliometer designed by Angstrom is the simplest and most accurate.

Angstrom’s pyrheliometer consists of two identical strips \( S_1 \) and \( S_2 \) of area \( A \). One junction of a thermocouple is connected to \( S_1 \) and the other junction is connected to \( S_2 \). A sensitive galvanometer is connected to the thermo couple.

Strip \( S_2 \) is connected to an external electrical circuit as shown in Fig. 8.18. When both the strips \( S_1 \) and \( S_2 \) are shielded from the solar radiation, galvanometer shows no deflection as both the junctions are at the same temperature. Now strip \( S_1 \) is exposed to the solar radiation and \( S_2 \) is shielded with a cover M. As strip \( S_1 \) receives heat radiations from the sun, its temperature rises and hence the galvanometer shows deflection. Now current is allowed to pass through the strip \( S_2 \) and it is adjusted so that galvanometer shows no deflection. Now, the strips \( S_1 \) and \( S_2 \) are again at the same temperature.

If the quantity of heat radiation that is incident on unit area in unit time on strip \( S_1 \) is \( Q \) and \( \alpha \) its absorption co-efficient, then the amount of heat radiations absorbed by the strip \( S_1 \) in unit time is \( Q\alpha \).

Also, heat produced in unit time in the strip \( S_2 \) is given by \( VI \), where \( V \) is the potential difference and \( I \) is the current flowing through it.

As heat absorbed = heat produced

\[
Q\alpha = VI \quad \text{(or)} \quad Q = \frac{VI}{A\alpha}
\]

Knowing the values of \( V, \ I, \ A \) and \( \alpha \), \( Q \) can be calculated.
Solved Problems

8.1 At what temperature will the RMS velocity of a gas be tripled its value at NTP?

**Solution**: At NTP, \( T_0 = 273 \text{ K} \)

RMS velocity, \( C = \sqrt{\frac{3RT}{M}} \)

\[ C = \sqrt{\frac{3R \times 273}{M}} \]  

... (1)

Suppose at the temperature \( T \), the RMS velocity is tripled, then

\[ 3C = \sqrt{\frac{3RT}{M}} \]  

... (2)

Divide (2) by (1)

\[ \frac{3C}{C} = \frac{\sqrt{\frac{3RT}{M}}}{\sqrt{\frac{3R \times 273}{M}}} \]

\[ 3 = \frac{T}{273} \]

\[ T = 273 \times 9 = 2457 \text{ K} \]

8.2 Calculate the number of degrees of freedom in 15 cm\(^3\) of nitrogen at NTP.

**Solution**: We know 22400 cm\(^3\) of a gas at NTP contains \( 6.02 \times 10^{23} \) molecules.

\[ \therefore \text{The number of molecules in 15 cm}\(^3\) of N\(_2\) at NTP} \]

\[ n = \frac{15}{22400} \times 6.023 \times 10^{23} = 4.033 \times 10^{20} \]

The number degrees of freedom of a diatomic gas molecule at 273 K, is \( f = 5 \)

\[ \therefore \text{Total degrees of freedom of 15 cm}\(^3\) of the gas} = nf \]

\[ \therefore \text{Total degrees of freedom} = 4.033 \times 10^{20} \times 5 = 2.016 \times 10^{21} \]
8.3 A gas is a mixture of 2 moles of oxygen and 4 moles of argon at temperature T. Neglecting vibrational modes, show that the energy of the system is $11 RT$ where $R$ is the universal gas constant.

**Solution:** Since oxygen is a diatomic molecule with 5 degrees of freedom, degrees of freedom of molecules in 2 moles of oxygen $= f_1 = 2 N \times 5 = 10 N$

Since argon is a monatomic molecule, degrees of freedom of molecules in 4 moles of argon $= f_2 = 4 N \times 3 = 12 N$

$\therefore$ Total degrees of freedom of the mixture $= f = f_1 + f_2 = 22 N$

As per the principle of law of equipartition of energy, energy associated with each degree of freedom of a molecule $= \frac{1}{2} kT$

$\therefore$ Total energy of the system $= \frac{1}{2} kT \times 22 N = 11 RT$

8.4 Two Carnot engines A and B are operating in series. The first one A receives heat at 600 K and rejects to a reservoir at temperature $T$. The second engine B receives the heat rejected by A and in turn rejects heat to a reservoir at 150 K. Calculate the temperature $T$ when (i) The work output of both the engines are equal, (ii) The efficiency of both the engines are equal.

**Solution:** (i) When the work outputs are equal:

For the first engine $W_1 = Q_1 - Q_2$

For the second engine $W_2 = Q_2 - Q_3$

Given (i.e) $W_1 = W_2$

$Q_1 - Q_2 = Q_2 - Q_3$

Divide by $Q_2$ on both sides

$\frac{Q_1}{Q_2} - 1 = 1 - \frac{Q_3}{Q_2}$

Also $\frac{Q_1}{Q_2} = \frac{600}{T}$

and $\frac{Q_2}{Q_3} = \frac{T}{150}$

$\therefore \frac{Q_1}{Q_2} = \frac{Q_3}{Q_2} = \frac{T_2}{T_1}$
\[
\frac{600}{T} - 1 = 1 - \frac{150}{T} \\
\frac{600 - T}{T} = \frac{T - 150}{T} \\
\therefore T = 375 K
\]

(ii) When efficiencies are equal

\[\eta_1 = 1 - \frac{Q_2}{Q_1} \quad \text{and} \quad \eta_2 = 1 - \frac{Q_3}{Q_2}\]

As \(\eta_1 = \eta_2\)

\[1 - \frac{Q_2}{Q_1} = 1 - \frac{Q_3}{Q_2}\]

\[1 - \frac{T}{600} = 1 - \frac{150}{T}\]

\[\frac{600 - T}{600} = \frac{T - 150}{T}\]

\[\frac{T}{600} = \frac{150}{T}\]

\[T^2 = 600 \times 150\]

\[\therefore T = 300 K\]

8.5 A Carnot engine whose low temperature reservoir is at 7°C has an efficiency of 50%. It is desired to increase the efficiency to 70%. By how many degrees should the temperature of the high temperature reservoir be increased?

\textbf{Data} : \(\eta_1 = 50\% = 0.5 ; \ T_2 = 7 + 273 = 280K ; \ \eta_2 = 70\% = 0.7\)

\textbf{Solution} : \(\eta_1 = 1 - \frac{T_2}{T_1} ; 0.5 = 1 - \frac{280}{T_1} ; \therefore T_1 = 560 K\)

\[\text{Let the temperature of the high temperature reservoir be } T'_1\]
\[ \eta_2 = 1 - \frac{T_2}{T_1'} \quad \text{and} \quad 0.7 = 1 - \frac{280}{T_1'} \quad \therefore T_1' = 933.3 \text{ K} \]

\[ \therefore \text{The temperature of the reservoir should be increased by} \quad 933.3 \text{ K} - 560 \text{ K} = 373.3 \text{ K} \]

8.6 A Carnot engine is operated between two reservoirs at temperature 177\(^\circ\) C and 77\(^\circ\) C. If the engine receives 4200 J of heat energy from the source in each cycle, calculate the amount of heat rejected to the sink in each cycle. Calculate the efficiency and work done by the engine.

**Data:**
- \( T_1 = 177\text{ }\text{o}\text{ C} = 177 + 273 = 450 \text{ K} \)
- \( T_2 = 77\text{ }\text{o}\text{ C} = 77 + 273 = 350 \text{ K} \)
- \( Q_1 = 4200 \text{ J} \quad Q_2 = ? \)

**Solution:**
\[
\frac{Q_2}{Q_1} = \frac{T_2}{T_1} \]

\[ Q_2 = Q_1 \cdot \frac{T_2}{T_1} = 4200 \times \frac{350}{450} = 3266.67 \text{ J} \]

Efficiency, \( \eta = 1 - \frac{T_2}{T_1} \)

\[ \eta = 1 - \frac{350}{450} = 0.2222 = 22.22\% \]

Work done
\[ W = Q_1 \cdot Q_2 = 4200 \cdot 3266.67 \]
\[ W = 933.33 \text{ J} \]

8.7 A Carnot engine has the same efficiency, when operated
(i) between 100 K and 500 K
(ii) between T K and 900 K
Find the value of T

**Solution:**
(i) Here \( T_1 = 500 \text{ K}; \ T_2 = 100 \text{ K} \)

\[ \eta = 1 - \frac{T_2}{T_1} = 1 - \frac{100}{500} = 1 - 0.2 = 0.8 \]

(ii) Now, \( T_1 = 900 \text{ K}; \ T_2 = T \text{ and } \eta = 0.8 \)
Again, \( \eta = 1 - \frac{T_2}{T_1} \)

\[ 0.8 = 1 - \frac{T}{900} \text{ or } \frac{T}{900} = 1 \cdot 0.8 = 0.2 \]

\( \therefore T = 180 \text{ K} \)

8.8 In a refrigerator heat from inside at 277 K is transferred to a room at 300 K. How many joule of heat will be delivered to the room for each joule of electric energy consumed ideally?

**Data :** \( T_1 = 300 \text{ K} ; T_2 = 277 \text{ K} \)

**Solution :** COP of a refrigerator

\[
\text{COP} = \frac{T_2}{T_1 - T_2} = \frac{277}{300 - 277} = 12.04 \quad \text{...(1)}
\]

Suppose for each joule of electric energy consumed an amount of heat \( Q_2 \) is extracted from the inside of the refrigerator. The amount of heat delivered to the room for each joule of electrical energy consumed is given by

\[ \delta Q_1 = Q_2 + W = Q_2 + 1 \quad \text{ (\( \because W = Q_1 - Q_2 \))} \]

\( \therefore \quad Q_1 - Q_2 = 1 \)

Also for a refrigerator, \( \text{COP} = \frac{Q_2}{\delta Q_1 - Q_2} = Q_2 \quad \text{...(2)} \)

From equations (1) and (2)

\( \text{(i.e) } Q_2 = 12.04 \)

\( \therefore \quad Q_1 = Q_2 + 1 = 12.04 + 1 = 13.04 \text{ J} \)

8.9 Two rods A and B of different material have equal length and equal temperature gradient. Each rod has its ends at temperatures \( T_1 \) and \( T_2 \). Find the condition under which rate of flow of heat through the rods A and B is same.

**Solution :** Suppose the two rods A and B have the same temperature difference \( T_1 - T_2 \) across their ends and the length of each rod is \( l \).

When the two rods have the same rate of heat conduction,
For the same rate of heat conduction, the areas of cross-section of the two rods should be inversely proportional to their coefficients of thermal conductivity.

A metal cube takes 5 minutes to cool from 60°C to 52°C. How much time will it take to cool to 44°C, if the temperature of the surroundings is 32°C?

**Solution:** While cooling from 60°C to 52°C

\[
\begin{align*}
\text{Rate of cooling} &= \frac{60 - 52}{5} = 1.6^\circ C/\text{minute} = \frac{1.6^\circ C}{60} \text{ per second} \\
\therefore \text{Average temperature while cooling} &= \frac{60 + 52}{2} = 56^\circ C \\
\therefore \text{Average temperature excess} &= 56 - 32 = 24^\circ C \\
\end{align*}
\]

According to Newton’s law of cooling,

\[
\text{Rate of cooling} \propto \text{Temperature excess}
\]

\[
\frac{1.6}{60} = K \times 24 \quad \ldots(1)
\]

Suppose that the cube takes \( t \) seconds to cool from 52°C to 44°C

\[
\begin{align*}
\therefore \text{Rate of cooling} &= \frac{52 - 44}{t} = \frac{8}{t} \\
\text{Average temperature while cooling} &= \frac{52 + 44}{2} = 48^\circ C \\
\therefore \text{Average temperature excess} &= 48 - 32 = 16^\circ C \\
\end{align*}
\]

According to Newton’s law, Rate of cooling = \( K \times (\text{Temperature excess}) \)

\[
\frac{8}{t} = K \times 16
\]

Dividing equation (1) by equation (2)

\[
\frac{1.6}{60} \times \frac{t}{8} = \frac{24}{16} = 450 \text{ s}
\]
Self evaluation

(The questions and problems given in this self evaluation are only samples. In the same way any question and problem could be framed from the text matter. Students must be prepared to answer any question and problem from the text matter, not only from the self evaluation.)

8.1 Avogadro number is the number of molecules in
(a) one litre of a gas at NTP
(b) one mole of a gas
(c) one gram of a gas
(d) 1 kg of a gas

8.2 First law of thermodynamics is a consequence of the conservation of
(a) momentum
(b) charge
(c) mass
(d) energy

8.3 At a given temperature, the ratio of the RMS velocity of hydrogen to the RMS velocity of oxygen is
(a) 4
(b) $\frac{1}{4}$
(c) 16
(d) 8

8.4 The property of the system that does not change during an adiabatic change is
(a) temperature
(b) volume
(c) pressure
(d) heat

8.5 For an ant moving on the horizontal surface, the number of degrees of freedom of the ant will be:
(a) 1
(b) 2
(c) 3
(d) 6
8.6 The translational kinetic energy of gas molecules for one mole of the gas is equal to:

(a) \( \frac{3}{2}RT \)  
(b) \( \frac{2}{3}kT \)  
(c) \( \frac{1}{2}RT \)  
(d) \( \frac{3}{2}kT \)

8.7 The internal energy of a perfect gas is

(a) partly kinetic and partly potential  
(b) wholly potential  
(c) wholly kinetic  
(d) depends on the ratio of two specific heats

8.8 A refrigerator with its power on, is kept in a closed room. The temperature of the room will

(a) rise  
(b) fall  
(c) remains the same  
(d) depend on the area of the room

8.9 A beaker full of hot water is kept in a room. If it cools from 80°C to 75°C in \( t_1 \) minutes, from 75°C to 70°C in \( t_2 \) minutes and from 70°C to 65°C in \( t_3 \) minutes then

(a) \( t_1 = t_2 = t_3 \)  
(b) \( t_1 < t_2 = t_3 \)  
(c) \( t_1 < t_2 < t_3 \)  
(d) \( t_1 > t_2 > t_3 \)

8.10 Which of the following will radiate heat to the large extent

(a) white polished surface  
(b) white rough surface  
(c) black polished surface  
(d) black rough surface

8.11 A block of ice in a room at normal temperature

(a) does not radiate  
(b) radiates less but absorbs more  
(c) radiates more than it absorbs  
(d) radiates as much as it absorbs
8.12 What are the postulates of Kinetic theory of gases?
8.13 Derive an expression for the average kinetic energy of the molecule of gas.
8.14 Two different gases have exactly the same temperature. Do the molecules have the same RMS speed?
8.15 Explain internal energy. What is its value in one complete cyclic process?
8.16 What are degrees of freedom?
8.17 State the law of equipartition of energy and prove that for a diatomic gas, the ratio of the two specific heats at room temperature is $\frac{7}{5}$.
8.18 Distinguish between isothermal and adiabatic process
8.19 Define isothermal process. Derive an expression for the work done during the process.
8.20 A gas has two specific heats, whereas liquid and solid have only one. Why?
8.21 Derive an expression for the work done in one cycle during an adiabatic process
8.22 Define molar specific heat at constant pressure.
8.23 Derive Meyer’s relation.
8.24 What is an indicator diagram?
8.25 Distinguish between reversible process and irreversible process with examples.
8.26 Is it possible to increase the temperature of a gas without the addition of heat? Explain.
8.27 On driving a scooter for a long time the air pressure in the tyre slightly increases why?
8.28 How is second law of thermodynamics different from first law of thermodynamics?
8.29 Define Clausius statement.
8.30 Describe the working of Carnot engine and derive its efficiency.

8.31 Give an example for a heat pump.

8.32 A heat engine with 100% efficiency is only a theoretical possibility. Explain.

8.33 What is Coefficient of Performance? Derive the relation between COP and efficiency.

8.34 Why are ventilators provided in our houses?

8.35 Define temperature gradient.

8.36 Define steady state in thermal conduction of heat.

8.37 What are the factors upon which coefficient of thermal conductivity depends?

8.38 Write the applications of Kirchoff’s law.

8.39 Define absorptive power.

8.40 Define Stefan’s law.

8.41 Explain Fery’s concept of a perfect black body.

8.42 State Wien’s displacement law.

8.43 State Newton’s law of cooling. Explain the experimental verification of Newton’s law of cooling.

8.44 Why does a piece of red glass when heated and taken out glow with green light?

8.45 Define solar constant.

8.46 Describe the working of pyrheliometer.

Problems

8.47 Calculate the kinetic energy of translational motion of a molecule of a diatomic gas at 320 K.

8.48 Calculate the rms velocity of hydrogen molecules at NTP (One mole of hydrogen occupies 22.4 litres at NTP).

8.49 The RMS speed of dust particles in air at NTP is $2.2 \times 10^{-2}$ ms$^{-1}$. Find the average mass of the particles.
8.50 Find the number of molecules in $10 \times 10^{-6}$ m$^3$ of a gas at NTP, if the mass of each molecule is $4 \times 10^{-26}$ kg and the RMS velocity is 400 m s$^{-1}$.

8.51 Calculate the molecular kinetic energy of translation of one mole of hydrogen at NTP, ($R = 8.31$ J mol$^{-1}$ K$^{-1}$).

8.52 Find the work done by 1 mole of perfect gas when it expands isothermally to double its volume. The initial temperature of the gas is 0°C ($R=8.31$ J mol$^{-1}$ K$^{-1}$).

8.53 A tyre pumped to a pressure of 3 atmosphere suddenly bursts. Calculate the fall in temperature if the temperature of air before expansion is 27°C and $\gamma = 1.4$.

8.54 A certain volume of dry air at NTP is expanded into three times its volume, under (i) isothermal condition (ii) adiabatic condition. Calculate in each case, the final pressure and final temperature, ($\gamma$ for air = 1.4).

8.55 A gas is suddenly compressed to $\frac{1}{2}$ of its original volume. If the original temperature is 300 K, find the increase in temperature (Assume $\gamma = 1.5$).

8.56 A system absorbs 8.4 kJ of heat and at the same time does 500 J of work. Calculate the change in internal energy of the system.

8.57 How many metres can a man weighing 60 kg, climb by using the energy from a slice of bread which produces a useful work of $4.2 \times 10^5$ J. Efficiency of human body is 28%.

8.58 The wavelength with maximum energy emitted from a certain star in our galaxy is $1.449 \times 10^{-5}$ cm. Calculate the temperature of star.

8.59 The surface temperature of a spherical hot body is 1000 K. Calculate the rate at which energy is radiated. ($\sigma = 5.67 \times 10^{-8}$ W m$^{-2}$ K$^{-4}$)

8.60 The opposite faces of the top of an electric oven are at a difference of temperature of 100°C and the area of the top surface and its
thickness are 300 cm² and 0.2 cm respectively. Find the quantity of heat that will flow through the top surface in one minute. 

\(K = 0.2 \text{ W m}^{-1} \text{ K}^{-1}\)

8.61 Compare the rate of loss of heat from a black metal sphere at 227°C with the rate of loss of heat from the same sphere at 127°C. The temperature of the surroundings is 27°C.

8.62 The ratio of radiant energies radiated per unit surface area by two bodies is 16 : 1. The temperature of hotter body is 1000 K. Calculate the temperature of the other body. Hint: \(E \propto (T^4 - T_0^4)\)

8.63 Calculate the surface temperature of the Sun (\(\lambda_m = 4753 \text{ Å}\)).

8.64 A hot solid takes 10 minutes to cool from 60°C to 50°C. How much further time will it take to cool to 40°C, if the room temperature is 20°C?

8.65 An object is heated and then allowed to cool when its temperature is 70°C, its rate of cooling is 3°C per minute and when the temperature is 60°C, the rate of cooling is 2.5°C per minute. Determine the temperature of the surroundings.
### Answers

<table>
<thead>
<tr>
<th>Question</th>
<th>Answer</th>
</tr>
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<tbody>
<tr>
<td>8.1</td>
<td>(b)</td>
</tr>
<tr>
<td>8.2</td>
<td>(d)</td>
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<tr>
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<td>(a)</td>
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<td>8.4</td>
<td>(d)</td>
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<td>(b)</td>
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<td>(c)</td>
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<tr>
<td>8.8</td>
<td>(a)</td>
</tr>
<tr>
<td>8.9</td>
<td>(c)</td>
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<td>8.10</td>
<td>(d)</td>
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<tr>
<td>8.11</td>
<td>(b)</td>
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<th>Value</th>
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<td>$6.624 \times 10^{-21}$ J</td>
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<td>8.48</td>
<td>1845 m s$^{-1}$</td>
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<tr>
<td>8.49</td>
<td>$2.335 \times 10^{17}$ kg</td>
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<td>8.50</td>
<td>$4.748 \times 10^{20}$</td>
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<tr>
<td>8.51</td>
<td>$3.403 \times 10^3$ J</td>
</tr>
<tr>
<td>8.52</td>
<td>1572.6 J</td>
</tr>
<tr>
<td>8.53</td>
<td>80.8 K</td>
</tr>
<tr>
<td>8.54</td>
<td>$3.376 \times 10^4$ N m$^{-2}$ ; 273 K ; $2.171 \times 10^4$ N m$^{-2}$ ; 176 K</td>
</tr>
<tr>
<td>8.55</td>
<td>124.2 K</td>
</tr>
<tr>
<td>8.56</td>
<td>7900 J</td>
</tr>
<tr>
<td>8.57</td>
<td>200 m</td>
</tr>
<tr>
<td>8.58</td>
<td>200000 K</td>
</tr>
<tr>
<td>8.59</td>
<td>$5.67 \times 10^4$ W m$^2$</td>
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<tr>
<td>8.60</td>
<td>18 K J</td>
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<tr>
<td>8.61</td>
<td>31 : 10</td>
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<tr>
<td>8.62</td>
<td>500 K</td>
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<tr>
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<td>6097 K</td>
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<tr>
<td>8.64</td>
<td>840 seconds</td>
</tr>
<tr>
<td>8.65</td>
<td>$10^9$ C</td>
</tr>
</tbody>
</table>
9. Ray Optics

Light rays and beams

A ray of light is the direction along which the light energy travels. In practice a ray has a finite width and is represented in diagrams as straight lines. A beam of light is a collection of rays. A search light emits a parallel beam of light (Fig. 9.1a). Light from a lamp travels in all directions which is a divergent beam. (Fig. 9.1b). A convex lens produces a convergent beam of light, when a parallel beam falls on it (Fig. 9.1c).

\[ \text{(a) Parallel beam} \quad \text{(b) Divergent Beam} \quad \text{(c) Convergent Beam} \]

9.1 Reflection of light

Highly polished metal surfaces reflect about 80% to 90% of the light incident on them. Mirrors in everyday use are therefore usually made of depositing silver on the backside of the glass. The largest reflector in the world is a curved mirror nearly 5 metres across, whose front surface is coated with aluminium. It is the Hale Telescope on the top of Mount Palomar, California, U.S.A. Glass by itself, will also reflect light, but the percentage is small when compared with the case of silvered surface. It is about 5% for an air-glass surface.

9.1.1 Laws of reflection

Consider a ray of light, AO, incident on a plane mirror XY at O. It is reflected along OB. Let the normal ON is drawn at the point of incidence. The angle AON between the incident ray and the normal is called angle of incidence, \( i \) (Fig. 9.2) the angle BON between the reflected ray and the normal is called angle of reflection, \( r \). Experiments
show that: (i) The incident ray, the reflected ray and the normal drawn to the reflecting surface at the point of incidence, all lie in the same plane.

(ii) The angle of incidence is equal to the angle of reflection. (i.e) \( i = r \).

These are called the laws of reflection.

9.1.2 Deviation of light by plane mirror

Consider a ray of light, \( AO \), incident on a plane mirror \( XY \) (Fig. 9.3) at \( O \). It is reflected along \( OB \). The angle \( AOX \) made by \( AO \) with \( XY \) is known as the glancing angle \( \alpha \) with the mirror. Since the angle of reflection is equal to the angle of incidence, the glancing angle \( BOY \) made by the reflected ray \( OB \) with the mirror is also equal to \( \alpha \).

The light has been deviated from a direction \( AO \) to a direction \( OB \). Since angle \( COY = \angle AOX \), it follows that

angle of deviation, \( d = 2\alpha \)

So, in general, the angle of deviation of a ray by a plane mirror or a plane surface is twice the glancing angle.

9.1.3 Deviation of light due to rotation of a mirror

Let us consider a ray of light \( AO \) incident on a plane mirror \( XY \) at \( O \). It is reflected along \( OB \). Let \( \alpha \) be the glancing angle with \( XY \) (Fig. 9.4). We know that the angle of deviation \( COB = 2\alpha \).

Suppose the mirror is rotated through an angle \( \theta \) to a position \( XY' \).
The same incident ray AO is now reflected along OP. Here the glancing angle with XY' is \((\alpha + \theta)\). Hence the new angle of deviation \(COP = 2(\alpha + \theta)\). The reflected ray has thus been rotated through an angle \(BOP\) when the mirror is rotated through an angle \(\theta\).

\[
|BOP| = |COP| - |COB|
\]

\[
|BOP| = 2(\alpha + \theta) - 2\alpha = 2\theta
\]

For the same incident ray, when the mirror is rotated through an angle, the reflected ray is rotated through twice the angle.

### 9.2 Image in a plane mirror

Let us consider a point object A placed in front of a plane mirror M as shown in the Fig. 9.5. Consider a ray of light AO from the point object incident on the mirror and reflected along OB. Draw the normal ON to the mirror at O.

The angle of incidence \(AON = \text{angle of reflection } BON\)

Another ray AD incident normally on the mirror at D is reflected back along DA. When BO and AD are produced backwards, they meet at I. Thus the rays reflected from M appear to come from a point I behind the mirror.

From the figure

\[
|AON| = |DAO|, \text{ alternate angles and } |BON| = |DIO|, \text{ corresponding angles it follows that } |DAO| = |DIO|.
\]

The triangles ODA and ODI are congruent

\[\therefore AD = ID\]

For a given position of the object, A and D are fixed points. Since \(AD = ID\), the point I is also fixed. It should be noted that AO = OI. So the object and its image in a plane mirror are at equal perpendicular distances from the mirror.
9.2.1 Virtual and real images

An object placed in front of a plane mirror has an image behind the mirror. The rays reflected from the mirror do not actually meet through I, but only appear to meet and the image cannot be received on the screen, because the image is behind the mirror. This type of image is called an unreal or virtual image (Fig. 9.6a).

If a convergent beam is incident on a plane mirror, the reflected rays pass through a point I in front of M, as shown in the Fig. 9.6b. In the Fig. 9.6a, a real object (divergent beam) gives rise to a virtual image. In the Fig. 9.6b, a virtual object (convergent beam) gives a real image. Hence plane mirrors not only produce virtual images for real objects but also produce real images for virtual objects.

9.2.2 Characteristics of the image formed by a plane mirror

(i) Image formed by a plane mirror is as far behind the mirror as the object is in front of it and it is always virtual.

(ii) The image produced is laterally inverted.

(iii) The minimum size of the mirror required to see the complete image of the object is half the size of the object.

(iv) If the mirror turns by an angle $\theta$, the reflected ray turns through an angle $2\theta$.

(v) If an object is placed between two plane mirrors inclined at an angle $\theta$, then the number of images formed is $n = \frac{360^\circ}{\theta} - 1$
9.3 Reflection at curved surfaces

In optics we are mainly concerned with curved mirrors which are the part of a hollow sphere (Fig. 9.7). One surface of the mirror is silvered. Reflection takes place at the other surface. If the reflection takes place at the concave surface, (which is towards the centre of the sphere) it is called concave mirror. If the reflection takes place at the convex surface, (which is away from the centre of the sphere) it is called convex mirror. The laws of reflection at a plane mirror are equally true for spherical mirrors also.

The centre of the sphere, of which the mirror is a part is called the centre of curvature (C).

The geometrical centre of the mirror is called its pole (P).

The line joining the pole of the mirror and its centre of curvature is called the principal axis.

The distance between the pole and the centre of curvature of the spherical mirror is called the radius of curvature of the mirror and is also equal to the radius of the sphere of which the mirror forms a part.

When a parallel beam of light is incident on a spherical mirror, the point where the reflected rays converge (concave mirror) or appear to
diverge from the point (convex mirror) on the principal axis is called the principal focus ($F$) of the mirror. The distance between the pole and the principal focus is called the focal length ($f$) of the mirror (Fig. 9.8).

### 9.3.1 Images formed by a spherical mirror

The images produced by spherical mirrors may be either real or virtual and may be either larger or smaller than the object. The image can be located by graphical construction as shown in Fig. 9.9 by adopting any two of the following rules.

1. A ray parallel to the principal axis after reflection by a concave mirror passes through the principal focus of the concave mirror and appear to come from the principal focus in a convex mirror.
2. A ray passing through the centre of curvature retraces its path after reflection.
3. A ray passing through the principal focus, after reflection is rendered parallel to the principal axis.
4. A ray striking the pole at an angle of incidence $i$ is reflected at the same angle $i$ to the axis.

![Fig. 9.9 Formation of images in concave mirror](image)

### 9.3.2 Image formed by a convex mirror

In a convex mirror irrespective of the position of the object, the image formed is always virtual, erect but diminished in size. The image lies between the pole and the focus (Fig. 9.10).

![Fig. 9.10 Image formed by convex mirror](image)
In general, real images are located in front of a mirror while virtual images behind the mirror.

### 9.3.3 Cartesian sign convention

The following sign conventions are used.

1. All distances are measured from the pole of the mirror (in the case of lens from the optic centre).

2. The distances measured in the same direction as the incident light, are taken as positive.

3. The distances measured in the direction opposite to the direction of incident light are taken as negative.

4. Heights measured perpendicular to the principal axis, in the upward direction are taken as positive.

5. Heights measured perpendicular to the principal axis, in the downward direction are taken as negative.

6. The size of the object is always taken as positive, but image size is positive for erect image and negative for an inverted image.

7. The magnification is positive for erect (and virtual) image, and negative for an inverted (and real) image.

### 9.3.4 Relation between $u$, $v$ and $f$ for spherical mirrors

A mathematical relation between object distance $u$, the image distance $v$ and the focal length $f$ of a spherical mirror is known as mirror formula.
(i) **Concave mirror - real image**

Let us consider an object $OO'$ on the principal axis of a concave mirror beyond $C$. The incident and the reflected rays are shown in the Fig 9.12. A ray $O'A$ parallel to principal axis is incident on the concave mirror at $A$, close to $P$. After reflections the ray passes through the focus $F$. Another ray $O'C$ passing through centre of curvature $C$, falls normally on the mirror and reflected back along the same path. A third ray $O'P$ incident at the pole $P$ is reflected along $PI'$. The three reflected rays intersect at the point $I'$. Draw perpendicular $I'I$ to the principal axis. $II'$ is the real, inverted image of the object $OO'$.

Right angled triangles, $II'P$ and $OO'P$ are similar.

$$\frac{II'}{OO'} = \frac{PI}{PO} \quad \ldots \ (1)$$

Right angled triangles $II'F$ and $APF$ are also similar ($A$ is close to $P$; hence $AP$ is a vertical line)

$$\frac{II'}{AP} = \frac{IF}{PF}$$

$AP = OO'$. Therefore the above equation becomes,

$$\frac{II'}{OO'} = \frac{IF}{PF} \quad \ldots \ (2)$$

Comparing the equations (1) and (2)

$$\frac{PI}{PO} = \frac{IF}{PF} \quad \ldots \ (3)$$

But, $IF = PI - PF$

Therefore equation (3) becomes,

$$\frac{PI}{PO} = \frac{PI - PF}{PF} \quad \ldots \ (4)$$

Using sign conventions, we have $PO = -u$, $PI = -v$ and $PF = f$
Substituting the values in the above equation, we get

\[
\frac{-v}{-u} = \frac{-v - (-f)}{-f} \quad \text{(or)}
\]

\[
\frac{v}{u} = \frac{v - f}{f} = \frac{v}{f} - 1
\]

Dividing by \( v \) and rearranging,

\[
\frac{1}{u} + 1 = \frac{1}{f}
\]

This is called *mirror equation*. The same equation can be obtained for virtual image also.

**(ii) Convex mirror - virtual image**

Let us consider an object \( OO' \) anywhere on the principal axis of a convex mirror. The incident and the reflected rays are shown in the Fig. 9.13. A ray \( O'A \) parallel to the principal axis incident on the convex mirror at \( A \) close to \( P \). After reflection the ray appears to diverge from the focus \( F \). Another ray \( O'C \) passing through centre of curvature \( C \), falls normally on the mirror and is reflected back along the same path. A third ray \( O'P \) incident at the pole \( P \) is reflected along \( PQ \). The three reflected rays when produced appear to meet at the point \( I' \). Draw perpendicular \( II' \) to the principal axis. \( II' \) is the virtual image of the object \( OO' \).

Right angled triangles, \( II'P \) and \( OO'P \) are similar.

\[
\therefore \frac{II'}{OO'} = \frac{PI}{FO}
\]

\[
\text{... (1)}
\]

Right angled triangles \( II'F \) and \( APF \) are also similar (\( A \) is close to \( P \); hence \( AP \) is a vertical line)

\[
\frac{II'}{AP} = \frac{IF}{PF}
\]
\[ AP = OO'. \] Therefore the above equation becomes,

\[ \frac{II'}{OO'} = \frac{IF}{PF} \quad \cdots (2) \]

Comparing the equations (1) and (2)

\[ \frac{PI}{PO} = \frac{IF}{PF} \quad \cdots (3) \]

But, \( IF = PF - PI \). Therefore equation (3) becomes,

\[ \frac{PI}{PO} = \frac{PF - PI}{PF} \]

Using sign conventions, we have \( PO = -u \), \( PI = +v \) and \( PF = +f \).

Substituting the values in the above equation, we get

\[ \frac{+v}{-u} = \frac{+f - (+v)}{+f} \quad \text{(or)} \quad -\frac{v}{u} = \frac{f - v}{f} = 1 - \frac{v}{f} \]

Dividing by \( v \) and rearranging we get,

\[ \frac{1}{u} + \frac{1}{v} = \frac{1}{f} \]

This is called \textit{mirror equation for convex mirror producing virtual image}.

### 9.3.5 Magnification

The linear or transverse magnification is defined as the ratio of the size of the image to that of the object.

\[ '\text{Magnification} = \frac{\text{size of the image}}{\text{size of the object}} = \frac{h_2}{h_1} \]

where \( h_1 \) and \( h_2 \) represent the size of the object and image respectively.

From Fig. 9.12 it is known that \( \frac{II'}{OO'} = \frac{PI}{PO} \)

Applying the sign conventions,

\[ II' = -h_2 \quad \text{(height of the image measured downwards)} \]

\[ OO' = +h_1 \quad \text{(height of the object measured upwards)} \]

\[ PI = -v \quad \text{(image distance against the incident light)} \]
PO = \( -u \) (object distance against the incident light)

Substituting the values in the above equation, we get

\[
m = \frac{-h_2}{h_1} = \frac{-v}{-u} \quad \text{(or)} \quad m = \frac{h_2}{h_1} = \frac{-v}{u}
\]

For an erect image, \( m \) is positive and for an inverted image, \( m \) is negative. This can be checked by substituting values for convex mirror also.

Using mirror formula, the equation for magnification can also be obtained as

\[
m = \frac{h_2}{h_1} = \frac{-v}{u} = \frac{f - v}{f} = \frac{f}{f - u}
\]

This equation is valid for both convex and concave mirrors.

### 9.4 Total internal reflection

When a ray of light \( AO \) passes from an optically denser medium to a rarer medium, at the interface \( XY \), it is partly reflected back into the same medium along \( OB \) and partly refracted into the rarer medium along \( OC \) (Fig. 9.14).

If the angle of incidence is gradually increased, the angle of refraction \( r \) will also gradually increase and at a certain stage \( r \) becomes 90°. Now the refracted ray \( OC \) is bent so much away from the normal and it grazes the surface of separation of two media. The angle of incidence in the denser medium at which the refracted ray just grazes the surface of separation is called the critical angle \( c \) of the denser medium.

If \( i \) is increased further, refraction is not possible and the incident

![Fig. 9.14 Total internal reflection](image-url)
ray is totally reflected into the same medium itself. This is called total internal reflection.

If \( \mu_d \) is the refractive index of the denser medium then, from Snell’s Law, the refractive index of air with respect to the denser medium is given by,

\[
\frac{\mu_d}{\mu_a} = \frac{\sin i}{\sin r}
\]

\[
\frac{1}{\mu_d} = \frac{\sin i}{\sin r} \quad (\because \mu_a = 1 \text{ for air})
\]

If \( r = 90^\circ \), \( i = c \)

\[
\frac{\sin c}{\sin 90^\circ} = \frac{1}{\mu_d} \quad (\text{or}) \quad \sin c = \frac{1}{\mu_d} \quad \text{or} \quad c = \sin^{-1}\left(\frac{1}{\mu_d}\right)
\]

If the denser medium is glass, \( c = \sin^{-1}\left(\frac{1}{\mu_d}\right) \)

Hence for total internal reflection to take place (i) light must travel from a denser medium to a rarer medium and (ii) the angle of incidence inside the denser medium must be greater than the critical angle i.e. \( i > c \).

**Table 9.1 Critical angle for some media**

*(NOT FOR EXAMINATION)*

<table>
<thead>
<tr>
<th>Medium</th>
<th>Refractive index</th>
<th>Critical angle</th>
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</thead>
<tbody>
<tr>
<td>Water</td>
<td>1.33</td>
<td>48.75°</td>
</tr>
<tr>
<td>Crown glass</td>
<td>1.52</td>
<td>41.14°</td>
</tr>
<tr>
<td>Dense flint glass</td>
<td>1.62</td>
<td>37.31°</td>
</tr>
<tr>
<td>Diamond</td>
<td>2.42</td>
<td>24.41°</td>
</tr>
</tbody>
</table>
9.4.1 Applications

(i) Diamond

Total internal reflection is the main cause of the brilliance of diamonds. The refractive index of diamond with respect to air is 2.42. Its critical angle is 24.41°. When light enters diamond from any face at an angle greater than 24.41° it undergoes total internal reflection. By cutting the diamond suitably, multiple internal reflections can be made to occur.

(ii) Optical fibres

The total internal reflection is the basic principle of optical fibre. An optical fibre is a very thin fibre made of glass or quartz having radius of the order of micrometer (10⁻⁶ m). A bundle of such thin fibres forms a ‘light pipe’ (Fig. 9.15a).

Fig. 9.15b shows the principle of light transmission inside an optical fibre. The refractive index of the material of the core is higher than that of the cladding. When the light is incident at one end of the fibre at a small angle, the light passes inside, undergoes repeated total internal reflections along the fibre and finally comes out. The angle of incidence is always larger than the critical angle of the core material with respect to its cladding. Even if the fibre is bent or twisted, the light can easily travel through the fibre.

Light pipes are used in medical and optical examination. They are also used to transmit communication signals.

9.5 Michelson’s method

A.A. Michelson, an American physicist, spent many years of his life in measuring the velocity of light and he devised a method in the year 1926 which is considered as accurate.
The experimental set up is shown in Fig. 9.16. Light from an arc source after passing through a narrow slit S is reflected from one face a of an octagonal mirror R. The ray after reflections at small fixed mirrors b and c is then rendered parallel by a concave mirror M₁ placed in the observing station on Mt. Wilson. This parallel beam of light travels a distance of 35 km and falls on another concave mirror M₂ placed at Mt. St Antonio, and it is reflected to a plane mirror d placed at the focus of the concave mirror M₂. The ray of light from d is rendered parallel after getting reflected by M₂ and travels back to the concave mirror M₁.

After reflections at M₁ and the plane mirrors e and f, the ray falls on the opposite face a₁ of the octagonal mirror. The final image which is totally reflected by a total reflecting prism P, is viewed through an eye piece E.

When the octagonal mirror is stationary, the image of the slit is seen through the eye piece. When it is rotated the image disappears. The speed of rotation of R is suitably adjusted so that the image is seen again clearly as when R is stationary. The speed of revolution is measured by stroboscope.

Let $D$ be the distance travelled by light from face $a$ to face $a_1$ and $n$ be the number of rotations made by R per second.
The time taken by R to rotate through 45° or \( \frac{1}{8} \) of a rotation = \( \frac{1}{8n} \)

During this time interval, the distance travelled by the light = D

\[ \therefore \text{The velocity of light} \quad c = \frac{\text{Distance travelled}}{\text{Time taken}} = \frac{D}{\frac{1}{8n}} = 8nD. \]

In general, if the number of faces in the rotating mirror is \( N \), the velocity of light = \( NnD \).

The velocity of light determined by him is \( 2.99797 \times 10^8 \) m s\(^{-1}\).

**Importance of velocity of light**

The value of velocity of light in vacuum is of great importance in science. The following are some of the important fields where the value of velocity of light is used.

1. **Frequency - wavelength relation**: From the relation \( c = \nu \lambda \), the frequency of electromagnetic radiations can be calculated if the wavelength is known and vice versa.

2. **Relativistic mass variation with velocity**: Theory of relativity has shown that the mass \( m \) of a moving particle varies with its velocity \( v \) according to the relation

\[
   m = \frac{m_0}{\sqrt{1 - \frac{v^2}{c^2}}}
\]

   Here \( m_0 \) is the rest mass of the particle.

3. **Mass - Energy relation**: \( E = mc^2 \) represents conversion of mass into energy and energy into mass. The energy released in nuclear fission and fusion is calculated using this relation.

4. **Measurement of large distance in Astronomy**: Light year is a unit of distance used in astronomy. A light year is the distance travelled by light in one year. It is equal to \( 9.46 \times 10^{15} \) metre.

5. **Refractive index**: The refractive index \( \mu \) of a medium is given by

\[
   \mu = \frac{\text{velocity of light in vacuum}}{\text{velocity of light in medium}} = \frac{c}{v}
\]
9.6 Refraction of light

When a ray of light travels from one transparent medium into another medium, it bends while crossing the interface, separating the two media. This phenomenon is called refraction.

Image formation by spherical lenses is due to the phenomenon of refraction. The laws of refraction at a plane surface are equally true for refraction at curved surfaces also. While deriving the expressions for refraction at spherical surfaces, we make the following assumptions.

(i) The incident light is assumed to be monochromatic and
(ii) the incident pencil of light rays is very narrow and close to the principal axis.

9.6.1 Cartesian sign convention

The sign convention followed in the spherical mirror is also applicable to refraction at spherical surface. In addition to this two more sign conventions to be introduced which are:

(i) The power of a converging lens is positive and that of a diverging lens is negative.

(ii) The refractive index of a medium is always said to be positive. If two refractions are involved, the difference in their refractive index is also taken as positive.

9.6.2 Refraction at a spherical surface

Let us consider a portion of a spherical surface AB separating two media having refracting indices $\mu_1$ and $\mu_2$ (Fig. 9.17). This is symmetrical about an axis passing through the centre C and cuts the surface at P. The point P is called the pole of the surface. Let R be the radius of curvature of the surface.

Consider a point object O on the axis in the first medium. Consider two rays $OP$ and $OD$ originating from O. The ray $OP$ falls

![Fig. 9.17 Refraction at a spherical surface](image-url)
normally on $AB$ and goes into the second medium, undeviated. The ray $OD$ falls at $D$ very close to $P$. After refraction, it meets at the point $I$ on the axis, where the image is formed. $CE$ is the normal drawn to the point $D$. Let $i$ and $r$ be the angle of incidence and refraction respectively.

Let $\hat{DOP} = \alpha$, $\hat{DCP} = \beta$, $\hat{DIC} = \gamma$

Since $D$ is close to $P$, the angles $\alpha$, $\beta$ and $\gamma$ are all small. From the Fig. 9.17.

$$\tan \alpha = \frac{DP}{PO}, \quad \tan \beta = \frac{DP}{PC} \quad \text{and} \quad \tan \gamma = \frac{DP}{PI}$$

$$\therefore \alpha = \frac{DP}{PO}, \quad \beta = \frac{DP}{PC} \quad \text{and} \quad \gamma = \frac{DP}{PI}$$

From the $\triangle ODC$, $i = \alpha + \beta$ \hspace{2cm} ...(1)

From the $\triangle DCI$, $\beta = r + \gamma$ or $r = \beta - \gamma$ \hspace{2cm} ...(2)

From Snell’s Law, $\frac{\mu_2}{\mu_1} = \frac{\sin i}{\sin r}$ and for small angles of $i$ and $r$, we can write, $\mu_1 i = \mu_2 r$ \hspace{2cm} ...(3)

From equations (1), (2) and (3)

we get $\mu_1 (\alpha + \beta) = \mu_2 (\beta - \gamma)$ or $\mu_1 \alpha + \mu_2 \gamma = (\mu_2 - \mu_1)\beta$ \hspace{2cm} ...(4)

Substituting the values of $\alpha$, $\beta$ and $\gamma$ in equation (4)

$$\frac{\mu_1}{PO} \left( \frac{DP}{PO} \right) + \frac{\mu_2}{PI} \left( \frac{DP}{PI} \right) = \left( \mu_2 - \mu_1 \right) \frac{DP}{PC}$$

$$\frac{\mu_1}{PO} \cdot \frac{\mu_2}{PI} \left( \frac{\mu_2 - \mu_1}{PC} \right)$$

\hspace{2cm} ...(5)

As the incident ray comes from left to right, we choose this direction as the positive direction of the axis. Therefore $u$ is negative, whereas $v$ and $R$ are positive substitute $PO = -u$, $PI = +v$ and $PC = +R$ in equation (5),

$$\frac{\mu_1}{-u} + \frac{\mu_2}{v} = \frac{\mu_2 - \mu_1}{R}$$

$$\frac{\mu_2}{v} - \frac{\mu_1}{u} = \frac{\mu_2 - \mu_1}{R}$$

\hspace{2cm} \ldots (6)
Equation (6) represents the general equation for refraction at a spherical surface.

If the first medium is air and the second medium is of refractive index \( \mu \), then
\[
\frac{\mu}{v} - \frac{1}{u} = \frac{\mu - 1}{R}
\]

...(7)

9.6.3 Refraction through thin lenses

A lens is one of the most familiar optical devices. A lens is made of a transparent material bounded by two spherical surfaces. If the distance between the surfaces of a lens is very small, then it is a thin lens.

As there are two spherical surfaces, there are two centres of curvature \( C_1 \) and \( C_2 \) and correspondingly two radii of curvature \( R_1 \) and \( R_2 \). The line joining \( C_1 \) and \( C_2 \) is called the principal axis of the lens. The centre \( P \) of the thin lens which lies on the principal axis is called the optic centre.

9.6.4 Lens maker’s formula and lens formula

Let us consider a thin lens made up of a medium of refractive index \( \mu_2 \) placed in a medium of refractive index \( \mu_1 \). Let \( R_1 \) and \( R_2 \) be the radii of curvature of two spherical surfaces \( ACB \) and \( ADB \) respectively and \( P \) be the optic centre.

Consider a point object \( O \) on the principal axis. The ray OP falls normally on the spherical surface and goes through the lens undeviated. The ray OA falls at A very close to \( P \). After refraction at the surface \( ACB \) the image is formed at \( I' \). Before it does so, it is again refracted by the surface \( ADB \). Therefore the final image is formed at \( I \) as shown in Fig. 9.18.
The general equation for the refraction at a spherical surface is given by
\[ \frac{\mu_2}{v} - \frac{\mu_1}{u} = \frac{\mu_2 - \mu_1}{R} \] ...
\(1\)

For the refracting surface ACB, from equation (1) we write
\[ \frac{\mu_2}{v'} - \frac{\mu_1}{u} = \frac{\mu_2 - \mu_1}{R_1} \] ...
\(2\)

The image \(I'\) acts as a virtual object for the surface ADB and the final image is formed at I. The second refraction takes place when light travels from the medium of refractive index \(\mu_2\) to \(\mu_1\).

For the refracting surface ADB, from equation (1) and applying sign conventions, we have
\[ \frac{\mu_1}{v} - \frac{\mu_2}{v'} = \left( \frac{\mu_2 - \mu_1}{R_2} \right) \] ...
\(3\)

Adding equations (2) and (3)
\[ \frac{\mu_1}{v} - \frac{\mu_2}{v'} = \left( \frac{\mu_2 - \mu_1}{R_2} \right) \]

Dividing the above equation by \(\mu_1\)
\[ \frac{1}{v} - \frac{1}{u} = \left( \frac{\mu_2}{\mu_1} - 1 \right) \left[ \frac{1}{R_1} - \frac{1}{R_2} \right] \] ...
\(4\)

If the object is at infinity, the image is formed at the focus of the lens.

Thus, for \(u = \infty, v = f\). Then the equation (4) becomes.
\[ \frac{1}{f} = \left( \frac{\mu_2}{\mu_1} - 1 \right) \left[ \frac{1}{R_1} - \frac{1}{R_2} \right] \] ...
\(5\)

If the refractive index of the lens is \(\mu\) and it is placed in air, \(\mu_2 = \mu\) and \(\mu_1 = 1\). So the equation (5) becomes
\[ \frac{1}{f} = (\mu - 1) \left[ \frac{1}{R_1} - \frac{1}{R_2} \right] \] ...
\(6\)

This is called the lens maker’s formula, because it tells what curvature will be needed to make a lens of desired focal length. This formula is true for concave lens also.

Comparing equation (4) and (5)
we get \[ \frac{1}{v} - \frac{1}{u} = \frac{1}{f} \] ... (7)

which is known as the lens formula.

9.6.5 Magnification

Let us consider an object \( OO' \) placed on the principal axis with its height perpendicular to the principal axis as shown in Fig. 9.19. The ray \( OP \) passing through the optic centre will go undeviated. The ray \( O'A \) parallel to the principal axis must pass through the focus \( F_2' \). The image is formed where \( O'P' \) and \( AP' \) intersect. Draw a perpendicular from \( I' \) to the principal axis. This perpendicular \( II' \) is the image of \( OO' \).

The linear or transverse magnification is defined as the ratio of the size of the image to that of the object.

\[ \therefore \text{Magnification } m = \frac{\text{Size of the image}}{\text{Size of the object}} = \frac{II'}{OO'} = \frac{h_2}{h_1} \]

where \( h_1 \) is the height of the object and \( h_2 \) is the height of the image.

From the similar right angled triangles \( OOP \) and \( III' \), we have

\[ \frac{II'}{OO'} = \frac{PI}{PO} \]

Applying sign convention,

\[ II' = -h_2; \quad OO' = +h_1; \]
\[ PI = +v; \quad PO = -u; \]

Substituting this in the above equation, we get magnification

\[ m = \frac{-h_2}{+h_1} = \frac{+v}{-u} \]

\[ \therefore m = +\frac{v}{u} \]

The magnification is negative for real image and positive for virtual image. In the case of a concave lens, it is always positive.
Using lens formula the equation for magnification can also be obtained as:

\[ m = \frac{h_2}{h_1} = \frac{v}{u} = \frac{f - v}{f} = \frac{f}{f + u} \]

This equation is valid for both convex and concave lenses and for real and virtual images.

**9.6.6 Power of a lens**

Power of a lens is a measure of the degree of convergence or divergence of light falling on it. The power of a lens \((P)\) is defined as the reciprocal of its focal length.

\[ P = \frac{1}{f} \]

The unit of power is dioptre \((D)\) : 1 \(D\) = 1 \(m^{-1}\). The power of the lens is said to be 1 dioptre if the focal length of the lens is 1 metre. \(P\) is positive for converging lens and negative for diverging lens. Thus, when an optician prescribes a corrective lens of power + 0.5 \(D\), the required lens is a convex lens of focal length + 2 m. A power of -2.0 \(D\) means a concave lens of focal length -0.5 m.

**9.6.7 Combination of thin lenses in contact**

Let us consider two lenses \(A\) and \(B\) of focal length \(f_1\) and \(f_2\) placed in contact with each other. An object is placed at \(O\) beyond the focus of the first lens \(A\) on the common principal axis. The lens \(A\) produces an image at \(I_1\). This image \(I_1\) acts as the object for the second lens \(B\). The final image is produced at \(I\) as shown in Fig. 9.20. Since the lenses are thin, a common optical centre \(P\) is chosen.

Let \(PO = u\), object distance for the first lens \((A)\), \(PI = v\), final image distance and \(PI_1 = v_1\), image distance for the first lens \((A)\) and also object distance for second lens \((B)\).
For the image $I_1$ produced by the first lens $A$,
\[
\frac{1}{v_1} - \frac{1}{u} = \frac{1}{f_1} \quad \ldots(1)
\]

For the final image $I$, produced by the second lens $B$,
\[
\frac{1}{v} - \frac{1}{v_2} = \frac{1}{f_2} \quad \ldots(2)
\]

Adding equations (1) and (2),
\[
\frac{1}{v} - \frac{1}{u} = \frac{1}{f_1} + \frac{1}{f_2} \quad \ldots(3)
\]

If the combination is replaced by a single lens of focal length $F$ such that it forms the image of $O$ at the same position $I$, then
\[
\frac{1}{v} - \frac{1}{u} = \frac{1}{F} \quad \ldots(4)
\]

From equations (3) and (4)
\[
\frac{1}{F} = \frac{1}{f_1} + \frac{1}{f_2} \quad \ldots(5)
\]

This $F$ is the focal length of the equivalent lens for the combination.

The derivation can be extended for several thin lenses of focal lengths $f_1, f_2, f_3, \ldots$ in contact. The effective focal length of the combination is given by
\[
\frac{1}{F} = \frac{1}{f_1} + \frac{1}{f_2} + \frac{1}{f_3} + \ldots \quad \ldots(6)
\]

In terms of power, equation (6) can be written as
\[
P = P_1 + P_2 + P_3 + \ldots \quad \ldots(7)
\]

Equation (7) may be stated as follows:
The power of a combination of lenses in contact is the algebraic sum of the powers of individual lenses.

The combination of lenses is generally used in the design of objectives of microscopes, cameras, telescopes and other optical instruments.

9.7 Prism

A prism is a transparent medium bounded by the three plane faces. Out of the three faces, one is grounded and the other two are
polished. The polished faces are called refracting faces. The angle between the refracting faces is called angle of prism, or the refracting angle. The third face is called base of the prism.

**Refraction of light through a prism**

Fig. 9.21 shows the cross section of a triangular prism ABC, placed in air. Let 'A' be the refracting angle of the prism. A ray of light PQ incident on the refracting face AB, gets refracted along QR and emerges along RS. The angle of incidence and refraction at the two faces are $i_1$, $r_1$, $r_2$ and $i_2$ respectively. The angle between the incident ray PQ and the emergent ray RS is called angle of deviation, $d$.

In the $\triangle QER$, the exterior angle $\angle FER = \angle EQR + \angle ERQ$

$$d = (i_1 - r_1) + (i_2 - r_2)$$

$$\therefore d = (i_1 + i_2) - (r_1 + r_2) \ldots (1)$$

In the quadrilateral AQOR, the angles at Q and R are right angles

$$\angle Q + \angle R = 180^\circ$$

$$\therefore A + \angle QOR = 180^\circ \ldots (2)$$

Also, from the $\triangle QOR$

$$r_1 + r_2 + \angle QOR = 180^\circ \ldots (3)$$

From equation (2) and (3)

$$r_1 + r_2 = A \ldots (4)$$

Substituting in (1),

$$d = i_1 + i_2 - A$$

or $$A + d = i_1 + i_2 \ldots (5)$$

For a given prism and for a light of given wavelength, the angle of deviation depends upon the angle of incidence.
As the angle of incidence $i$ gradually increases, the angle of deviation $d$ decreases, reaches a minimum value $D$ and then increases. $D$ is called the angle of minimum deviation. It will be seen from the graph (Fig. 9.22) that there is only one angle of incidence for which the deviation is a minimum.

At minimum deviation position the incident ray and emergent ray are symmetric with respect to the base of the prism. (i.e) the refracted ray $QR$ is parallel to the base of the prism.

At the minimum deviation \( i_1 = i_2 = i \) and \( r_1 = r_2 = r \)

\[ \therefore \text{ from equation (4)} \quad 2r = A \quad \text{or} \quad r = \frac{A}{2} \]

and from equation (5) \( 2i = A + D \) or \( i = \frac{A + D}{2} \)

The refractive index is \( \mu = \frac{\sin i}{\sin r} \)

\[ \therefore \mu = \frac{\sin \left( \frac{A + D}{2} \right)}{\sin \left( \frac{A}{2} \right)} \]

9.8 Dispersion of light

Dispersion is the splitting of white light into its constituent colours. This band of colours of light is called its spectrum.

In the visible region of spectrum, the spectral lines are seen in the order from violet to red. The colours are given by the word VIBGYOR (Violet, Indigo, Blue, Green, Yellow, Orange and Red) (Fig. 9.23)
The origin of colour after passing through a prism was a matter of much debate in physics. Does the prism itself create colour in some way or does it only separate the colours already present in white light?

Sir Isaac Newton gave an explanation for this. He placed another similar prism in an inverted position. The emergent beam from the first prism was made to fall on the second prism (Fig. 9.24). The resulting emergent beam was found to be white light. The first prism separated the white light into its constituent colours, which were then recombined by the inverted prism to give white light. Thus it can be concluded that the prism does not create any colour but it only separates the white light into its constituent colours.

Dispersion takes place because the refractive index of the material of the prism is different for different colours (wavelengths). The deviation and hence the refractive index is more for violet rays of light than the corresponding values for red rays of light. Therefore the violet ray travels with a smaller velocity in glass prism than red ray. The deviation and the refractive index of the yellow ray are taken as the mean values. Table 9.2 gives the refractive indices for different wavelength for crown glass and flint glass.

**Table 9.2 Refractive indices for different wavelengths**

<table>
<thead>
<tr>
<th>Colour</th>
<th>Wave length (nm)</th>
<th>Crown glass</th>
<th>Flint glass</th>
</tr>
</thead>
<tbody>
<tr>
<td>Violet</td>
<td>396.9</td>
<td>1.533</td>
<td>1.663</td>
</tr>
<tr>
<td>Blue</td>
<td>486.1</td>
<td>1.523</td>
<td>1.639</td>
</tr>
<tr>
<td>Yellow</td>
<td>589.3</td>
<td>1.517</td>
<td>1.627</td>
</tr>
<tr>
<td>Red</td>
<td>656.3</td>
<td>1.515</td>
<td>1.622</td>
</tr>
</tbody>
</table>

The speed of light is independent of wavelength in vacuum. Therefore vacuum is a non-dispersive medium in which all colours travel with the same speed.
9.8.1 Dispersive power

The refractive index of the material of a prism is given by the relation

\[ \mu = \frac{\sin \frac{A + D}{2}}{\sin \frac{A}{2}} \]

Here \( A \) is the angle of the prism and \( D \) is the angle of minimum deviation.

If the angle of prism is small of the order of 10°, the prism is said to be small angled prism. When rays of light pass through such prisms the angle of deviation also becomes small.

If \( A \) be the refracting angle of a small angled prism and \( \delta \) the angle of deviation, then the prism formula becomes

\[ \mu = \frac{\sin \left(\frac{A + \delta}{2}\right)}{\sin \frac{A}{2}} \]

For small angles \( A \) and \( \delta \),

\[ \sin \frac{A + \delta}{2} = \frac{A + \delta}{2} \quad \text{and} \quad \sin \frac{A}{2} = \frac{A}{2} \]

\[ \therefore \mu = \frac{A + \delta}{2} \]

\[ \mu A = \frac{A}{2} + \delta \]

\[ \delta = (\mu - 1)A \]

... (1)

If \( \delta_v \) and \( \delta_r \) are the deviations produced for the violet and red rays and \( \mu_v \) and \( \mu_r \) are the corresponding refractive indices of the material of the small angled prism then,

for violet light,

\[ \delta_v = (\mu_v - 1)A \]

... (2)

for red light, \( \delta_r = (\mu_r - 1)A \)

... (3)

From equations (2) and (3)

\[ \delta_v - \delta_r = (\mu_v - \mu_r)A \]

... (4)

Fig. 9.25 Dispersive power
\( \delta_y - \delta_r \) is called the angular dispersion which is the difference in deviation between the extreme colours (Fig. 9.25).

If \( \delta_y \) and \( \mu_y \) are the deviation and refractive index respectively for yellow ray (mean wavelength) then,

for yellow light, \( \delta_y = (\mu_y - 1)A \) ... (5)

Dividing equation (4) by (5) we get

\[
\frac{\delta_v - \delta_r}{\delta_y} = \frac{\mu_v - \mu_r}{\mu_y - 1}
\]

The expression \( \frac{\delta_v - \delta_r}{\delta_y} \) is known as the dispersive power of the material of the prism and is denoted by \( \omega \).

\[
\therefore \omega = \frac{\mu_v - \mu_r}{\mu_y - 1}
\]

The dispersive power of the material of a prism is defined as the ratio of angular dispersion for any two wavelengths (colours) to the deviation of mean wavelength.

9.9 Spectrometer

The spectrometer is an optical instrument used to study the spectra of different sources of light and to measure the refractive indices of materials (Fig. 9.26). It consists of basically three parts. They are collimator, prism table and Telescope.

![Fig. 9.26 Spectrometer (NEED NOT DRAW IN THE EXAMINATION)](image-url)
**Collimator**

The collimator is an arrangement to produce a parallel beam of light. It consists of a long cylindrical tube with a convex lens at the inner end and a vertical slit at the outer end of the tube. The distance between the slit and the lens can be adjusted such that the slit is at the focus of the lens. The slit is kept facing the source of light. The width of the slit can be adjusted. The collimator is rigidly fixed to the base of the instrument.

**Prism table**

The prism table is used for mounting the prism, grating etc. It consists of two circular metal discs provided with three levelling screws. It can be rotated about a vertical axis passing through its centre and its position can be read with verniers $V_1$ and $V_2$. The prism table can be raised or lowered and can be fixed at any desired height.

**Telescope**

The telescope is an astronomical type. It consists of an eyepiece provided with cross wires at one end of the tube and an objective lens at its other end co-axially. The distance between the objective lens and the eyepiece can be adjusted so that the telescope forms a clear image at the cross wires, when a parallel beam from the collimator is incident on it.

The telescope is attached to an arm which is capable of rotation about the same vertical axis as the prism table. A circular scale graduated in half degree is attached to it.

Both the telescope and prism table are provided with radial screws for fixing them in a desired position and tangential screws for fine adjustments.

**9.9.1 Adjustments of the spectrometer**

The following adjustments must be made before doing the experiment with spectrometer.

(i) **Adjustment of the eyepiece**

The telescope is turned towards an illuminated surface and the eyepiece is moved to and fro until the cross wires are clearly seen.
(ii) Adjustment of the telescope

The telescope is adjusted to receive parallel rays by turning it towards a distant object and adjusting the distance between the objective lens and the eyepiece to get a clear image on the cross wire.

(iii) Adjustment of the collimator

The telescope is brought along the axial line with the collimator. The slit of the collimator is illuminated by a source of light. The distance between the slit and the lens of the collimator is adjusted until a clear image of the slit is seen at the cross wires of the telescope. Since the telescope is already adjusted for parallel rays, a well defined image of the slit can be formed, only when the light rays emerging from the collimator are parallel.

(iv) Levelling the prism table

The prism table is adjusted or levelled to be in horizontal position by means of levelling screws and a spirit level.

9.9.2 Determination of the refractive index of the material of the prism

The preliminary adjustments of the telescope, collimator and the prism table of the spectrometer are made. The refractive index of the prism can be determined by knowing the angle of the prism and the angle of minimum deviation.

(i) Angle of the prism (A)

The prism is placed on the prism table with its refracting edge facing the collimator as shown in Fig 9.27. The slit is illuminated by a sodium vapour lamp.

The parallel rays coming from the collimator fall on the two faces AB and AC.

The telescope is rotated to the position $T_1$ until the image of the slit, formed by the reflection at the face AB is made to coincide with the vertical cross wire of the telescope. The readings of the verniers are noted. The telescope is then rotated to the position $T_2$ where the image of the slit formed by the reflection at the face AC coincides with the vertical cross wire. The readings are again noted.
The difference between these two readings gives the angle rotated by the telescope. This angle is equal to twice the angle of the prism. Half of this value gives the angle of the prism A.

**(ii) Angle of minimum deviation (D)**

The prism is placed on the prism table so that the light from the collimator falls on a refracting face, and the refracted image is observed through the telescope (Fig. 9.28). The prism table is now rotated so that the angle of deviation decreases. A stage comes when the image stops for a moment and if we rotate the prism table further in the same direction, the image is seen to recede and the angle of deviation increases. The vertical cross wire of the telescope is made to coincide with the image of the slit where it turns back. This gives the minimum deviation position. The readings of the verniers are noted. Now the prism is removed and the telescope is turned to receive the direct ray and the vertical cross wire is made to coincide with the image. The readings of the verniers are noted. The difference between the two readings gives the angle of minimum deviation D.

The refractive index of the material of the prism μ is calculated using the formula \[ \mu = \frac{\sin \left( \frac{A + D}{2} \right)}{\sin \frac{A}{2}}. \]

The refractive index of a liquid may be determined in the same way using a hollow glass prism filled with the given liquid.

### 9.10 Rainbow

One of the spectacular atmospheric phenomena is the formation of rainbow during rainy days. The rainbow is also an example of dispersion of sunlight by the water drops in the atmosphere.

When sunlight falls on small water drops suspended in air during or after a rain, it suffers refraction, internal reflection and dispersion.
If the Sun is behind an observer and the water drops in front, the observer may observe two rainbows, one inside the other. The inner one is called primary rainbow having red on the outer side and violet on the inner side and the outer rainbow is called secondary rainbow, for which violet on the outer side and red on the inner side.

Fig. 9.29 shows the formation of primary rainbow. It is formed by the light from the Sun undergoing one internal reflection and two refractions and emerging at minimum deviation. It is however, found that the intensity of the red light is maximum at an angle of 43° and that of the violet rays at 41°. The other coloured arcs occur in between violet and red (due to other rain drops).

The formation of secondary rainbow is also shown in Fig. 9.31. It is formed by the light from the Sun undergoing two internal reflections and two refractions and also emerging at minimum deviation. In this case the inner red edge subtends an angle of 51° and the outer violet edge subtends an angle of 54°. This rainbow is less bright and narrower than the primary rainbow. Both primary and secondary rainbows exhibit all the colours of the solar spectrum.

From the ground level an arc of the rainbow is usually visible. A complete circular rainbow may be seen from an elevated position such as from an aeroplane.
Solved Problems

9.1 A man 2 m tall standing in front of a plane mirror whose eye is 1.90 m above the ground. What is the minimum size of the mirror required to see complete image?

Solution:

\[ M \quad - \quad \text{Mirror} \]

\[ FH \quad - \quad \text{Man} \]

\[ H \quad - \quad \text{Head} \]

\[ E \quad - \quad \text{Eye} \]

\[ F \quad - \quad \text{Feet} \]

A ray \( HA \) from the head, falls at \( A \) on the mirror and reflected to \( E \) along \( AE \). \( AD \) is the perpendicular bisector of \( HE \).

\[ \therefore AC = \frac{1}{2} \times 0.10 = 0.05 \text{ m.} \]

A ray \( FB \) from the feet, falls at \( B \) and reflected to \( E \) along \( BE \). \( BN \) is the perpendicular bisector of \( EF \).

\[ \therefore CB = \frac{1}{2} \times 1.90 = 0.95 \text{ m.} \]

\[ \therefore \text{The size of the mirror} = AC + CB = 0.05 \text{ m} + 0.95 \text{ m} \]

\[ \text{Size of the mirror} = 1 \text{ m} \]

9.2 An object of length 2.5 cm is placed at a distance of 1.5 times the focal length \( f \) from a concave mirror. Find the length of the image. Is the image is erect or inverted?

Data: \( f = -f; u = -1.5 \times f; h_1 = 2.5 \text{ cm}; h_2 = ? \)

Solution:

We know,

\[ \frac{1}{f} = \frac{1}{u} + \frac{1}{v} \]

\[ \frac{1}{v} = \frac{1}{f} - \frac{1}{u} = -\frac{1}{f} \cdot -1.5f \]

\[ h_2 = \frac{h_1 \times f}{u} \]

\[ h_2 = \frac{2.5 \times f}{-1.5f} \]

\[ h_2 = -\frac{2.5}{1.5} \]

\[ h_2 = -1.67 \text{ cm} \]

The image is inverted.
\[ \frac{1}{v} = \frac{1}{1.5f} - \frac{1}{f} \]

\[ v = -3f \]

Magnification,
\[ m = \frac{-v}{u} = \frac{-3f}{-1.5f} = -2 \]

But \[ \frac{h_2}{h_1} = m = -2 \]

\[ \therefore h_2 = -5 \text{ cm} \]

The length of the image is 5.0 cm. The –ve sign indicates that the image is inverted.

9.3 In Michelson’s method to determine the velocity of light in air, the distance travelled by light between reflections from the opposite faces of the octagonal mirror is 150 km. The image appears stationary when the minimum speed of rotation of the octagonal mirror is 250 rotations per second. Calculate the velocity of light.

**Data :**

\[ D = 150 \text{ km} = 150 \times 10^3 \text{ m}; \quad n = 250 \text{ rps}; \quad N = 8; \quad C = ? \]

**Solution :**

In Michelson’s method, the velocity of light is

\[ C = NnD \]

\[ C = 8 \times 250 \times 150 \times 10^3 \]

\[ C = 3 \times 10^8 \text{ m s}^{-1} \]

9.4 The radii of curvature of two surfaces of a double convex lens are 10 cm each. Calculate its focal length and power of the lens in air and liquid. Refractive indices of glass and liquid are 1.5 and 1.8 respectively.

**Data :** \[ R_1 = 10 \text{ cm}, R_2 = -10 \text{ cm}; \quad \mu_g = 1.5 \text{ and } \mu_l = 1.8 \]

**Solution :** In air

\[ \frac{1}{f_a} = (\frac{\mu_g - 1}{R_1} - \frac{1}{R_2}) = (1.5 - 1) \left[ \frac{1}{10} - \frac{1}{10} \right] \]

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\[ f_a = 10 \text{ cm} \]
\[ P_a = \frac{1}{f_a} = \frac{1}{10 \times 10^2} \]
\[ P_a = 10 \text{ dioptres} \]

**In liquid**

\[ \frac{1}{f_l} = (\frac{\mu_g}{\mu_l} - 1) \left[ \frac{1}{R_1} - \frac{1}{R_2} \right] \]

\[ \frac{1}{f_l} = \left( \frac{\mu_g}{\mu_l} - 1 \right) \left( \frac{1}{R_1} - \frac{1}{R_2} \right) = \left( \frac{1.5}{1.8} - 1 \right) \left( \frac{1}{10} + \frac{1}{10} \right) = \frac{1}{6} \times \frac{2}{10} \]

\[ f_l = -30 \text{ cm} \]

\[ P_l = \frac{1}{f_l} = -\frac{1}{30 \times 10^{-2}} \]

\[ P_l = -3.33 \text{ dioptres} \]

9.5 A needle of size 5 cm is placed 45 cm from a lens produced an image on a screen placed 90 cm away from the lens. Identify the type of the lens and calculate its focal length and size of the image.

**Data:** \( h_1 = 5 \text{ cm}, \ u = -45 \text{ cm}, \ v = 90 \text{ cm}, \ f = ? \ h_2 = ? \)

**Solution:** We know that

\[ \frac{1}{f} - \frac{1}{v} - \frac{1}{u} = \frac{1}{90} - \frac{1}{-45} \]

\[ \therefore f = 30 \text{ cm} \]

Since \( f \) is positive, the lens is converging

Since \[ \frac{h_2}{h_1} = \frac{v}{u} \]

\[ \frac{90}{-45} = -2 \]

\[ \therefore h_2 = -10 \text{ cm} \]

(The –ve sign indicates that the image is real and inverted)
Self evaluation

(The questions and problems given in this self evaluation are only samples. In the same way any question and problem could be framed from the text matter. Students must be prepared to answer any question and problem from the text matter, not only from the self evaluation.)

9.1 The number of images of an object held between two parallel plane mirrors.
(a) infinity (b) 1
(c) 3 (d) 0

9.2 Radius of curvature of concave mirror is 40 cm and the size of image is twice as that of object, then the object distance is
(a) 20 cm (b) 10 cm
(c) 30 cm (d) 60 cm

9.3 A ray of light passes from a denser medium strikes a rarer medium at an angle of incidence $i$. The reflected and refracted rays are perpendicular to each other. The angle of reflection and refraction are $r$ and $r'$. The critical angle is
(a) $\tan^{-1}(\sin i)$ (b) $\sin^{-1}(\tan i)$
(c) $\tan^{-1}(\sin r)$ (d) $\sin^{-1}(\tan r')$

9.4 Light passes through a closed tube which contains a gas. If the gas inside the tube is gradually pumped out, the speed of light inside the tube
(a) increases (b) decreases
(c) remains constant (d) first increases and then decreases

9.5 In Michelson’s experiment, when the number of faces of rotating mirror increases, the velocity of light
(a) decreases (b) increases
(c) does not change (d) varies according to the rotation

9.6 If the velocity of light in a medium is $\frac{2}{3}$ times of the velocity of light in vacuum, then the refractive index of that medium is.
(a) $\frac{3}{2}c$ (b) $2c/3$
(c) $2/3$ (d) 1.5
9.7 Two lenses of power +12 and −2 dioptre are placed in contact. The focal length of the combination is given by
(a) 8.33 cm  (b) 12.5 cm
(c) 16.6 cm  (d) 10 cm

9.8 A converging lens is used to form an image on a screen. When the lower half of the lens is covered by an opaque screen then,
(a) half of the image will disappear
(b) complete image will be formed
(c) no image is formed
(d) intensity of the image is high

9.9 Two small angled prism of refractive indices 1.6 and 1.8 produced same deviation, for an incident ray of light, the ratio of angle of prism
(a) 0.88  (b) 1.33
(c) 0.56  (d) 1.12

9.10 Rainbow is formed due to the phenomenon of
(a) refraction and absorption
(b) dispersion and focussing
(c) refraction and scattering
(d) dispersion and total internal reflection

9.11 State the laws of reflection.

9.12 Show that the reflected ray turns by $2\theta$ when mirror turns by $\theta$.

9.13 Explain the image formation in plane mirrors.

9.14 Draw graphically the image formation in spherical mirrors with different positions of the object and state the nature of the image.

9.15 What is the difference between the virtual images produced by (i) plane mirror (ii) concave mirror (iii) convex mirror

9.16 The surfaces of the sun glasses are curved, yet their power may be zero. Why?
9.17 Prove the mirror formula for reflection of light from a concave mirror producing (i) real image (ii) virtual image.

9.18 With the help of ray diagram explain the phenomenon of total internal reflection. Give the relation between critical angle and refractive index.

9.19 Write a note on optical fibre.

9.20 Explain Michelson’s method of determining velocity of light.

9.21 Give the importance of velocity of light.

9.22 Derive lens maker’s formula for a thin biconvex lens.

9.23 Define power of a lens. What is one dioptre?

9.24 Establish the relation \( \frac{1}{F} = \frac{1}{f_1} + \frac{1}{f_2} \) of thin lenses in contact.

9.25 Derive the relation \( \mu = \frac{\sin \frac{A+D}{2}}{\sin \frac{A}{2}} \).

9.26 Does a beam of white light disperse through a hollow prism?

9.27 Derive an equation for dispersive power of a prism.

9.28 Describe a spectrometer.

9.29 Explain how will you determine the angle of the minimum deviation of a prism using spectrometer.

9.30 Write a note on formation of rainbows.

**Problems**

9.31 Light of wavelength 5000 Å falls on a plane reflecting surface. Calculate the wavelength and frequency of reflected light. For what angle of incidence, the reflected ray is normal to the incident ray?
9.32 At what distance from a convex mirror of focal length 2.5 m should a boy stand, so that his image has a height equal to half the original height?

9.33 In a Michelson’s experiment the distance travelled by the light between two reflections from the octagon rotating mirror is 4.8 km. Calculate the minimum speed of the mirror so that the image is formed at the non-rotating position.

9.34 If the refractive index of diamond be 2.5 and glass 1.5, then how faster does light travel in glass than in diamond?

9.35 An object of size 3 cm is kept at a distance of 14 cm from a concave lens of focal length 21 cm. Find the position of the image produced by the lens?

9.36 What is the focal length of a thin lens if the lens is in contact with 2.0 dioptre lens to form a combination lens which has a focal length of –80 cm?

9.37 A ray passes through an equilateral prism such that the angle of incidence is equal to the angle of emergence and the later is equal to 3/4 of the angle of prism. Find the angle of deviation.

9.38 The refractive indices of flint glass of equilateral prism for 400 nm and 700 nm are 1.66 and 1.61 respectively. Calculate the difference in angle of minimum deviation.

9.39 White light is incident on a small angled prism of angle 5°. Calculate the angular dispersion if the refractive indices of red and violet rays are 1.642 and 1.656 respectively.

9.40 A thin prism of refractive index 1.5 deviates a ray by a minimum angle of 5°. When it is kept immersed in oil of refractive index 1.25, what is the angle of minimum deviation?
Answers

9.1 (a)  9.2 (b)  9.3 (b)  
9.4 (a)  9.5 (c)  9.6 (d)  
9.7 (d)  9.8 (b)  9.9 (b)  
9.10 (d)  

9.31  5000 Å ; 6 × 10^{14} \text{ Hz} ; 45^\circ  
9.32  2.5 m  
9.33  7.8 \times 10^3 \text{ rps}  
9.34  1.66 \text{ times}  
9.35  \sim 8.4 \text{ cm}  
9.36  \sim 30.8 \text{ cm}  
9.37  30^\circ  
9.38  4^\circ  
9.39  0.07^\circ  
9.40  2^\circ
10. Magnetism

The word magnetism is derived from iron ore magnetite (Fe₃O₄), which was found in the island of magnesia in Greece. It is believed that the Chinese had known the property of the magnet even in 2000 B.C. and they used magnetic compass needle for navigation in 1100 AD. But it was Gilbert who laid the foundation for magnetism and had suggested that Earth itself behaves as a giant bar magnet. The field at the surface of the Earth is approximately $10^{-4}$ T and the field extends up to a height of nearly five times the radius of the Earth.

10.1 Earth's magnetic field and magnetic elements

A freely suspended magnetic needle at a point on Earth comes to rest approximately along the geographical north - south direction. This shows that the Earth behaves like a huge magnetic dipole with its magnetic poles near its geographical poles. Since the north pole of the magnetic needle approximately points towards geographic north ($N_0$) it is appropriate to call the magnetic pole near $N_0$ as the magnetic south pole of Earth ($S_m$). Also, the pole near $S_0$ is the magnetic north pole of the Earth ($N_m$). (Fig. 10.1)

The Earth’s magnetic field at any point on the Earth can be completely defined in terms of certain quantities called magnetic elements of the Earth, namely

(i) Declination or the magnetic variation $\theta$.
(ii) Dip or inclination $\delta$ and
(iii) The horizontal component of the Earth’s magnetic field $B_h$.
Causes of the Earth’s magnetism

The exact cause of the Earth’s magnetism is not known even today. However, some important factors which may be the cause of Earth’s magnetism are:

(i) Magnetic masses in the Earth.
(ii) Electric currents in the Earth.
(iii) Electric currents in the upper regions of the atmosphere.
(iv) Radiations from the Sun.
(v) Action of moon etc.

However, it is believed that the Earth’s magnetic field is due to the molten charged metallic fluid inside the Earth’s surface with a core of radius about 3500 km compared to the Earth’s radius of 6400 km.

10.1.1 Bar magnet

The iron ore magnetite which attracts small pieces of iron, cobalt, nickel etc. is a natural magnet. The natural magnets have irregular shape and they are weak. A piece of iron or steel acquires magnetic properties when it is rubbed with a magnet. Such magnets made out of iron or steel are artificial magnets. Artificial magnets can have desired shape and desired strength. If the artificial magnet is in the form of a rectangular or cylindrical bar, it is called a bar magnet.

10.1.2 Basic properties of magnets

(i) When the magnet is dipped in iron filings, they cling to the ends of the magnet. The attraction is maximum at the two ends of the magnet. These ends are called poles of the magnet.

(ii) When a magnet is freely suspended, it always points along north-south direction. The pole pointing towards geographic north is called north pole N and the pole which points towards geographic south is called south pole S.

(iii) Magnetic poles always exist in pairs. (i.e) isolated magnetic pole does not exist.

(iv) The magnetic length of a magnet is always less than its geometric length, because the poles are situated a little inwards from the free ends of the magnet. (But for the purpose of calculation the
geometric length is always taken as magnetic length.)

(v) Like poles repel each other and unlike poles attract each other. North pole of a magnet when brought near north pole of another magnet, we can observe repulsion, but when the north pole of one magnet is brought near south pole of another magnet, we observe attraction.

(vi) The force of attraction or repulsion between two magnetic poles is given by Coulomb’s inverse square law.

Note: In recent days, the concept of magnetic poles has been completely changed. The origin of magnetism is traced only due to the flow of current. But anyhow, we have retained the conventional idea of magnetic poles in this chapter. Pole strength is denoted by \( m \) and its unit is ampere metre.

**Magnetic moment**

Since any magnet has two poles, it is also called a magnetic dipole.

The magnetic moment of a magnet is defined as the product of the pole strength and the distance between the two poles.

If \( m \) is the pole strength of each pole and \( 2l \) is the distance between the poles, the magnetic moment

\[ \mathbf{M} = m \, (2 \mathbf{l}) \]

Magnetic moment is a vector quantity. It is denoted by \( \mathbf{M} \). Its unit is \( \text{A m}^2 \). Its direction is from south pole to north pole.

**Magnetic field**

Magnetic field is the space in which a magnetic pole experiences a force or it is the space around a magnet in which the influence of the magnet is felt.

**Magnetic induction**

Magnetic induction is the fundamental character of a magnetic field at a point.

Magnetic induction at a point in a magnetic field is the force experienced by unit north pole placed at that point. It is denoted by \( \mathbf{B} \). Its unit is \( \frac{\text{N}}{\text{A m}} \). It is a vector quantity. It is also called as magnetic flux density.
If a magnetic pole of strength $m$ placed at a point in a magnetic field experiences a force $F$, the magnetic induction at that point is

$$\vec{B} = \frac{\vec{F}}{m}$$

**Magnetic lines of force**

A magnetic field is better studied by drawing as many number of magnetic lines of force as possible.

*A magnetic line of force is a line along which a free isolated north pole would travel when it is placed in the magnetic field.*

**Properties of magnetic lines of force**

(i) Magnetic lines of forces are closed continuous curves, extending through the body of the magnet.

(ii) The direction of line of force is from north pole to south pole outside the magnet while it is from south pole to north pole inside the magnet.

(iii) The tangent to the magnetic line of force at any point gives the direction of magnetic field at that point. (i.e) it gives the direction of magnetic induction ($\vec{B}$) at that point.

(iv) They never intersect each other.

(v) They crowd where the magnetic field is strong and thin out where the field is weak.

**Magnetic flux and magnetic flux density**

The number of magnetic lines of force passing through an area $A$ is called magnetic flux. It is denoted by $\phi$. Its unit is weber. It is a scalar quantity.

The number of magnetic lines of force crossing unit area kept normal to the direction of line of force is magnetic flux density. Its unit is Wb m$^{-2}$ or tesla or N A$^{-1}$m$^{-1}$.

$$\therefore \text{Magnetic flux } \phi = \vec{B} \cdot \vec{A}$$

**Uniform and non-uniform magnetic field**

Magnetic field is said to be uniform if the magnetic induction has the same magnitude and the same direction at all points.
the points in the region. It is represented by drawing parallel lines (Fig. 10.2).

An example of uniform magnetic field over a wide area is the Earth's magnetic field.

If the magnetic induction varies in magnitude and direction at different points in a region, the magnetic field is said to be non-uniform. The magnetic field due to a bar magnet is non-uniform. It is represented by convergent or divergent lines (Fig. 10.3).

10.2 Force between two magnetic poles

In 1785, Coulomb made use of his torsion balance and discovered the law governing the force between the two magnetic poles.

Coulomb's inverse square law

Coulomb's inverse square law states that the force of attraction or repulsion between the two magnetic poles is directly proportional to the product of their pole strengths and inversely proportional to the square of the distance between them.

If \( m_1 \) and \( m_2 \) are the pole strengths of two magnetic poles separated by a distance of \( d \) in a medium, then

\[
F \propto m_1 m_2 \quad \text{and} \quad F \propto \frac{1}{d^2}
\]

\[
\therefore F \propto \frac{m_1 m_2}{d^2}
\]

\[
F = k \frac{m_1 m_2}{d^2}
\]

where \( k \) is the constant of proportionality and \( k = \frac{\mu}{4\pi} \) where \( \mu \) is the permeability of the medium.

But \( \mu = \mu_o \times \mu_r \)
∴ \( \mu_r = \frac{\mu}{\mu_0} \)

where \( \mu_r \) - relative permeability of the medium
\( \mu_0 \) - permeability of free space or vacuum.

Let \( m_1 = m_2 = 1 \)
and \( d = 1 \) m
\( k = \frac{\mu_0}{4\pi} \)

In free space, \( \mu_0 = 4\pi \times 10^{-7} \text{ H m}^{-1} \)

\[ F = \frac{10^{-7} \times m_1 \times m_2}{d^2} \]

\[ F = \frac{10^{-7} \times 1 \times 1}{1^2} \]
\[ F = 10^{-7} \text{ N} \]

Therefore, unit pole is defined as that pole which when placed at a distance of 1 metre in free space or air from an equal and similar pole, repels it with a force of \( 10^{-7} \text{ N} \).

10.3 Magnetic induction at a point along the axial line due to a magnetic dipole (Bar magnet)

\( NS \) is the bar magnet of length \( 2l \) and of pole strength \( m \).
\( P \) is a point on the axial line at a distance \( d \) from its mid point \( O \) (Fig. 10.4).

According to inverse square law, \( F = \frac{\mu_0}{4\pi} \frac{m_1 m_2}{d^2} \)

∴ Magnetic induction \( (B_1) \) at \( P \) due to north pole of the magnet,
\[ B_1 = \frac{\mu_0}{4\pi} \frac{m}{NP^2} \text{ along } NP \left( \therefore B = \frac{F}{m} \right) \]

\[ = \frac{\mu_0}{4\pi} \frac{m}{(d-l)^2} \text{ along } NP \]

Magnetic induction \( (B_2) \) at \( P \) due to south pole of the magnet,
\[ B_2 = \frac{\mu_0}{4\pi} \frac{m}{SP^2} \text{ along } PS \]
B_2 = \frac{\mu_o}{4\pi} \frac{m}{(d+l)^2} \text{ along } PS

\therefore \text{ Magnetic induction at } P \text{ due to the bar magnet,} \quad B = B_1 - B_2

B = \frac{\mu_o m}{4\pi} \left( \frac{1}{(d-l)^2} - \frac{1}{(d+l)^2} \right)

B = \frac{\mu_o m}{4\pi} \left( \frac{(d+l)^2 - (d-l)^2}{(d^2 - l^2)^2} \right)

B = \frac{\mu_o m}{4\pi} \left( \frac{4ld}{(d^2 - l^2)^2} \right)

B = \frac{\mu_o m}{4\pi} \left( \frac{2l \times 2d}{(d^2 - l^2)^2} \right)

B = \frac{\mu_o}{4\pi} \left( \frac{2M}{d^2 - l^2} \right)

\text{where } M = 2ml \text{ (magnetic dipole moment).}

\text{For a short bar magnet, } l \text{ is very small compared to } d, \text{ hence } l^2 \text{ is neglected.}\n
\therefore B = \frac{\mu_o}{4\pi} \left( \frac{2M}{d^3} \right)

\text{The direction of } B \text{ is along the axial line away from the north pole.}

\textbf{10.4 Magnetic induction at a point along the equatorial line of a bar magnet}

\text{NS is the bar magnet of length } 2l \text{ and pole strength } m. \text{ P is a point on the equatorial line at a distance } d \text{ from its mid point } O \text{ (Fig. 10.5).}
Magnetic induction \((B_1)\) at \(P\) due to north pole of the magnet,
\[
B_1 = \frac{\mu_0}{4\pi} \frac{m}{NP^2} \text{ along } NP
\]
\[
= \frac{\mu_0}{4\pi} \frac{m}{(d^2 + l^2)} \text{ along } NP
\]
\[
(\therefore NP^2 = NO^2 + OP^2)
\]
Magnetic induction \((B_2)\) at \(P\) due to south pole of the magnet,
\[
B_2 = \frac{\mu_0}{4\pi} \frac{m}{PS^2} \text{ along } PS
\]
\[
= \frac{\mu_0}{4\pi} \frac{m}{(d^2 + l^2)} \text{ along } PS
\]
Resolving \(B_1\) and \(B_2\) into their horizontal and vertical components.
Vertical components \(B_1 \sin \theta\) and \(B_2 \sin \theta\) are equal and opposite and therefore cancel each other (Fig. 10.6).

The horizontal components \(B_1 \cos \theta\) and \(B_2 \cos \theta\) will get added along \(PT\).

Resultant magnetic induction at \(P\) due to the bar magnet is
\[
B = B_1 \cos \theta + B_2 \cos \theta. \quad \text{(along } PT)\]
\[
B = \frac{\mu_0}{4\pi} \frac{m}{d^2 + l^2} \cdot \frac{l}{\sqrt{d^2 + l^2}} + \frac{\mu_0}{4\pi} \frac{m}{d^2 + l^2} \cdot \frac{l}{\sqrt{d^2 + l^2}}
\]
\[
= \frac{\mu_0}{4\pi} \frac{2ml}{(d^2 + l^2)^{3/2}}
\]
\[
B = \frac{\mu_0}{4\pi} \frac{M}{(d^2 + l^2)^{3/2}}. \quad \text{(where } M = 2ml)\]

For a short bar magnet, \(l^2\) is neglected.
\[
\therefore B = \frac{\mu_0}{4\pi} \frac{M}{d^3}
\]
The direction of \(B\) is along \(PT\) parallel to \(NS\).
10.5 Mapping of magnetic field due to a bar magnet

A bar magnet is placed on a plane sheet of a paper. A compass needle is placed near the north pole of the magnet. The north and south poles of the compass are marked by pencil dots. The compass needle is shifted and placed so that its south pole touches the pencil dot marked for north pole. The process is repeated and a series of dots are obtained. The dots are joined as a smooth curve. This curve is a magnetic line of force. Even though few lines are drawn around a bar magnet the magnetic lines exists in all space around the magnet.

(i) Magnet placed with its north pole facing geographic north

A sheet of paper is fixed on a drawing board. Using a compass needle, the magnetic meridian is drawn on it. A bar magnet is placed on the magnetic meridian such that its north pole points towards geographic north. Using a compass needle, magnetic lines of force are drawn around the magnet. (Fig. 10.7)

The magnetic lines of force is due to the combined effect of the magnetic field due to the bar magnet and Earth. It is found that when the compass is placed at points P and P' along the equatorial line of the magnet, the compass shows no deflection. They are called “neutral points.” At these points the magnetic field due to the magnet along its equatorial line (B) is exactly balanced by the horizontal component of the Earth’s magnetic field. (Bₕ)

Hence, neutral points are defined as the points where the resultant magnetic field due to the magnet and Earth is zero.

Hence, at neutral points

\[
\frac{\mu_0}{4\pi} \frac{M}{(d^2 + l^2)^{3/2}} = B_h
\]
(ii) Magnet placed with its south pole facing geographic north

A sheet of paper is fixed on a drawing board. Using a compass needle, the magnetic meridian is drawn on it. A bar magnet is placed on a magnetic meridian such that its south pole facing geographic north. Using a compass needle, the magnetic lines of force are drawn around the magnet as shown in Fig. 10.8.

The magnetic lines of force is due to the combined effect of the magnetic field due to the bar magnet and Earth. It is found that when the compass is placed at points \(P\) and \(P'\) along the axial line of the magnet, the compass shows no deflection. They are called neutral points. At these points the magnetic field \((B)\) due to the magnet along its axial line is exactly balanced by the horizontal component of the Earth’s magnetic field \((B_h)\).

Hence at neutral points, \(B = B_h\)

\[
\therefore \frac{\mu_0}{4\pi} \frac{2Md}{(d^2 - l^2)^2} = B_h
\]

10.6 Torque on a bar magnet placed in a uniform magnetic field

Consider a bar magnet NS of length \(2l\) and pole strength \(m\) placed in a uniform magnetic field of induction \(B\) at an angle \(\theta\) with the direction of the field (Fig. 10.9).

Due to the magnetic field \(B\), a force \(mB\) acts on the north pole along the direction of the field and a force \(mB\) acts on the south pole along the direction opposite to the magnetic field.
These two forces are equal and opposite, hence constitute a couple. The torque $\tau$ due to the couple is

$$\tau = \text{one of the forces} \times \text{perpendicular distance between them}$$

$$\tau = F \times NA$$

$$= mB \times NA$$

$$= mB \times 2l \sin \theta$$

$$\therefore \tau = MB \sin \theta \quad \text{...(2)}$$

Vectorially,

$$\tau = \vec{M} \times \vec{B}$$

The direction of $\tau$ is perpendicular to the plane containing $\vec{M}$ and $\vec{B}$.

If $B = 1$ and $\theta = 90^\circ$

Then from equation (2), $\tau = M$

Hence, moment of the magnet $M$ is equal to the torque necessary to keep the magnet at right angles to a magnetic field of unit magnetic induction.

10.7 Tangent law

A magnetic needle suspended, at a point where there are two crossed magnetic fields acting at right angles to each other, will come to rest in the direction of the resultant of the two fields.

$B_1$ and $B_2$ are two uniform magnetic fields acting at right angles to each other. A magnetic needle placed in these two fields will be subjected to two torques tending to rotate the magnet in opposite directions. The torque $\tau_1$ due to the two equal and opposite parallel forces $mB_1$ and $mB_1$ tend to set the magnet parallel to $B_1$. Similarly the torque $\tau_2$ due to the two equal and opposite parallel forces $mB_2$ and $mB_2$ tends to set the magnet parallel to $B_2$. In a position where the torques balance each other, the
magnet comes to rest. Now the magnet makes an angle $\theta$ with $B_2$ as shown in the Fig. 10.10.

The deflecting torque due to the forces $mB_1$ and $mB_1$

$$\tau_1 = mB_1 \times NA$$
$$= mB_1 \times NS \cos \theta$$
$$= mB_1 \times 2l \cos \theta$$
$$= 2l mB_1 \cos \theta$$

$$\therefore \tau_1 = MB_1 \cos \theta$$

Similarly the restoring torque due to the forces $mB_2$ and $mB_2$

$$\tau_2 = mB_2 \times SA$$
$$= mB_2 \times 2l \sin \theta$$
$$= 2lm \times B_2 \sin \theta$$
$$\tau_2 = MB_2 \sin \theta$$

At equilibrium,

$$\tau_1 = \tau_2$$

$$\therefore MB_1 \cos \theta = MB_2 \sin \theta$$

$$\therefore B_1 = B_2 \tan \theta$$

This is called Tangent law

Invariably, in the applications of tangent law, the restoring magnetic field $B_2$ is the horizontal component of Earth’s magnetic field $B_h$.

### 10.8 Deflection magnetometer

Deflection magnetometer consists of a small magnetic needle pivoted on a sharp support such that it is free to rotate in a horizontal plane. A light, thin, long aluminium pointer is fixed perpendicular to the magnetic needle. The pointer also rotates along with the needle (Fig. 10.11). There is a circular scale divided into four quadrants and each quadrant is graduated from $0^\circ$ to $90^\circ$. A plane mirror fixed below the scale ensures, reading without...
parallax error, as the image of the pointer is made to coincide exactly with pointer itself. The needle, aluminium pointer and the scale are enclosed in a box with a glass top. There are two arms graduated in centimetre and their zeroes coincide at the centre of the magnetic needle.

10.8.1 End-on (or) Tan A position

The magnetic field at a point along the axial line of a bar magnet is perpendicular to the horizontal component of Earth’s magnetic field. If a magnetometer and a bar magnet are placed in such way that this condition is satisfied, then this arrangement is called Tan A position.

To achieve this, the arms of the deflection magnetometer are placed along East-West direction (i.e) perpendicular to the magnetic meridian. The bar magnet is placed along East-West direction (i.e) parallel to the arms, as shown in the Fig. 10.12.

When a bar magnet of magnetic moment $M$ and length $2l$ is placed at a distance $d$ from the centre of the magnetic needle, the needle gets deflected through an angle $\theta$ due to the action of two magnetic fields.

(i) the field $B$ due to the bar magnet acting along its axis and 
(ii) the horizontal component of Earth’s magnetic field $B_h$.

The magnetic field at a distance $d$ acting along the axial line of the bar magnet,

$$B = \frac{\mu_0}{4\pi} \frac{2Md}{(d^2-l^2)^2}$$

According to Tangent law,

$$B = B_h \tan \theta$$

$$\frac{\mu_0}{4\pi} \frac{2Md}{(d^2-l^2)^2} = B_h \tan \theta$$

Comparison of magnetic moments of two bar magnets

(i) Deflection method

The deflection magnetometer is placed in Tan A position (Fig. 10.13). A bar magnet of magnetic moment $M_1$ and length $2l_1$ is placed at a distance
$d_1$ from the centre of the magnetic needle, on one side of the compass box. Since, the sensitivity of the magnetometer is more at $45^\circ$, the distance of the bar magnet should be chosen such that the deflection lies between $30^\circ$ and $60^\circ$. The readings corresponding to the ends of the aluminium pointer are noted as $\theta_1$ and $\theta_2$. The magnet is reversed pole to pole and kept at the same distance. Two more readings $\theta_3$ and $\theta_4$ are noted. By placing the magnet on the other side of the compass box at the same distance, four more readings $\theta_5$, $\theta_6$, $\theta_7$, and $\theta_8$ are noted as above. The mean of the eight readings gives a value $\theta_I$.

The experiment is repeated as above for the second bar magnet of magnetic moment $M_2$ and length $2l_2$ by placing at a distance $d_2$. Now the mean of the eight readings gives a value of $\theta_{II}$.

Applying tangent law, for the first magnet,

$$\frac{\mu_0}{4\pi} \frac{2M_1d_1}{(d_1^2 - l_1^2)^2} = B_h \tan \theta_I$$

...(1)

and for the second magnet.

$$\frac{\mu_0}{4\pi} \frac{2M_2d_2}{(d_2^2 - l_2^2)^2} = B_h \tan \theta_{II}$$

...(2)

From the above equations (1) and (2), we get

$$\frac{M_1}{M_2} = \frac{(d_2^2 - l_2^2)^2}{(d_1^2 - l_1^2)^2} \frac{d_2}{d_1} \frac{\tan \theta_I}{\tan \theta_{II}}$$

...(3)

Special case

If the magnets are placed at the same distance, then $d_1 = d_2 = d$

$$\frac{M_1}{M_2} = \frac{(d^2 - l_1^2)^2}{(d^2 - l_2^2)^2} \frac{\tan \theta_I}{\tan \theta_{II}}$$

In addition, if $l_1$ and $l_2$ are small compared to the distance $d$

then $\frac{M_1}{M_2} = \frac{\tan \theta_I}{\tan \theta_{II}}$
(ii) Null deflection method

The deflection magnetometer is placed in Tan A position (Fig. 10.14). A bar magnet of magnetic moment $M_1$ and length $2l_1$ is placed on one side of the compass box at a distance $d_1$ from the centre of the magnetic needle. The second bar magnet of magnetic moment $M_2$ and length $2l_2$ is placed on the other side of the compass box such that like poles of the magnets face each other. The second magnet is adjusted so that the deflection due to the first magnet is nullified and the aluminium pointer reads $0^\circ - 0^\circ$. The distance of the second magnet is $x_1$. The first magnet is reversed pole to pole and placed at the same distance $d_1$. The second magnet is also reversed and adjusted such that the aluminium pointer reads $0^\circ - 0^\circ$. The distance of the second magnet is $x_2$.

The experiment is repeated by interchanging the magnets. Two more distances $x_3$ and $x_4$ are noted. The mean of $x_1$, $x_2$, $x_3$ and $x_4$ is taken as $d_2$.

As the magnetic fields due to the two bar magnets at the centre of the magnetic needle are equal in magnitude but opposite in direction, (i.e) $B_1 = B_2$

$$\frac{\mu_0}{4\pi} \frac{2M_1d_1}{(d_1^2 - l_1^2)^2} = \frac{\mu_0}{4\pi} \frac{2M_2d_2}{(d_2^2 - l_2^2)^2}$$

$$\therefore \frac{M_1}{M_2} = \frac{(d_2^2 - l_2^2)^2}{(d_1^2 - l_1^2)^2} \frac{d_2}{d_1}$$

If the bar magnets are short, $l_1$ and $l_2$ are negligible compared to the distance $d_1$ and $d_2$

$$\therefore \frac{M_1}{M_2} = \frac{d_1^3}{d_2^3}$$
10.8.2 Broad–side on (or) Tan B position

The magnetic field at a point along the equatorial line of a bar magnet is perpendicular to the horizontal component of Earth's magnetic field. If the magnetometer and a bar magnet are placed in such way that this condition is satisfied, then this arrangement is called Tan B position.

To achieve this, the arms of the deflection magnetometer are placed along the North - South direction (i.e) along the magnetic meridian. The magnet is placed along East - West direction (i.e) parallel to the aluminium pointer as shown in the Fig. 10.15.

When a bar magnet of magnetic moment \( M \) and length \( 2l \) is placed at a distance \( d \) from the centre of the magnetic needle, the needle gets deflected through an angle \( \theta \) due to the action of the following two magnetic fields.

(i) The field \( B \) due to the bar magnet along its equatorial line
(ii) The horizontal component of Earth’s magnetic field \( B_h \).

The magnetic field at a distance \( d \) along the equatorial line of the bar magnet,

\[
B = \frac{\mu_0 M}{4\pi (d^2 + l^2)^{3/2}}
\]

According to tangent law

\[
B = B_h \tan \theta
\]

(i.e) \( \frac{\mu_0}{4\pi} \frac{M}{(d^2 + l^2)^{3/2}} = B_h \tan \theta \)

If the magnet is short, \( l \) is small compared to \( d \) and hence \( l^2 \) is neglected.

\[
\frac{\mu_0}{4\pi} \frac{M}{d^3} = B_h \tan \theta
\]
Comparison of magnetic moments of two bar magnets

(i) Deflection method

The deflection magnetometer is placed in Tan B position. A bar magnet of magnetic moment $M_1$ and length $2l_1$ is placed at a distance $d_1$ from the centre of the magnetic needle, on one side of the compass box (Fig. 10.16). Since, the sensitivity of the magnetometer is more at 45°, the distance of the bar magnet should be chosen such that the deflection lies between 30° and 60°. The readings corresponding to the ends of the aluminium pointer are noted as $\theta_1$ and $\theta_2$. The magnet is reversed pole to pole and kept at the same distance. Two more readings $\theta_3$ and $\theta_4$ are noted. By placing the magnet on the other side of the compass box at the same distance, four more readings $\theta_5$, $\theta_6$, $\theta_7$ and $\theta_8$ are noted as above. The mean of the eight readings gives a value $\theta_i$.

The experiment is repeated as above for the second bar magnet of magnetic moment $M_2$ and length $2l_2$ by placing at a distance $d_2$. Now the mean of the eight readings gives a value of $\theta_{II}$.

Applying tangent law, for the first magnet,

$$\frac{\mu_0}{4\pi} \frac{M_1}{(d_1^2 + l_1^2)^{3/2}} = B_h \tan \theta_i \tag{1}$$

and for the second magnet

$$\frac{\mu_0}{4\pi} \frac{M_2}{(d_2^2 + l_2^2)^{3/2}} = B_h \tan \theta_{II} \tag{2}$$

From the above equations (1) and (2), we get

$$\frac{M_1}{M_2} = \frac{\left(\frac{d_1^2 + l_1^2}{d_2^2 + l_2^2}\right)^{3/2} \tan \theta_i}{\tan \theta_{II}} \tag{3}$$

Special case

If the magnets are placed at the same distance, then $d_1 = d_2 = d$
\[
\frac{M_1}{M_2} = \frac{(d^2 + l_1^2)^{3/2}}{(d^2 + l_2^2)^{3/2}} \cdot \tan \theta_I \cdot \tan \theta_H
\]

In addition, if \( l_1 \) and \( l_2 \) are small compared to the distance \( d \),

\[
\frac{M_1}{M_2} = \frac{\tan \theta_I}{\tan \theta_H}
\]

(ii) Null deflection method

The deflection magnetometer is placed in Tan B position (Fig. 10.17). A bar magnet of magnetic moment \( M_1 \) and length \( 2l_1 \) is placed on one side of the compass box at a distance \( d_1 \) from the centre of the magnetic needle. The second bar magnet of magnetic moment \( M_2 \) and length \( 2l_2 \) is placed on the other side of the compass box such that like poles of the magnets face in the opposite direction. The second magnet is adjusted so that the deflection due to the first magnet is nullified and the aluminium pointer reads \( 0^\circ \) - \( 0^\circ \). The distance of the second magnet is \( x_1 \). The first magnet is reversed pole to pole and placed at the same distance \( d_1 \). The second magnet is also reversed and adjusted such that the aluminium pointer reads \( 0^\circ \) - \( 0^\circ \). The distance of the second magnet is \( x_2 \).

The experiment is repeated by interchanging the magnets. Two more distances \( x_3 \) and \( x_4 \) are noted. The mean of \( x_1, x_2, x_3 \) and \( x_4 \) is taken as \( d_2 \).

Since the magnetic fields due to the two bar magnets at the centre of the magnetic needle are equal in magnitude but opposite in direction,

\[
\therefore \quad B_1 = B_2
\]

\[
\frac{\mu_0}{4\pi} \frac{M_1}{(d_1^2 + l_1^2)^{3/2}} = \frac{\mu_0}{4\pi} \frac{M_2}{(d_2^2 + l_2^2)^{3/2}}
\]

\[
\therefore \quad \frac{M_1}{M_2} = \frac{(d_1^2 + l_1^2)^{3/2}}{(d_2^2 + l_2^2)^{3/2}}
\]
If the bar magnets are short, \( l_1 \) and \( l_2 \) are negligible compared to the distance \( d_1 \) and \( d_2 \)

\[
\frac{M_1}{M_2} = \frac{d_1^3}{d_2^3}
\]

10.9 Magnetic properties of materials

The study of magnetic properties of materials assumes significance since these properties decide whether the material is suitable for permanent magnets or electromagnets or cores of transformers etc. Before classifying the materials depending on their magnetic behaviour, the following important terms are defined.

(i) Magnetising field or magnetic intensity

The magnetic field used to magnetise a material is called the magnetising field. It is denoted by \( H \) and its unit is \( \text{A m}^{-1} \).

(Note: Since the origin of magnetism is linked to the current, the magnetising field is usually defined in terms of ampere turn which is out of our purview here.)

(ii) Magnetic permeability

Magnetic permeability is the ability of the material to allow the passage of magnetic lines of force through it.

Relative permeability \( \mu_r \) of a material is defined as the ratio of number of magnetic lines of force per unit area \( B \) inside the material to the number of lines of force per unit area in vacuum \( B_0 \) produced by the same magnetising field.

\[
\therefore \text{Relative permeability } \mu_r = \frac{B}{B_0}
\]

\[
\mu_r = \frac{\mu H}{\mu_0 H} = \frac{\mu}{\mu_0}
\]

(since \( \mu_r \) is the ratio of two identical quantities, it has no unit.)

\[
\therefore \text{The magnetic permeability of the medium } \mu = \mu_0 \mu_r \text{ where } \mu_0 \text{ is the permeability of free space.}
\]

Magnetic permeability \( \mu \) of a medium is also defined as the ratio of magnetic induction \( B \) inside the medium to the magnetising field \( H \) inside the same medium.

\[
\therefore \mu = \frac{B}{H}
\]
(iii) Intensity of magnetisation

Intensity of magnetisation represents the extent to which a material has been magnetised under the influence of magnetising field \( H \).

*Intensity of magnetisation of a magnetic material is defined as the magnetic moment per unit volume of the material.*

\[
I = \frac{M}{V}
\]

Its unit is A m\(^{-1}\).

For a specimen of length \( 2l \), area \( A \) and pole strength \( m \),

\[
I = \frac{2lm}{2lA}
\]

\[
\therefore I = \frac{m}{A}
\]

Hence, intensity of magnetisation is also defined as the pole strength per unit area of the cross section of the material.

(iv) Magnetic induction

When a soft iron bar is placed in a uniform magnetising field \( H \), the magnetic induction inside the specimen \( B \) is equal to the sum of the magnetic induction \( B_o \) produced in vacuum due to the magnetising field and the magnetic induction \( B_m \) due to the induced magnetisation of the specimen.

\[
B = B_o + B_m
\]

But \( B_o = \mu_o H \) and \( B_m = \mu_o I \)

\[
B = \mu_o H + \mu_o I
\]

\[
\therefore B = \mu_o (H + I)
\]

(v) Magnetic susceptibility

Magnetic susceptibility \( \chi_m \) is a property which determines how easily and how strongly a specimen can be magnetised.

*Susceptibility of a magnetic material is defined as the ratio of intensity of magnetisation \( I \) induced in the material to the magnetising field \( H \) in which the material is placed.*
Thus $\chi_m = \frac{I}{H}$

Since $I$ and $H$ are of the same dimensions, $\chi_m$ has no unit and is dimensionless.

**Relation between $\chi_m$ and $\mu_r$**

$$\chi_m = \frac{I}{H}$$

$\therefore I = \chi_m H$

We know

$B = \mu_o (H + I)$

$B = \mu_o (H + \chi_m H)$

$B = \mu_o H (1 + \chi_m)$

If $\mu$ is the permeability, we know that $B = \mu H$.

$\therefore \mu H = \mu_o H (1 + \chi_m)$

$$\frac{\mu}{\mu_o} = (1 + \chi_m)$$

$\therefore \mu_r = 1 + \chi_m$

### 10.10 Classification of magnetic materials

On the basis of the behaviour of materials in a magnetising field, the materials are generally classified into three categories namely, (i) Diamagnetic, (ii) Paramagnetic and (iii) Ferromagnetic

**Properties of diamagnetic substances**

Diamagnetic substances are those in which the net magnetic moment of atoms is zero.

1. The susceptibility has a low negative value. (For example, for bismuth $\chi_m = -0.00017$).

2. Susceptibility is independent of temperature.

3. The relative permeability is slightly less than one.

4. When placed in a non uniform magnetic field they have a tendency to move...
away from the field. (i.e) from the stronger part to the weaker part of the field. They get magnetised in a direction opposite to the field as shown in the Fig. 10.18.

5. When suspended freely in a uniform magnetic field, they set themselves perpendicular to the direction of the magnetic field (Fig. 10.19).

Examples: Bi, Sb, Cu, Au, Hg, H₂O, H₂ etc.

(ii) Properties of paramagnetic substances

Paramagnetic substances are those in which each atom or molecule has a net non-zero magnetic moment of its own.

1. Susceptibility has a low positive value.
   (For example: \( \chi_m \) for aluminium is +0.00002).
2. Susceptibility is inversely proportional to absolute temperature
   (i.e) \( \chi_m \propto \frac{1}{T} \). As the temperature increases susceptibility decreases.
3. The relative permeability is greater than one.
4. When placed in a non-uniform magnetic field, they have a tendency to move from weaker part to the stronger part of the field. They get magnetised in the direction of the field as shown in Fig. 10.20.

5. When suspended freely in a uniform magnetic field, they set themselves parallel to the direction of magnetic field (Fig. 10.21).

Examples: Al, Pt, Cr, O₂, Mn, CuSO₄ etc.
(iii) Properties of ferromagnetic substances

Ferromagnetic substances are those in which each atom or molecule has a strong spontaneous net magnetic moment. These substances exhibit strong paramagnetic properties.

1. The susceptibility and relative permeability are very large.
   (For example: \( \mu_r \) for iron = 200,000)
2. Susceptibility is inversely proportional to the absolute temperature.
   \[ \chi \propto \frac{1}{T} \]
   As the temperature increases the value of susceptibility decreases. At a particular temperature, ferromagnetics become paramagnetics. This transition temperature is called the Curie temperature. For example, the Curie temperature of iron is about 1000 K.
3. When suspended freely in uniform magnetic field, they set themselves parallel to the direction of magnetic field.
4. When placed in a non-uniform magnetic field, they have a tendency to move from the weaker part to the stronger part of the field. They get strongly magnetised in the direction of the field.

Examples: Fe, Ni, Co and a number of their alloys.

10.11 Hysteresis

Consider an iron bar being magnetised slowly by a magnetising field \( H \) whose strength can be changed. It is found that the magnetic induction \( B \) inside the material increases with the strength of the magnetising field and then attains a saturated level. This is depicted by the path OP in the Fig. 10.22.

If the magnetising field is now decreased slowly, then magnetic induction also decreases but it does not follow the path PO. Instead, when \( H = 0 \), \( B \) has non zero value equal to \( OQ \). This implies that some
magnetism is left in the specimen. The value of magnetic induction of a substance, when the magnetising field is reduced to zero, is called remanance or residual magnetic induction of the material. OQ represents the residual magnetism of the material. Now, if we apply the magnetising field in the reverse direction, the magnetic induction decreases along QR till it becomes zero at R. Thus to reduce the residual magnetism (remanent magnetism) to zero, we have to apply a magnetising field OR in the opposite direction.

The value of the magnetising field $H$ which has to be applied to the magnetic material in the reverse direction so as to reduce its residual magnetism to zero is called its coercivity.

When the strength of the magnetising field $H$ is further increased in the reverse direction, the magnetic induction increases along RS till it acquires saturation at a point S (points P and S are symmetrical). If we now again change the direction of the field, the magnetic induction follows the path STUP. This closed curve PQRSTUP is called the ‘hysteresis loop’ and it represents a cycle of magnetisation. The word ‘hysteresis’ literally means lagging behind. We have seen that magnetic induction $B$ lags behind the magnetising field $H$ in a cycle of magnetisation. *This phenomenon of lagging of magnetic induction behind the magnetising field is called hysteresis.*

**Hysteresis loss**

In the process of magnetisation of a ferromagnetic substance through a cycle, there is expenditure of energy. The energy spent in magnetising a specimen is not recoverable and there occurs a loss of energy in the form of heat. This is so because, during a cycle of magnetisation, the molecular magnets in the specimen are oriented and reoriented a number of times. This molecular motion results in the production of heat. It has been found that *loss of heat energy per unit volume of the specimen in each cycle of magnetisation is equal to the area of the hysteresis loop.*

The shape and size of the hysteresis loop is characteristic of each material because of the differences in their retentivity, coercivity, permeability, susceptibility and energy losses etc. By studying hysteresis loops of various materials, one can select suitable materials for different purposes.
10.11.1 Uses of ferromagnetic materials

(i) Permanent magnets

The ideal material for making permanent magnets should possess high retentivity (residual magnetism) and high coercivity so that the magnetisation lasts for a longer time. Examples of such substances are steel and alnico (an alloy of Al, Ni and Co).

(ii) Electromagnets

Material used for making an electromagnet has to undergo cyclic changes. Therefore, the ideal material for making an electromagnet has to be one which has the least hysteresis loss. Moreover, the material should attain high values of magnetic induction $B$ at low values of magnetising field $H$. Soft iron is preferred for making electromagnets as it has a thin hysteresis loop (Fig. 10.23) [small area, therefore less hysteresis loss] and low retentivity. It attains high values of $B$ at low values of magnetising field $H$.

(iii) Core of the transformer

A material used for making transformer core and choke is subjected to cyclic changes very rapidly. Also, the material must have a large value of magnetic induction $B$. Therefore, soft iron that has thin and tall hysteresis loop is preferred. Some alloys with low hysteresis loss are: radio-metals, pern-alloy and mumetal.

(iv) Magnetic tapes and memory store

Magnetisation of a magnet depends not only on the magnetising field but also on the cycle of magnetisation it has undergone. Thus, the value of magnetisation of the specimen is a record of the cycles of magnetisation it has undergone. Therefore, such a system can act as a device for storing memory.

Ferro magnetic materials are used for coating magnetic tapes in a cassette player and for building a memory store in a modern computer. Examples: Ferrites (Fe, Fe$_2$O, MnFe$_2$O$_4$ etc.).
Solved Problems

10.1 A short bar magnet is placed with its north pole pointing north. The neutral point is 10 cm away from the centre of the magnet. If \( B = 4 \times 10^{-5} \text{T} \), calculate the magnetic moment of the magnet.

**Data :** \( d = 10 \times 10^{-2} \text{m} \); \( B = 4 \times 10^{-5} \text{ T} \); \( M = ? \)

**Solution :** When the north pole of a bar magnet points north, the neutral points will lie on its equatorial line.

\[ B = \mu_0 \frac{M}{4\pi d^3} \]

\[ \therefore M = B \times d^3 \times 10^{-7} = 4 \times 10^{-5} (10 \times 10^{-2})^3 \times 10^{-7} \]

\[ M = 0.4 \text{ A m}^2 \]

10.2 A bar magnet is suspended horizontally by a torsionless wire in magnetic meridian. In order to deflect the magnet through 30° from the magnetic meridian, the upper end of the wire has to be rotated by 270°. Now this magnet is replaced by another magnet. In order to deflect the second magnet through the same angle from the magnetic meridian, the upper end of the wire has to be rotated by 180°. What is the ratio of the magnetic moments of the two bar magnets. (Hint : \( \tau = C\theta \))

**Solution :** Let \( C \) be the deflecting torque per unit twist and \( M_1 \) and \( M_2 \) be the magnetic moments of the two magnets.

The deflecting torque is \( \tau = C\theta \)

The restoring torque is \( \tau = MB \sin \theta \)

In equilibrium

deflecting torque = restoring torque

For the Magnet – I

\[ C (270° - 30°) = M_1 B_h \sin \theta \] \[ \ldots (1) \]

For the magnet – II

\[ C (180° - 30°) = M_2 B_h \sin \theta \] \[ \ldots (2) \]
Dividing (1) by (2)

\[
\frac{M_1}{M_2} = \frac{240^\circ}{150^\circ} = \frac{8}{5}
\]

10.3 A short bar magnet of magnetic moment \( 5.25 \times 10^{-2} \) A m\(^2\) is placed with its axis perpendicular to the Earth’s field direction. At what distance from the centre of the magnet on (i) its equatorial line and (ii) its axial line, is the resultant field inclined at \( 45^\circ \) with the Earth’s field. Magnitude of the Earth’s field at the place is \( 0.42 \times 10^{-4} \) T.

**Data :**

- \( M = 5.25 \times 10^{-2} \) A m\(^2\)
- \( \theta = 45^\circ \)
- \( B_h = 0.42 \times 10^{-4} \) T
- \( d = ? \)

**Solution :** From Tangent Law

\[
\frac{B}{B_h} = \tan \theta
\]

\[
B = B_h \tan \theta = 0.42 \times 10^{-4} \times \tan 45^\circ
\]

\[
B = 0.42 \times 10^{-4} \text{ T}
\]

(i) For the point on the equatorial line

\[
B = \frac{\mu_0}{4\pi} \frac{M}{d^3}
\]

\[
d^3 = \frac{\mu_0}{4\pi} \frac{M}{B}
\]

\[
d^3 = \frac{4\pi \times 10^{-7} \times 5.25 \times 10^{-2}}{4\pi \times 0.42 \times 10^{-4}}
\]

\[
= 12.5 \times 10^{-5} = 125 \times 10^{-6}
\]

\[\therefore d = 5 \times 10^{-2} \text{ m}\]

(ii) For the point on the axial line

\[
B = \frac{\mu_0}{4\pi} \frac{2M}{d^3} \quad \text{(or)} \quad d^3 = \frac{\mu_0}{4\pi} \frac{2M}{B}
\]

\[
d^3 = \frac{4\pi \times 10^{-7}}{4\pi} \times \frac{2 \times 5.25 \times 10^{-2}}{0.42 \times 10^{-4}}
\]

\[
d^3 = 250 \times 10^{-6} = 2 \times 125 \times 10^{-6}
\]
\[ d = 2^{1/3} \times (5 \times 10^{-2}) \]
\[ d = 6.3 \times 10^{-2} \text{ m}. \]

10.4 A bar magnet of mass 90 g has magnetic moment 3 A\( m^2 \). If
the intensity of magnetisation of the magnet is \( 2.7 \times 10^5 \text{ A m}^{-1} \), find the density of the material of the magnet.

Data : \( m = 90 \times 10^{-3} \text{ kg}; M = 3 \text{ A m}^2 \)
\[ I = 2.7 \times 10^5 \text{ A m}^{-1}; \rho = ? \]

Solution : Intensity of magnetisation, \( I = \frac{M}{V} \)

But, volume \( V = \frac{m}{\rho} \)

\[ I = \frac{M\rho}{m} \]

\[ \rho = \frac{IM}{M} = \frac{2.7 \times 10^5 \times 90 \times 10^{-3}}{3} = 8100 \]

\[ \rho = 8100 \text{ kg m}^{-3} \]

10.5 A magnetising field of 50 A\( m^{-1} \) produces a magnetic field of
induction 0.024 T in a bar of length 8 cm and area of cross
section 1.5 cm\(^2\). Calculate (i) the magnetic permeability (ii) the
magnetic susceptibility.

Data : \( H = 50 \text{ A m}^{-1}, B = 0.024 \text{ T} = 2.4 \times 10^{-2} \text{ T}, \)
\[ 2l = 8 \times 10^{-2} \text{ m}, \quad A = 1.5 \times 10^{-4} \text{ m}^2 \]
\[ \mu = ?; \chi_m = ? \]

Solution : Permeability \( \mu = \frac{B}{H} = \frac{2.4 \times 10^{-2}}{50} = 4.8 \times 10^{-4} \text{ H m}^{-1} \)

Susceptibility, \( \chi_m = \mu_r - 1 = \frac{\mu}{\mu_0} - 1 \)

\[ \chi_m = \frac{4.8 \times 10^{-4}}{4\pi \times 10^{-7}} - 1 = 381.16 \]
Self evaluation

(The questions and problems given in this self evaluation are only samples. In the same way any question and problem could be framed from the text matter. Students must be prepared to answer any question and problem from the text matter, not only from the self evaluation.)

10.1 Two magnetic poles kept separated by a distance \( d \) in vacuum experience a force of 10 N. The force they would experience when kept inside a medium of relative permeability 2, separated by the same distance is

- (a) 20 N
- (b) 10 N
- (c) 5 N
- (d) 40 N

10.2 The magnetic moment of a magnet is 5 A m\(^2\). If the pole strength is 25 A m, what is the length of the magnet?

- (a) 10 cm
- (b) 20 cm
- (c) 25 cm
- (d) 1.25 cm

10.3 A long magnetic needle of length \( 2l \), magnetic moment \( M \) and pole strength \( m \) is broken into two at the middle. The magnetic moment and pole strength of each piece will be

- (a) \( M, m \)
- (b) \( \frac{M}{2}, \frac{m}{2} \)
- (c) \( M, \frac{m}{2} \)
- (d) \( \frac{M}{2}, m \)

10.4 Two short magnets have equal pole strengths but one is twice as long as the other. The shorter magnet is placed 20 cm in tan A position from the compass needle. The longer magnet must be placed on the other side of the magnetometer for zero deflection at a distance

- (a) 20 cm
- (b) 20 \((2)^{1/3}\) cm
- (c) 20 \((2)^{2/3}\) cm
- (d) 20 \((2)^{1}\) cm
10.5 The direction of a magnet in tan B position of a deflection magnetometer is
(a) North − South  (b) East − West
(c) North − West   (d) South − West

10.6 The relative permeability of a specimen is 10001 and magnetising field strength is 2500 A m\(^{-1}\). The intensity of magnetisation is
(a) \(0.5 \times 10^{-7}\) A m\(^{-1}\)  (b) \(2.5 \times 10^{-7}\) A m\(^{-1}\)
(c) \(2.5 \times 1.0^{+7}\) A m\(^{-1}\)  (d) \(2.5 \times 10^{-1}\) A m\(^{-1}\)

10.7 For which of the following substances, the magnetic susceptibility is independent of temperature?
(a) diamagnetic
(b) paramagnetic
(c) ferromagnetic
(d) diamagnetic and paramagnetic

10.8 At curie point, a ferromagnetic material becomes
(a) non-magnetic  (b) diamagnetic
(c) paramagnetic  (d) strongly ferromagnetic

10.9 Electromagnets are made of soft iron because soft iron has
(a) low susceptibility and low retentivity
(b) high susceptibility and low retentivity
(c) high susceptibility and high retentivity
(d) low permeability and high retentivity

10.10 State Coulomb’s inverse square law.

10.11 Obtain the expressions for the magnetic induction at a point on the (i) axial line and (ii) equatorial line of a bar magnet.

10.12 Find the torque experienced by a magnetic needle in a uniform magnetic field.

10.13 State and prove tangent law.
10.14 What is tan A position? How will you set up the deflection magnetometer in tan A position?

10.15 Explain the theory of tan A position. Explain how will you compare the magnetic moments of two bar magnets in this position.

10.16 What is tan B position? How will you set up the deflection magnetometer in tan B position?

10.17 Explain the theory of tan B position. Explain how will you compare the magnetic moments of two bar magnets in this position.

10.18 Define the terms (i) magnetic permeability (ii) intensity of magnetisation and (iii) magnetic susceptibility.

10.19 Distinguish between dia, para and ferro magnetic substances. Give one example for each.

10.20 Explain the hysteresis cycle.

Problems

10.21 The force acting on each pole of a magnet placed in a uniform magnetic induction of $5 \times 10^{-4}$ T is $6 \times 10^{-3}$ N. If the length of the magnet is 8 cm, calculate the magnetic moment of the magnet.

10.22 Two magnetic poles, one of which is twice stronger than the other, repel one another with a force of $2 \times 10^{-5}$ N, when kept separated at a distance of 20 cm in air. Calculate the strength of each pole.

10.23 Two like poles of unequal pole strength are placed 1 m apart. If a pole of strength 4 A m is in equilibrium at a distance 0.2 m from one of the poles, calculate the ratio of the pole strengths of the two poles.

10.24 A magnet of pole strength $24.6 \times 10^{-2}$ A m and length 10 cm is placed at $30^\circ$ with a magnetic field of 0.01 T. Find the torque acting on the magnet.

10.25 The magnetic moment of a bar magnet of length 10 cm is $9.8 \times 10^{-1}$ A m$^2$. Calculate the magnetic field at a point on its axis at a distance of 20 cm from its midpoint.
10.26 Two mutually perpendicular lines are drawn on a table. Two small magnets of magnetic moments 0.108 and 0.192 A m$^2$ respectively are placed on these lines. If the distance of the point of intersection of these lines is 30 cm and 40 cm respectively from these magnets, find the resultant magnetic field at the point of intersection.

10.27 The intensity of magnetisation of an iron bar of mass 72 g, density 7200 kg m$^{-3}$ is 0.72 A m$^{-1}$. Calculate the magnetic moment.

10.28 A magnet of volume 25 cm$^3$ has a magnetic moment of $12.5 \times 10^{-4}$ A m$^2$. Calculate the intensity of magnetisation.

10.29 A magnetic intensity of $2 \times 10^3$ A/m produces a magnetic induction of $4\pi$ Wb/m$^2$ in a bar of iron. Calculate the relative permeability and susceptibility.
**Answers**

<table>
<thead>
<tr>
<th>Question</th>
<th>Answer</th>
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<tbody>
<tr>
<td>10.1</td>
<td>(a)</td>
</tr>
<tr>
<td>10.2</td>
<td>(b)</td>
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<tr>
<td>10.3</td>
<td>(d)</td>
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<tr>
<td>10.4</td>
<td>(b)</td>
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<tr>
<td>10.5</td>
<td>(b)</td>
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<tr>
<td>10.6</td>
<td>(c)</td>
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<tr>
<td>10.7</td>
<td>(a)</td>
</tr>
<tr>
<td>10.8</td>
<td>(c)</td>
</tr>
<tr>
<td>10.9</td>
<td>(b)</td>
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<tr>
<td>10.21</td>
<td>0.96 A m²</td>
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<tr>
<td>10.22</td>
<td>2 A m, 4 A m</td>
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<td>1 :16</td>
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<td>10.24</td>
<td>$1.23 \times 10^{-4}$ N m</td>
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<tr>
<td>10.25</td>
<td>$2.787 \times 10^{-5}$ T</td>
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<tr>
<td>10.26</td>
<td>$10^{-6}$ T</td>
</tr>
<tr>
<td>10.27</td>
<td>$7.2 \times 10^{-6}$ A m²</td>
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<td>10.28</td>
<td>50 A m⁻¹</td>
</tr>
<tr>
<td>10.29</td>
<td>5000, 4999</td>
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1 Declination

A vertical plane passing through the axis of a freely suspended magnetic needle is called magnetic meridian and the vertical plane passing through the geographic north-south direction (axis of rotation of Earth) is called geographic meridian (Fig.).

In the Fig. 1 the plane $PQRS$ represents the magnetic meridian and the plane $PQR'S'$ represents the geographic meridian.

Declination at a place is defined as the angle between magnetic meridian and the geographic meridian at that place.

Determination of declination

A simple method of determining the geographical meridian at a place is to erect a pole of 1 to 1.5 m height on the ground and a circle is drawn with the pole O as centre and its height as radius as shown in the Fig. 2.

Mark a point $P_1$ on the circle before noon, when the tip of the shadow of the pole just touches the circle.

Again mark a point $P_2$ when the tip of the shadow touches the circle in the afternoon. The line $POQ$ drawn bisecting the angle $P_1OP_2$ is the geographical meridian at that place.

Magnetic meridian is drawn by freely suspending a magnetic needle provided with two pins fixed vertically at its ends.

When the needle is at rest, draw a line $AB$ joining the tips of the two pins. The magnetic needle is reversed upside down. Pins are fixed at the ends of the needle. When the magnet is at rest, draw a line $CD$ joining the tips of pins. $O$ is the point of intersection of $AB$ and $CD$. The line RS bisecting the angle $BOD$ is the magnetic meridian at that place (Fig. 3).
Now the angle between geographic meridian \( PQ \) and the magnetic meridian \( RS \) is the angle of declination \( \theta \) (Fig. 4).

2 Dip

Dip is defined as the angle between the direction of Earth’s magnetic field and the direction of horizontal component of earth’s magnetic field. It is the angle by which the total Earth’s magnetic field dips or comes out of the horizontal plane. It is denoted by \( \delta \). The value of dip varies from place to place. It is \( 0^\circ \) along the equator and \( 90^\circ \) at the poles. At Chennai the value of dip is about \( 9^\circ 7' \).

At a place the value of dip is measured by an instrument called dip circle.

Dip circle

A magnetic needle \( NS \) is pivoted at the centre of a circular vertical scale \( V \) by means of a horizontal rod. The needle is free to move over this circular scale. The scale has four segments and each segment is graduated from \( 0^\circ \) to \( 90^\circ \) such that it reads \( 0^\circ - 0^\circ \) along the horizontal and \( 90^\circ - 90^\circ \) along the vertical. The needle and the scale are enclosed in a rectangular box \( A \) with glass window. The box is mounted on a vertical pillar \( P \) on a horizontal base, which is provided with levelling screws. The base has a circular scale graduated from \( 0^\circ \) to \( 360^\circ \) (Fig. 5). The box can be rotated about a vertical axis and its position can be read on the circular scale with the help of a vernier (not shown in the figure).

The levelling screws are adjusted such that the base is horizontal and the scale inside the box is vertical. The box is rotated so that the ends of the magnetic needle \( NS \) read \( 90^\circ - 90^\circ \) on the vertical scale.

The needle, in this position is along the vertical component of the Earth’s field. The horizontal component of Earth’s field being perpendicular to the plane, does not affect the needle. This shows that the vertical scale and the needle are in a plane at right angles to the magnetic meridian. Now the box is rotated through an angle of \( 90^\circ \) with the help of the horizontal circular scale. The magnetic needle comes to rest exactly in the magnetic meridian. The reading of the magnetic needle gives the angle of dip at that place.