

Engineering Physics (BT-201)

B.Tech. I Year

Unit I - QUANTUM PHYSICS

OUTLINE

- WAVE PARTICLE DUALITY
- GROUP AND PARTICLE VELOCITIES
- UNCERTAINTY PRINCIPLE
- ELEMENTARY PROOF AND APPLICATIONS (Uncertainty Principle)
- COMPTON SCATTERING
- ENERGY AND MOMENTUM OPERATORS
- TIME DEPENDENT AND TIME INDEPENDENT SCHRÖDINGER
WAVE EQUATION
- ONE DIMENSIONAL SQUARE POTENTIAL WELL

QUANTUM MECHANICS

We know that classical mechanics can successfully explain the motion of astronomical bodies (such as stars, planets satellites etc.) means Newton's law of motion, as well as macroscopic bodies (such as motion under). Except this motion of charged particles in e.m. fields, elastic vibrations in solids, propagation of sound waves in glass etc. can also be explained successfully by classical mechanics, but some phenomenon like black body radiation, photo-electric effect, Compton effect, specific heat of solids at low temperature, stability of atoms, emission and absorption of light etc. could not be explained, which is explained by quantum mechanics.

Wave particle duality or dual nature of light

Light obeys the phenomena of interference, diffraction, polarization, photoelectric effect, Compton Effect etc. The phenomena of interference, diffraction and polarization can be explained by assuming that light is a form of wave. By the wave theory of light, it has been proved that light possesses a wave nature. However, some other phenomena like photoelectric effect, Compton Effect and discrete emission and radiation can be explained only with the help of the quantum theory of light. According to the quantum theory, light radiation travels in the form of energy bundles called quanta of energy $h\nu$, where ν is the frequency of radiation. Hence according to the quantum theory, light possesses a corpuscular (particle) nature. Therefore, sometimes light obeys the wave theory and sometimes the corpuscular theory, Hence, light has dual nature.

DE-BROGLIE HYPOTHESIS OR DE-BROGLIE MATTER WAVES

Louis De-Broglie suggested that the dual nature is not only of light, but each moving material particle has the dual nature. He assumed a wave should be with each moving particle, (all micro particles) which is called the matter waves. Although these waves can travel through vaccum like e.m. waves, but these are different from e.m. waves, because these waves associated with all types of charged and neutral particles.

PROOF OF DE-BROGLIE WAVES

According to Plank's quantum theory light is in form of small bundles of energy ($h\nu$) called quanta or photons.

If we consider a photon (quantum) to be a wave of frequency ν then its energy

$$E = h\nu \quad \dots\dots(i) \quad \text{Or} \quad E = hc/\lambda \quad [\text{i.e. } c = \nu \lambda]$$

Where c = velocity of light (or photon) in vaccum

λ = wavelength of photon or radiation

h = Plank's constant = 6.625×10^{-34} J-sec

Now if we consider photon as a particle of mass m , then from Einstein's mass-energy equivalence (or theory of relativity) energy of photon

$$E = mc^2 \dots\dots(ii) \quad \text{from eq.(i) \& (ii) } mc^2 = h\nu = hc/\lambda$$

$$\text{or } mc = h/\lambda$$

$$\text{or } \lambda = h/mc \quad [\text{i.e. } p = mc]$$

$$\text{or } \lambda = h/p \dots\dots(a) \quad \text{where } \lambda \text{ is a De-Broglie wavelength.}$$

Thus this is an evidence of De-Broglie nature of radiation (photon) the wavelength (λ) to particle like nature of radiation (photon) the momentum (p). It means this equation shows dual nature of light.

We know that the kinetic energy of the particle

$$E = \frac{1}{2} (mv^2) = \frac{1}{2} (m^2v^2/m) = P^2/2m$$

$$\text{or } P^2 = 2mE$$

$$\text{or } P = \sqrt{2mE}$$

So, De-Broglie wavelength $\lambda = h/\sqrt{2mE} \dots \dots \dots (b)$

According to kinetic theory of gases the average kinetic energy of the material particle

$$\frac{1}{2} (mv^2) = \frac{3}{2} kT$$

$$\text{Or } m^2v^2 = 3mkT$$

$$\text{Or } P^2 = 3mkT$$

$$\text{Or } P = \sqrt{3mkT}$$

So De-Broglie wavelength $\lambda = h/\sqrt{3mkT} \dots \dots \dots (c)$

where $k = \text{Boltzmann's constant} = 1.38 \times 10^{-23} \text{ J/K}$

$T = \text{Absolute temperature}$

If an particle accelerated through a potential difference of V volts, then

Work done by electric field = increase (gain) in kinetic energy

$$\text{or } qV = \frac{1}{2} mv^2$$

$$\text{or } mqV = \frac{1}{2} m^2v^2$$

$$\text{or } 2mqV = m^2v^2 = P^2$$

$$\text{or } \sqrt{2mqV} = P$$

So, De-Broglie wavelength $\lambda = h/\sqrt{2mqV}$

In case of electron wavelength of wave associated with the moving electron

$$\lambda = h/\sqrt{2meV}$$

$$\text{or } \lambda = \sqrt{(150/V)} \text{ \AA} = 12.27/\sqrt{V} \text{ \AA}.$$

PROPERTIES OF DE-BROGLIE WAVES

1. We know that the wavelength of matter waves (de-Broglie waves) associated by a moving particle is $\lambda = h/mv$

It means $\lambda \propto (1/m)$ and $\lambda \propto (1/v)$.

For a particle at rest means $v = 0$ so $\lambda = \infty$ and if $v = \infty$ then $\lambda = 0$. Here $\lambda = 0$ means that the matter waves are generated only when the particles are in motion.

2. Matter waves are independent of charge because it generated by any moving particle.

3. Energy of particle is given by $E = hv$ or $v = E/h$, where v is the frequency of wave.

We know that $E = mc^2$, where c is the velocity of light.

So we can write $v = mc^2/h$(1)

We know that wavelength of a wave associated with the particle of mass m , moving with velocity v is, $\lambda = h / mv$ (2)

If de-Broglie's wave velocity (phase velocity) is v_p , then

$$v_p = v\lambda$$

$$\text{or } v_p = mc^2/h * h / mv$$

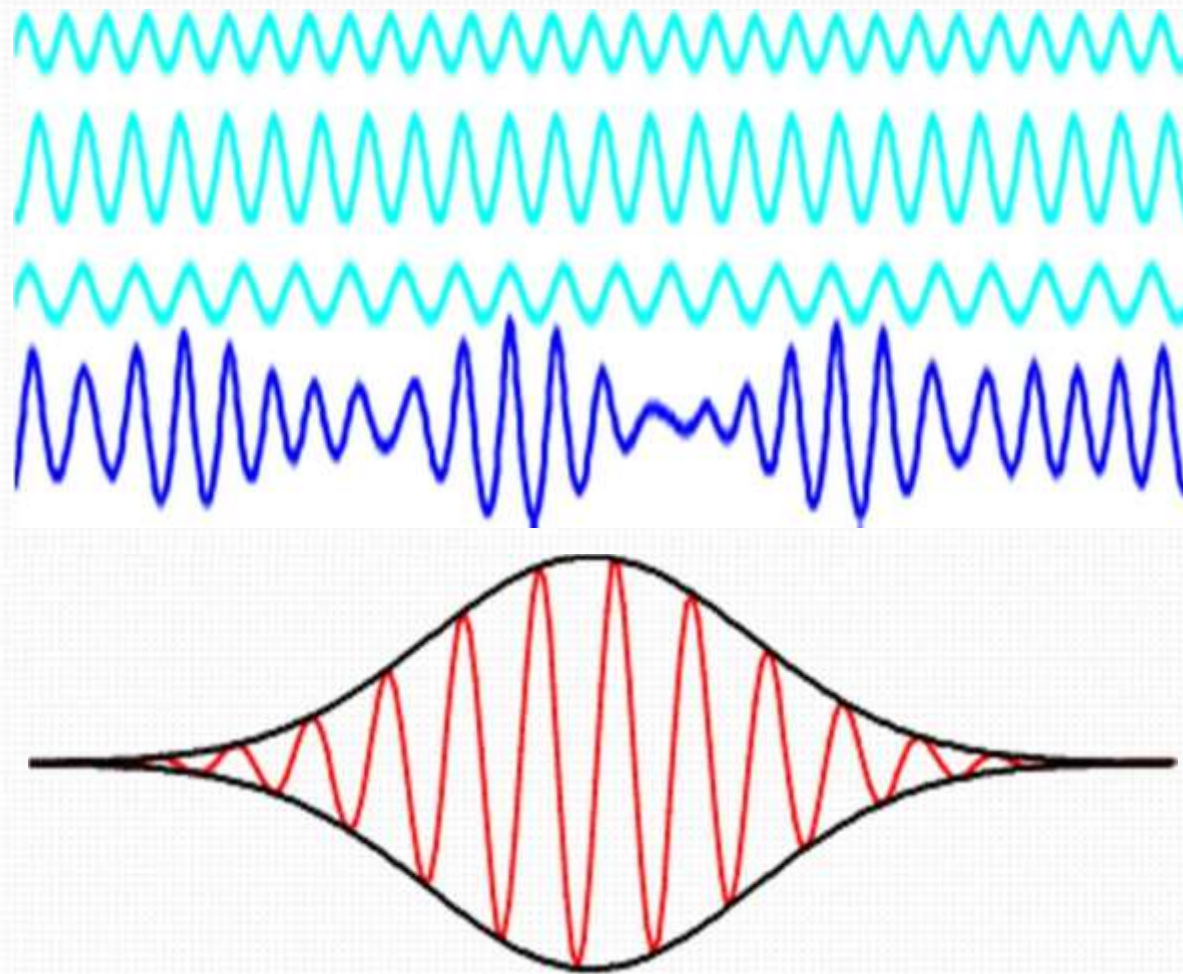
$$\text{or } v_p = c^2 / v$$

- From Einstein's theory of velocity, the speed of light is maximum speed that can be attained by a particle in nature.
- Thus from equation $v_p = c^2 / v$, the velocity of de-Broglie's wave associated with the particle would travel faster than the particle itself,
- Hence it is evident that a particle will not equivalent to a single wave, but equivalent to a group of wave, called wave packet or wave group.

Wave- Packet or Relation B/W group velocity and phase velocity

- A wave packet consist of a group of several waves of slightly different velocities & wavelength and formed by the superposition of waves situated on and around the center wavelength given by the de-Broglie formula.
- The amplitude and phase of the component waves are such that they interfere constructively in limited region where the particle is found and outside this region they interfere destructively, so amplitude falls to zero rapidly.
- Thus when several waves of slightly different wavelength travel along a straight line in one direction. The resultant waves obtained due to their superposition in form of group of waves which is called the wave packet.

Wave- Packet



The velocity of component or constituents waves of a wave packet is called phase velocity and the velocity of the wave packet is called the group velocity.

Consider the wave group arise from the combination of two waves that have the same amplitude A , but differ by an amount $O\omega$ in angular frequency and an amount Ok in wave number. These two waves can be represented as,

$$y_1 = A \cos (\omega t - kx)$$

$$y_2 = A \cos [(\omega + O\omega)t - (k + Ok)x]$$

Where $\omega = 2\pi\nu$ is angular frequency & $k = 2\pi/\lambda$ is wave number (propagation constant)

The motion of wave packet or the displacement equation of wave packet obtained due to their superposition will be at any point x and time t ,

CONT..

$$Y = y_1 + y_2$$

$$\text{or } Y = A [\cos(\omega t - kx) + \cos\{(\omega + O\omega)t - (k + Ok)x\}]$$

$$\text{or } Y = \frac{2A \cos(\omega t - kx) + \{(\omega + O\omega)t - (k + Ok)x\}}{2} \cdot \frac{\cos[(\omega t - kx) - \{(\omega + O\omega)t - \{(k + Ok)x\}]}{2}$$

$$[\text{Because } \cos C + \cos D = 2 \cos \{C+D\}/2 \cdot \cos\{C-D\}/2]$$

$$\text{Or, } Y = \frac{2A \cos[(\omega t - Kx) + (\omega t + O\omega t - Kx - OKx)]}{2} \times \frac{\cos[(\omega t + O\omega t - Kx - OKx - \omega t + Kx)]}{2}$$

$$\text{Or, } Y = \frac{2A \cos[2(\omega t - Kx) + (O\omega t - OKx)]}{2} \times \frac{\cos(O\omega t - OKx)}{2}$$

$$\text{Or, } Y = \frac{2A \cos[(2\omega + O\omega)t - (2K + OK)x]}{2} \times \frac{\cos(O\omega t - OKx)}{2}$$

Because $O\omega \ll \omega$ and $OK \ll K$, so we can write $2\omega + O\omega = 2\omega$ and $2K + OK = 2K$

$$\text{Or, } Y = 2A \cos[(2\omega t - 2Kx)/2] \times \cos(O\omega t - OKx)/2]$$

$$Y = 2A \cos(\omega t - Kx) \times \cos(O\omega t - OKx)/2$$

$$\text{Or, } Y = 2A \cos(O\omega t - OKx)/2 \times \cos(\omega t - Kx)$$

This equation represents a wave of angular frequency ω & propagation constant K .

The amplitude of this wave is $2A \cos(O\omega t - OKx)/2$

The velocity of group wave (v_g) is the velocity with which the maximum amplitude moves and maximum amplitude will be $2A$ i.e.,

$$2A \cos(O\omega t - OKx)/2 = 2A$$

It means $\cos(O\omega t - OKx)/2 = 0$ means $O\omega t - OKx = 0$

$$\text{Or } O\omega t = OKx, \text{ or } (O\omega / OK)t = x$$

Because we know that $x/t = v$ (distance / time = velocity). So in case of wave packet group

$$\text{velocity } v_g = \frac{dx}{dt} = \frac{d}{dt} \frac{O\omega t}{OK} = O\omega / OK$$

$$\text{Or, we can write } v_g = \frac{d\omega}{dK} \text{-----(1)}$$

$$\text{Because when } = \lim_{OK \rightarrow 0} \left(\frac{O\omega/2}{OK/2} \right)$$

Because the displacement equation of a wave $y = A \cos (\omega t - Kx)$

So, phase or wave velocity $v_p = dx/dt$

So, A will be maximum when $\cos (\omega t - Kx) = 1$ Or, $\omega t - Kx = 0$ or $x = \omega t / K$

$$\text{So, phase velocity } v_p = \frac{dx}{dt} = \frac{d}{dt} \frac{\omega t}{K} \text{ or } v_p = \omega / K \text{-----(2)}$$

So, from equation (1) $v_g = d\omega/dK$

$$\text{Means } v_g = \frac{d}{dK} (Kv_p)$$

$$\therefore v_p = \omega/K$$

$$\text{or } v_g = v_p \times 1 + K \cdot dv_p/dK$$

$$\text{or } v_g = v_p + (2\pi/\lambda) \cdot dv_p/d(2\pi/\lambda)$$

$$\text{or } v_g = v_p - (2\pi/\lambda) \cdot \frac{dv_p}{(2\pi/\lambda^2)d\lambda}$$

$$\therefore d(1/x) = -(1/x)dx$$

$$\text{or } v_g = v_p - \frac{\lambda \cdot dv_p}{d\lambda}$$

Thus group velocity (v_g) depends on the phase velocity (v_p) and variation of phase velocity with the wavelength ($dv_p/d\lambda$).

DIFFERENT

CASES
In the non-dispersive medium: If $dv_p/d\lambda = 0$ i.e., the medium is such that in it the phase velocity doesn't depend on the wavelength then,

$$v_p = v_g$$

i.e., the group velocity and phase velocity are equal. Such a medium is which the $v_p = v_g$ is called non-dispersive medium. For example; $v_p = v_g$ for electromagnetic wave in vacuum and elastic wave in a homogenous medium.

In the dispersive medium: In a dispersion medium $dv_p/d\lambda$ is positive and $v_g < v_p$. For example; v_g of electromagnetic wave in a dielectric substance is less than v_p .

But if the value of $dv_p/d\lambda$ is negative in a dispersive medium, then $v_p > v_g$. For example; in electric conductors the group velocity v_g is more than the phase velocity v_p .

Relationship between group velocity & phase velocity for the De-Broglie wave associated with a particle of rest mass (m_0) moving with a velocity v).

➤ (Relativistic Case)

Because we know that angular frequency,

$$\omega = 2\pi\nu = 2\pi.(mc^2/h)$$

$$[\because E = mc^2 = h\nu \text{ or } \nu = mc^2/h]$$

$$\text{Or } \omega = \frac{2\pi n_0 c^2}{h \sqrt{1 - \frac{v^2}{c^2}}} \text{-----(1)}$$

$$\left[\because n = \frac{n_0}{\sqrt{1 - \frac{v^2}{c^2}}} \right]$$

And propagation constant $K = 2\pi/\lambda = 2\pi.(mv/h)$

$$[\because \lambda = h/p = h/mv]$$

$$K = \frac{2\pi n_0 v}{h \sqrt{1 - \frac{v^2}{c^2}}} \text{-----(2)}$$

So phase velocity $v_p = \omega/K = c^2/v$, and

Group velocity $v_g = \frac{d\omega}{dK} = \frac{d\omega/dv}{dK/dv}$

$$\text{Now } \frac{dm}{dv} = \frac{d}{dv} \left[\frac{2nn_0c^2}{h\sqrt{1-\frac{v^2}{c^2}}} \right] \quad \text{or} \quad \frac{d}{dv} \left[\frac{2nn_0c^2}{h} \left(1 - \frac{v^2}{c^2}\right)^{-\frac{1}{2}} \right]$$

$$\frac{dm}{dv} = \frac{2nn_0c^2}{h} \left[\frac{d}{dv} \left(1 - \frac{v^2}{c^2}\right)^{-\frac{1}{2}} \right] \quad \text{After Differentiation}$$

$$\frac{dm}{dv} = \frac{2nn_0c^2}{h} \left[\frac{v}{c^2} \left(1 - \frac{v^2}{c^2}\right)^{-\frac{3}{2}} \right] \quad \text{Similary for } \frac{dK}{dv} = \frac{2nn_0}{h} \left[\left(1 - \frac{v^2}{c^2}\right)^{-\frac{3}{2}} \right]$$

So we can write Group velocity $v_g = \frac{dm}{dK} = \frac{\frac{dm}{dv}}{\frac{dK}{dv}} = v$

$$v_p = \frac{c^2}{v_g}$$

So, the De – Broglie wave group associated with a moving body travels with the same velocity as the moving particle. It is evident that a material particle in motion is equivalent to a group of waves or a wave packet.

Relation between V_p and V_g for a non- relativistic free particle (non-relativistic case)

Suppose v_g & v_p represents the group and phase velocity respectively for a non-relativistic free particle of mass m .

Let λ is the De- Broglie wavelength and ν is the frequency of the wave then, the phase velocity, $v_p = \nu\lambda$ -----(1)

According to De-Broglie hypothesis, $\lambda = h/mv_g$, -----(2)

Total energy, $E = \text{kinetic energy} = mv_g^2/2$

[Because potential energy of freely moving particle is constant]

Also $E = h\nu$. Or, $\nu = E/h = mv_g^2/2h$ -----(3)

So from equation (1), (2) and (3),

Because phase velocity $v_p = \nu\lambda = mv_g^2/2h \times h/mv_g$

$$v_p = v_g / 2$$

Hence for a non- relativistic free particle the phase velocity is half of the group velocity.

Heisenberg's Uncertainty Principle

“It is impossible to determine simultaneously the position and momentum of the particle with any desired accuracy.”

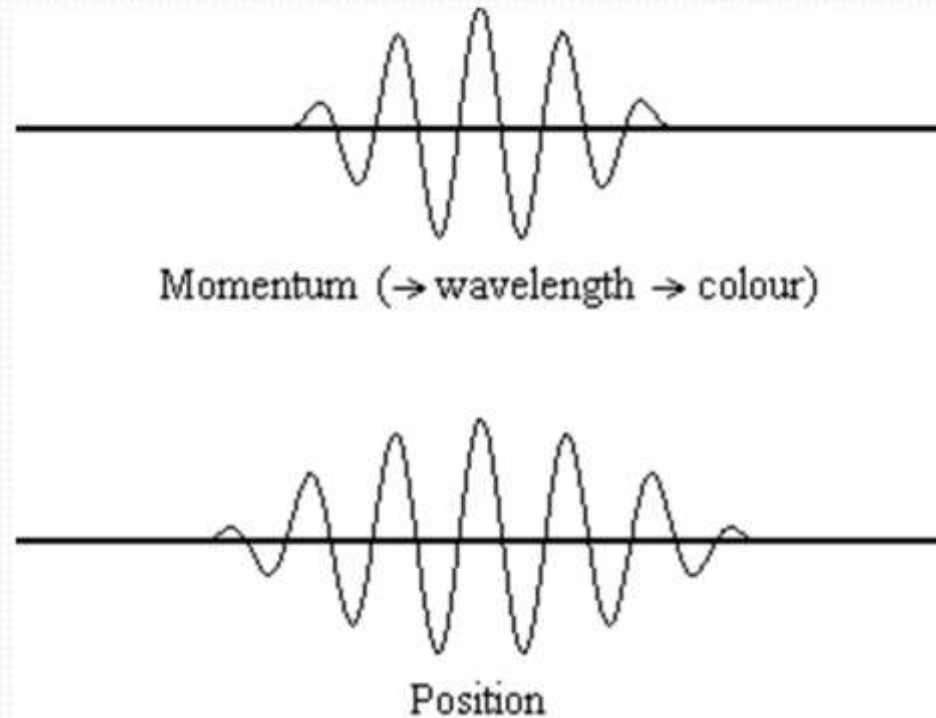
This definition is known as Heisenberg's uncertainty principle.

- This limitation is critical when dealing with small particles such as electrons.
- But it does not matter for ordinary-sized objects such as cars or airplanes.
- To locate an electron, you might strike it with a photon.
- The electron has such a small mass that striking it with a photon affects its motion in a way that cannot be predicted accurately.
- The very act of measuring the position of the electron changes its momentum, making its momentum uncertain.



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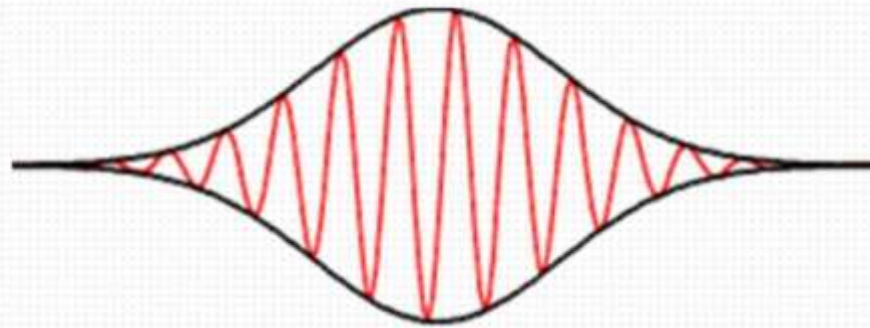
If we want accuracy in position, we must use short wavelength photons because the best resolution we can get is about the wavelength of the radiation used. Short wavelength radiation implies high frequency, high energy photons. When these collide with the electrons, they transfer more momentum to the target. If we use longer wavelength, i.e. less energetic photons, we compromise resolution and position.



Proof of Heisenberg's principle

1. Position and Momentum uncertainty:-

Heisenberg's principle can be proved by assuming that a particle in motion is equal to wave group and $v_g = v$. We know that a moving particle is equal to a wave group rather than a particle. It means there is limit for the measurement of particle properties.



Let a particle surrounded by a wave group. Let this wave group arises from the combination of two waves that have same amplitude A but differ by an amount $\Delta\omega$ in angular frequency and an amount Δk in propagation constant.

These two waves can be represented as :-

$$y_1 = A \cos(\omega t - kx)$$

$$\text{and } y_2 = A \cos[(\omega + \Delta\omega)t - (k + \Delta k)x]$$

The displacement of wave group at any time & at any point x will be.

$$Y = y_1 + y_2$$

$$\text{or } Y = A [\cos(\omega t - kx) + \cos\{(\omega + O\omega)t - (k + Ok)x\}]$$

$$\text{or } Y = 2A \cos \left[\frac{(\omega t - kx) + \{(\omega + O\omega)t - (k + Ok)x\}}{2} \right] \cdot \cos \left[\frac{(\omega t - kx) - \{(\omega + O\omega)t - (k + Ok)x\}}{2} \right]$$

$$[\text{Because } \cos C + \cos D = 2 \cos \{C+D\}/2 \cdot \cos\{C-D\}/2]$$

$$\text{Or, } Y = 2A \cos \left[\frac{(2\omega + O\omega)t - (2K + OK)x}{2} \right] \times \cos(O\omega t - OKx)/2$$

Because $O\omega \ll \omega$ and $OK \ll K$, so we can write $2\omega + O\omega = 2\omega$ and $2K + OK = 2K$

$$\text{Or, } Y = 2A \cos(O\omega t - OKx)/2 \times \cos(\omega t - Kx)$$

This equation represents a wave (wave packet) of angular frequency ω & propagation constant k . The amplitude of this wave is $2A \cos(\omega t - kx)/2$

Because the particle is moving with a velocity equal to the velocity of wave packet so the position of the particle can be anywhere in the wave packet so we can say the uncertainty in the position of the particle equals to the length of wave packet. Means distance between two consecutive modes. Node is formed when amplitude = 0, means

$$2A \cos(\omega t - kx)/2 = 0$$

$$\text{Or } \cos(\omega t - kx)/2 = 0$$

It is possible when $(\omega t - kx)/2 = \pi/2, 3\pi/2, 5\pi/2, \dots$

At a particular time t at positions x_1 and x_2 ,

$$(O\omega t - OK x_1)/2 = \pi/2 \text{ ----- (1)}$$

$$(O\omega t - OK x_2)/2 = 3\pi/2 \text{ ----- (2)}$$

From equations (1) and (2), we can write

$$OK/2.(x_2 - x_1) = \pi$$

If the error in the measurement of the position of the particle is $(x_2 - x_1) = O_x$, then

$$OK.O_x = 2\pi \text{ ----- (3)}$$

If O_p is error in the measurement of the momentum of the particle and OK is propagation constant then,

$$OK = 2\pi/O\lambda = (2\pi.O_p)/h$$

So from Equation (3) we can

write,

$$(2\pi.O_p)/h.O_x = 2\pi$$

$$\text{Or } O_p.O_x = h$$

$$\text{Or } O_p.O_x \geq \hbar/2\pi$$

It means that if we make simultaneous measurement of two quantities namely position and momentum of the particle the product of the fundamental errors is approximately equal to Planck's constant.

Application of Uncertainty principle

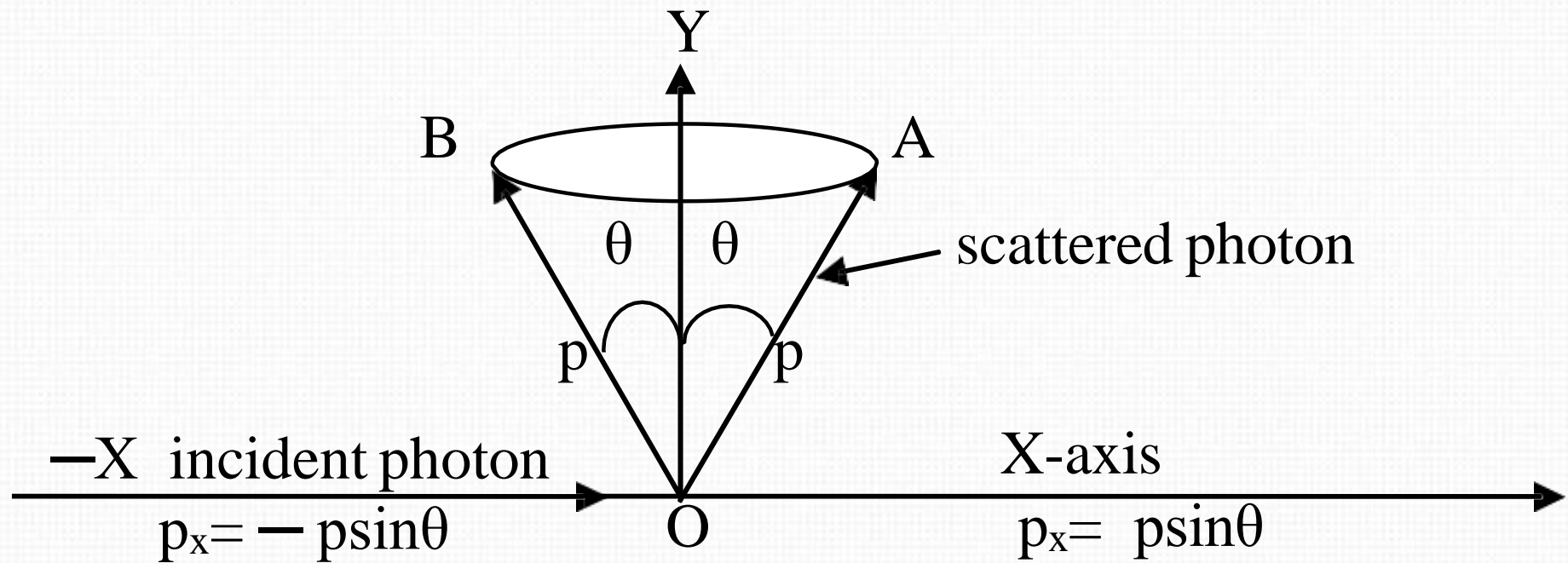
Determination of position of α -particle by γ -ray microscope:

For the determination of position of an particle (electron), it should be targeted (illuminated) by light (photons). From physical optics it is known that the exactness of position increase with decrease in wavelength of light. So uncertainty in position, Δx , so we take γ -rays microscope to be small, wavelength λ should be small. We know that from the theory of resolving power, the minimum distance between two objects x_1 and x_2 to be seen clearly or the uncertainty in the position is

$$(\Delta x) = \Delta x = \lambda / 2 \sin \theta \text{ -----(1)}$$

Where λ is wavelength of γ -rays and θ is semi-vertical angle of the cone of the objective with the object.

Suppose γ -rays incident on a stationary electron at O. The scattered γ -rays to be seen by the microscope which reaches through the objective of the microscope AB. The γ -rays consist of photons of energy $h\nu$ and momentum $h\nu/c$. Due to collision of photons with electron there will be a transfer of momentum from photon to electron. Therefore change in momentum of photon which will reach the objective of microscope.



So the momentum of photon in the direction of motion of electron will be $p \sin\theta$ and in the direction opposite of motion of electron will be $p \sin(-\theta)$.

$$\begin{aligned} \text{Total change in momentum } \Delta p &= p \sin\theta - p \sin(-\theta) \\ \Delta p &= 2p \sin\theta \end{aligned} \quad (2)$$

From equation (1) (2), we get

$$\begin{aligned} \Delta p \cdot \Delta x &= \lambda/2 \sin\theta \cdot 2p \sin\theta \\ \text{Or } \Delta p \cdot \Delta x &= h/p \sin\theta \cdot p \sin\theta \\ \Delta p \cdot \Delta x &= h \\ \text{Or } \Delta p \cdot \Delta x &\geq \hbar/2\pi \end{aligned}$$

Thus in the determination of position of a particle by γ -ray microscope the product of uncertainty in position and uncertainty in momentum is of the order of h . It is in accordance with Heisenberg's Uncertainty principle.

2. Diffraction of an electron beam by a single slit: Consider an electron beam travelling in X-direction which is incident on a narrow slit AB. Since the electron beam has a wave behavior, so we get diffraction pattern on the screen. Because it is quite uncertain to say that from which place of the slit the electron passes. If the width of the slit is O_y , the maximum uncertainty in the position of electron screen = O_y (in Y-direction)

Obviously narrow the slit less the uncertainty in position. But according to the theory of diffraction at a single slit, for half angular width of principal maxima

$$Oy.\sin\theta = \lambda \text{ for the wave of wavelength } \lambda.$$

$$\text{Or } \sin\theta = \lambda/Oy$$

So maximum uncertainty in the position of electron is,

$$Oy = \lambda/\sin\theta \dots\dots\dots(1)$$

If the wavelength of wave associated with electron is λ , then from de-Broglie $\lambda = h/p$.

The momentum of electron in parallel to the slit (i.e. in Y- direction) can have any value between $p\sin\theta$ and $p\sin(-\theta)$ because the diffracted electron can be found anywhere within the principle maxima (angular spread from $-\theta$ to $+\theta$).

Therefore uncertainty in momentum in direction parallel to the slit

$$\Delta p_y = p\sin\theta - p\sin(-\theta) = p\sin\theta + p\sin\theta = 2p\sin\theta$$

$$\Delta p_y = 2(h/p).\sin\theta \dots\dots\dots(2)$$

From (1) and (2), we get

$$\Delta p_y. Oy = \lambda/\sin\theta.2(h/p).\sin\theta = 2h \geq \hbar$$

3. Non-existence of electron in nucleus: From Rutherford's experiment we know that size of nucleus is equal to 10^{-14} meter. If electron exists in the nucleus then uncertainty in position of electron is Δx is the same as the size of nucleus. It means

$$\Delta x = 10^{-14} \text{ meter.}$$

So minimum uncertainty in momentum $\Delta p_{\min} = h / \Delta x = 6.625 \times 10^{-34} \text{ J} \cdot \text{sec} / 10^{-14} \text{ meter}$

$$\text{Or } \Delta p_{\min} = 6.625 \times 10^{-34} \text{ kg} \cdot \text{m}^2 / \text{sec}^2 \cdot \text{sec} = 6.625 \times 10^{-20} \text{ kg} \cdot \text{m} / \text{sec} \quad 10^{-14} \text{ meter}$$

For electron of minimum momentum inside the nucleus minimum energy $OE_{\min} = (p^2 c^2 + m_0 c^4)^{1/2}$

$$OE_{\min} = [(3.31 \times 10^{-20})^2 (3 \times 10^8)^2 + (9.1 \times 10^{-31})^2 (3 \times 10^8)^4]^{1/2}$$

$$OE_{\min} = [(10.95 \times 10^{-40}) (9 \times 10^{16}) + (82.81 \times 10^{-61})$$

$$(9 \times 10^{32})]^{1/2} \quad OE_{\min} = [(98.55 \times 10^{-24}) + (414.05 \times 10^{-29})]^{1/2}$$

Because second term is much smaller than first in above equation, so we can neglect it.

$$OE_{\min} = 9.93 \times 10^{-12}$$

$$\text{joule Or } OE_{\min} =$$

$$9.93 \times 10^{-12} \text{ eV} \quad 1.6 \times 10^{-19}$$

$$\text{Or } OE_{\min} = 9.93 \times 10^{-12} \text{ MeV}$$

$$1.6 \times 10^{-19} \times 10^6$$

$$\text{Or } OE_{\min} = 52.26 \text{ MeV}$$

Because maximum kinetic energy of a β^- particle emitted from radioactive nuclei is of the order of 4 MeV.

Conditions For Acceptable Wave Function

Because $\psi^*\psi = |\psi|^2$ is a real quantity. Thus it is clear that $|\psi|^2$ is a real quantity and is a measure of probability density. Hence probability of finding the particle in a small volume $dv = |\psi|^2 dv = |\psi|^2 dx dy dz$. Since total probability of finding the particle in any position is unity.

$$\text{So } \int_{-\infty}^{+\infty} |\psi|^2 dv = 1 \text{ or } \int_{-\infty}^{+\infty} \psi^* \psi dv = 1$$

This condition is known as normalization condition. The function satisfied the condition is called normalized wave function.

So ψ must be normalized, single valued because at any instant t there can be only one probability for the particle to be at a point; ψ must be finite and continuous.

Expectation Value: -

$$\text{The expectation value of a quantity } \langle f(r) \rangle = \frac{\iiint \psi^* f(r) \psi dv}{\iiint \psi^* \psi dv}$$

-

$$\langle f(r) \rangle = \frac{\iiint f(r) |\psi|^2 dv}{\iiint |\psi|^2 dv}$$

If the wave function is normalized then $\iiint |\psi|^2 dv = 1$, so

$$\langle f(r) \rangle = \iiint f(r) |\psi|^2 dv$$

Schrodinger's Wave Equation

It is a differential equation of the de-Broglie waves associated with the particle and describes the motion of particle.

If a wave function associated with a particle which is moving with velocity v in +ve direction, then displacement of wave is given by

$$\begin{aligned}
 \psi &= Ae^{-i\omega(t - x/v)} \\
 &= Ae^{-i(\omega t - x \cdot 2\pi v / v\lambda)} \quad \text{.....} \\
 &= Ae^{-i(\omega t - x \cdot 2\pi / \lambda)} \quad \text{(Because } v = v\lambda \text{ and } \omega = 2\pi\nu) \\
 &= Ae^{-i(\omega t - k \cdot x)} \quad \text{(Because } k = 2\pi / \lambda) \\
 &= Ae^{-i 2\pi (vt - x/\lambda)} \\
 &= Ae^{-i 2\pi (Et/h - xP/h)} \quad \text{(Because } E = h\nu \text{ and } \lambda = h/p) \\
 &= Ae^{-i 2\pi/h (Et - P \cdot x)} \\
 &= Ae^{-i / \hbar (Et - P \cdot x)}
 \end{aligned}$$

This is a wave equation for a freely moving particle. Now differentiating equation (2) with respect to get

$$\begin{aligned}
 \frac{\partial \psi}{\partial t} &= - (iE / \hbar) \cdot Ae^{-i / \hbar (Et - P \cdot x)} \\
 \frac{\partial \psi}{\partial t} &= - (iE / \hbar) \cdot \psi \quad \text{(Because } \psi = Ae^{-i / \hbar (Et - P \cdot x)}) \\
 \frac{\partial \psi}{\partial t} &= - (iE / \hbar) \cdot \psi \\
 \text{or } E\psi &= (i\hbar) \frac{\partial \psi}{\partial t} \quad \text{.....(3)}
 \end{aligned}$$

So energy operator $E = (i\hbar) \cdot \partial/\partial t$ (4)

Now partially differentiating equation (2) with respect to x, we get

$$\partial\psi/\partial x = (iP/\hbar) \cdot Ae^{-i/\hbar (Et - P \cdot x)}$$

$$\partial\psi/\partial x = (iP/\hbar) \cdot \Psi \quad (\text{Because } \psi = Ae^{-i/\hbar (Et - P \cdot x)})$$

$$) \text{ or } P\psi = (\hbar/i) \cdot \partial\psi/\partial x$$

$$\text{or } P\psi = (-i\hbar) \cdot \partial\psi/\partial x \text{(5)}$$

So momentum operator $P = (-i\hbar) \cdot \partial/\partial x$

.....(6) We know that total energy of the particle

$$E = K.E. + P.E.$$

$$\text{` } E = mv^2/2 +$$

$$V \text{ or } E = p^2/2m$$

$$+ V$$

Total energy in terms of wave function or operating wave function on above equation, $E\psi = (p^2/2m) \psi + V\psi$ (7)

Putting the value of E & P from equation (4) & (6) we

$$\text{get, } [(i\hbar) \cdot \partial/\partial t] \psi = [(-i\hbar) \cdot \partial/\partial x]^2 \psi / 2m + V\psi$$

$$(i\hbar) \cdot \partial\psi/\partial t = [(\hbar^2/-1) \cdot \partial^2/\partial x^2] \psi / 2m +$$

$$V\psi \quad (i\hbar) \cdot \partial\psi/\partial t = (\hbar^2/-2m) \cdot \partial^2\psi/\partial x^2 + V\psi$$

$$(i\hbar) \cdot \partial\psi/\partial t = -(\hbar^2/2m) \cdot \partial^2\psi/\partial x^2 + V\psi$$

This is Schrodinger time dependent equation for one dimension. For three dimensional

$$\text{case, } (i\hbar) \cdot \partial\psi/\partial t = -(\hbar^2/2m) \cdot [\partial^2/\partial x^2 + \partial^2/\partial y^2 + \partial^2/\partial z^2] \psi + V\psi$$

$$\text{or } (i\hbar) \cdot \partial\psi/\partial t = -(\hbar^2/2m) \cdot \psi + V\psi \quad (\text{Where } = \partial^2/\partial x^2 + \partial^2/\partial y^2 + \partial^2/\partial z^2)$$

For free moving particle $V = 0$ so we can write, $(i\hbar) \cdot \partial\psi/\partial t = - (\hbar^2/2m) \cdot \psi$

Now putting the value of P from equation (4) to equation (7) we get,

$$E\psi = -(\hbar^2/2m) \cdot \partial^2\psi/\partial x^2 + V\psi$$

$$\text{or } (\hbar^2/2m) \cdot \partial^2\psi/\partial x^2 + (E - V)\psi = 0$$

$$\text{or } \partial^2\psi/\partial x^2 + (2m/\hbar^2) \cdot (E - V)\psi = 0 \dots\dots\dots(8)$$

This is Schrodinger time independent equation for one dimension. For three dimensional case, $[\partial^2/\partial x^2 + \partial^2/\partial y^2 + \partial^2/\partial z^2]\psi + (2m/\hbar^2) \cdot (E - V)\psi = 0$

$$\psi + (2m/\hbar^2) \cdot (E - V)\psi = 0$$

For free moving particle $V = 0$ so we can write,

$$\psi + (2m/\hbar^2) \cdot E\psi = 0$$

Q. Show that the function $\psi = A x e^{(-x^2/2)}$ is the eigen function of the operator A



UNIT II

WAVE OPTICS

Introduction:

Optics is a branch of physics which deals with the “Theory of light and its propagation in a given medium”. The branch of optics is divided into two parts

1. Ray or Geometrical optics and
2. Physical or Wave optics

Ray optics: Ray optics deals with image formation by optical systems. It was supported by Newton's corpuscular theory. It deals with the particles propagating in a medium called “light- photons with energy $E_n = nh\nu$ in a medium”

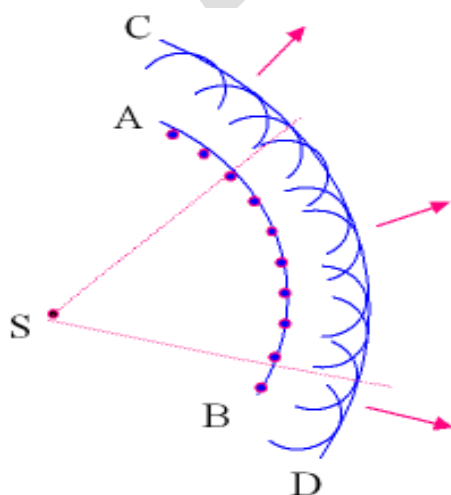
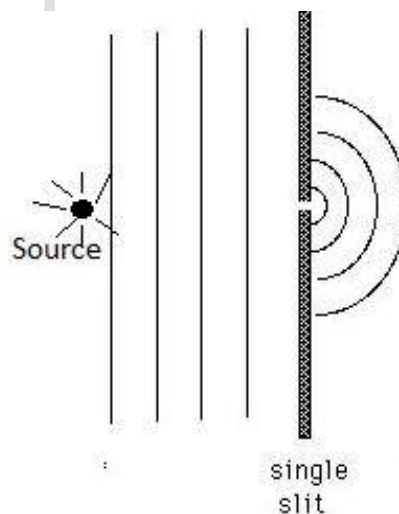
Physical optics: Physical optics deals with nature of light. Huygens proposed the wave theory of light. According to this, a luminous body is a source of disturbance in a hypothetical medium called “ether”

HUYGEN'S PRINCIPLE

Huygens' principle provides a geometrical method of finding the shape and position of the wave front at a certain instant from its shape and position of some earlier instant.

Huygens' principle is stated in the following two parts.

- Each point on the wave front acts as a centre of new disturbance and emits its own set of spherical waves called secondary wavelets. These secondary wavelets travel in all directions with the velocity of light so long as they move in the same medium.
- The envelope or the locus of these wavelets in the forward direction gives the position of the new wavefront at any subsequent times.



Huygens' Principle:

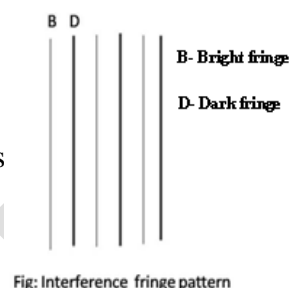
Each wavefront is the envelope of the wavelets. Each point on a wavefront acts as an independent source to generate wavelets for the next wavefront. AB and CD are two wavefronts.

Superposition Principle: According to superposition principle “The resultant or total displacement of the medium acted upon by two or more waves simultaneously equal to sum or difference of displacements of individual waves”.

$$R \text{ or } Y = y_1 \pm y_2 \pm \dots \pm y_n$$

Interference of light:

It is defined as “Modification of Resultant Intensity of light obtained by the superposition of two or more light waves”. This theory of interference of light was developed by Thomas Young in his experimental study of light. The resultant intensity consists of series of bright & dark fringes that appear on the screen which are known as interference pattern or interference fringes which correspond to maximum & minimum intensities of light.



Interference pattern may be either straight or circular or parabolic (Arc) fringes

Coherent Sources: If the phase difference between two light waves emitted from two sources is zero or has a constant value then the sources are said to be coherent.

Incoherent sources: If the phase difference between two light waves coming from two sources changes with time, the sources are called as “In-coherent Sources”.

Types of interference: Interference of light based on Young's double slit experiment is divided into two types.

1. Constructive interference
2. Destructive interference

Constructive interference: When the two light waves reach a point in phase the resultant displacement (R_{Max} or Y_{Max}) is always equal to algebraic sum of individual displacements of the light waves. It is known as constructive interference.

$$R_{\text{Max}} \text{ or } Y_{\text{Max}} = y_1 + y_2 + \dots + y_n$$

Destructive interference: When the two light waves reach a point i.e. out of phase. The resultant displacement (R_{Min} or Y_{Min}) is always equal to the difference of displacements of the light waves. It is known as destructive interference.

$$R_{\text{Min}} \text{ or } Y_{\text{Min}} = y_1 - y_2 - \dots - y_n$$

Conditions for interference of light:

- 1) The two light sources emitting light waves should be coherent.
- 2) The separation between the two sources should be small.
- 3) The distance between the two sources and the screen should be large.

- 4) To view interference fringes, the back-ground (good-contrast) should be dark.
- 5) Amplitudes of two light waves are nearly equal.
- 6) The sources should be mono-chromatic.
- 7) The sources should be narrow i.e. they must be small.

Interference in thin films due to reflected light:

When light is incident on a (plane parallel) thin film, some portion gets reflected from the upper surface and the remaining portion is transmitted into the film. Again, some portion of the transmitted light is reflected back into the film by lower surface and emerges through the upper surface. These reflected light beams superimpose with each other, producing interference and forming interference patterns.

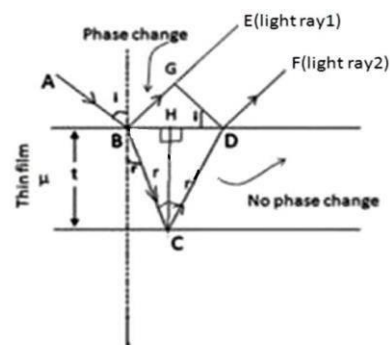


Fig-1: Interference in thin films due to reflection

Consider a thin film of thickness 't' and refractive index μ . Let a monochromatic light ray AB be incident at an angle of incidence of 'i' on the upper surface of the film. BE and BC are the reflected and transmitted light rays. Let the angle of refraction be 'r'. The ray BC will be reflected into the film and emerge through the film in the form of the light ray DF.

These two light rays superimpose depending upon path difference between them producing interference patterns. To know the path difference, let us draw the normal DG to BE. From D and G onwards, the light rays travel equal distances. By the time the light ray travels from B to G, the transmitted light ray has to travel from B to C and C to D.

The path difference between light rays (1) and (2) is

$$\Delta = \mu(BC + CD) \text{ (in film)} - BG \text{ (in air)} \text{ ----- (1)}$$

$$\text{From triangle BCH, } BC = \frac{HC}{\cos r} = \frac{t}{\cos r}$$

$$\therefore BC = \frac{HC}{\cos r} = \frac{t}{\cos r}$$

$$\text{From triangle DCH, } \cos r = \frac{HC}{CD}$$

$$\Rightarrow \cos r = \frac{t}{CD}$$

$$CD = \frac{t}{\cos r}$$

$$BC = CD = \frac{t}{\cos r}$$

$$\Rightarrow (BC+CD) = \frac{2t}{\cos r} \quad \text{----- (2)}$$

To calculate BG, $BD = BH + HD$

$$\text{Triangle BHC, } \tan r = \frac{BH}{CH} = \frac{BH}{t}$$

$$\therefore BH = t \tan r$$

$$\text{Similarly } HD = t \tan r$$

$$BD = BH + HD = 2t \tan r \quad (\because BH = HD)$$

$$\text{From triangle BGD, } \sin i = \frac{BG}{BD}$$

$$\Rightarrow BG = BD \sin i$$

$$\therefore BG = 2t \tan r \sin i$$

From Snell's law,

$$\mu = \frac{\sin i}{\sin r}$$

$$\sin i = \mu \sin r$$

$$BG = 2\mu t \tan r \sin r$$

Substituting the above values in equation-(1)

$$\text{path difference} = \mu \left(\frac{2t}{\cos r} \right) - 2\mu t (\tan r) \sin r$$

$$= \frac{2\mu t}{\cos r} (1 - \sin^2 r)$$

$$= \frac{2\mu t}{\cos r} (\cos^2 r)$$

$$\boxed{\text{path difference} = 2\mu t \cos r}$$

In the above case, the ray BE is reflected at the surface of a denser medium and hence it suffers an additional path difference equal to $-\frac{\lambda}{2}$.

The ray BC suffers reflection at C, at the surface of a rarer medium and hence on emerging out at D, it suffers no additional path change.

Hence the net path difference between the two reflected rays DF and BE

$$\text{Total path difference} = 2\mu t \cos r + \frac{\lambda}{2}$$

When the path difference is an integral multiple of λ , then the rays (1) and (2) will be in phase and appear as constructive interference.

$$\text{i.e., } 2\mu t \cos r + \frac{\lambda}{2} = n\lambda \quad (\text{Condition for the bright fringe})$$

$$2\mu t \cos r = n\lambda - \frac{\lambda}{2}$$

$$2\mu t \cos r = (2n-1) \frac{\lambda}{2} \quad \text{where } n=0, 1, 2, \dots$$

When the path difference is half integral multiple of λ , the rays (1) and (2) meet in out of phase and undergo destructive interference.

$$\text{i.e., } 2\mu t \cos r + \frac{\lambda}{2} = (2n-1) \frac{\lambda}{2} \quad (\text{Condition for the dark fringe})$$

$$2\mu t \cos r = (2n-1) \frac{\lambda}{2} - \frac{\lambda}{2}$$

$$2\mu t \cos r = n\lambda \quad \text{where } n=0, 1, 2, \dots$$

Depending on the condition the interference pattern consists of bright and dark fringes.

Newton's Rings:

Newton's rings are one of the best examples for the interference in a non uniform thin film. When a plano convex lens with its convex surface is placed on a plane glass plate, an air film of increasing thickness is formed between the two. The thickness of the film at the point of contact is zero.

If monochromatic light is allowed to fall normally and the film is viewed in the reflected light, alternate dark and bright rings concentric around the point of contact between the lens and glass plate are seen. These circular rings were discovered by Newton and are called as Newton's rings.

The plano-convex lens (L) of larger radius of curvature is placed with its convex surface on a plane glass plate (P). The lens makes the contact with the plate at 'O'. The monochromatic light falls on a glass plate G held at an angle of 45° with the vertical. The glass plate G reflects normally a part of the incident light towards the air film enclosed by the lens L and the glass plate P. A part of the light is reflected by the

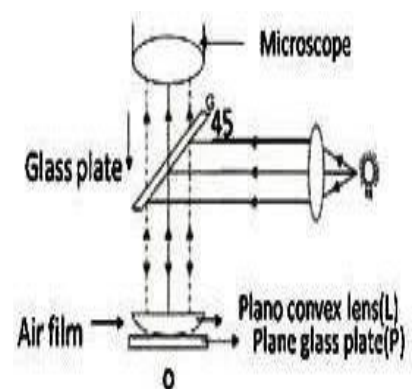


Fig 2 Newton's rings experiment

curved surface of the lens L and a part is transmitted which is reflected back from the plane surface of the plate. These reflected rays interfere and give rise to an interference pattern in the form of circular rings. These rings are seen near the upper surface of the air film through the microscope.

Explanation:

Newton's rings are formed due to interference between the light rays reflected from the top and bottom surfaces of an air film between the plate and the lens. A part of the incident light is reflected at a point in the form of the ray (1) with any additional phase change. The other part is refracted and again reflected in the form of the ray (2) with additional phase change of π or path change of $\lambda/2$.

As the rings are observed in the reflected light, the path difference between the rays is

$$\text{path difference} = 2\mu t \cos r + \frac{\lambda}{2}$$

For an air film $\mu = 1$ and for normal incidence $r = 0$,

then,
$$\text{path difference} = 2t + \frac{\lambda}{2}$$

At the point of contact $t = 0$ the path difference is $\lambda/2$, i.e., the reflected light at the point of contact suffers phase change of π . Then the incident and reflected rays are out of phase and interfere destructively. Hence the central spot is dark.

The condition for the bright ring is $2t + \frac{\lambda}{2} = n\lambda$

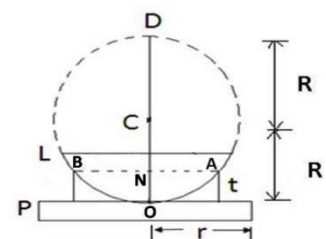
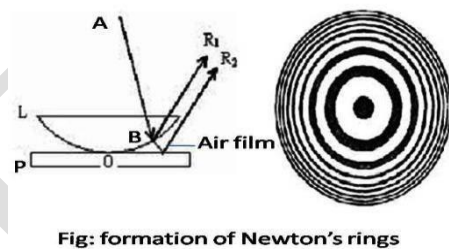
$$2t = (2n - 1) \frac{\lambda}{2} \text{ where } n = 1, 2, 3, \dots$$

The condition for dark ring is $2t + \frac{\lambda}{2} = (2n + 1) \frac{\lambda}{2}$

$$2t = n\lambda \text{ where } n = 0, 1, 2, 3, \dots$$

Theory:

To find the diameters of dark and bright rings, let 'L' be the lens placed on a glass plate P. The convex surface of the lens is the part of spherical surface with centre at 'C'. Let R be the radius of



curvature and R be the radius of Newton's ring corresponding to the film thickness ' t '.

From the property of a circle,

$$NA \times NB = NO \times ND$$

$$r \times r = t \times (2R - t)$$

$$r^2 = 2Rt - t^2$$

where R is the radius of curvature of plano lens and ' t ' is the maximum thickness of air film

As ' t ' is small, t^2 will be negligible, $r^2 = 2Rt$

$$t = \frac{r^2}{2R}$$

Condition for bright ring is, $2t = (2n-1) \frac{\lambda}{2}$

$$2 \frac{r^2}{2R} = (2n-1) \frac{\lambda}{2}$$

$$r^2 = \frac{(2n-1)\lambda R}{2}$$

Replacing r by $D/2$ the diameter of n^{th} bright ring will be

$$\frac{D^2}{4R} = \frac{(2n-1)\lambda}{2}$$

$$\therefore D \propto \sqrt{\text{odd natural number}} \quad (\text{for bright ring})$$

$$D \propto \sqrt{\text{natural number}} \quad (\text{for dark ring})$$

Thus, the diameters of dark rings are proportional to the square root of natural numbers and diameters of bright rings are proportional to odd natural numbers.

Determination of wavelength of light source:

Let R be the radius of curvature of a Plano convex lens, λ be the wavelength of light used. Let D_m and D_n be the diameters of m^{th} and n^{th} dark rings respectively.

$$D_m^2 = 4m\lambda R$$

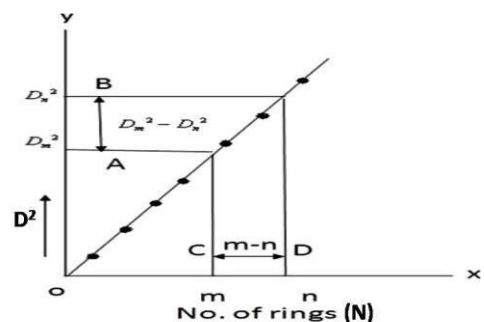


Fig: Plot of D^2 w.r.t. no. of rings

$$D_n^2 = 4n\lambda R$$

$$D_m^2 - D_n^2 = 4(m-n)\lambda R$$

$$\lambda = \frac{D_m^2 - D_n^2}{4(m-n)R}$$

$$R = \frac{D_m^2 - D_n^2}{4(m-n)\lambda}$$

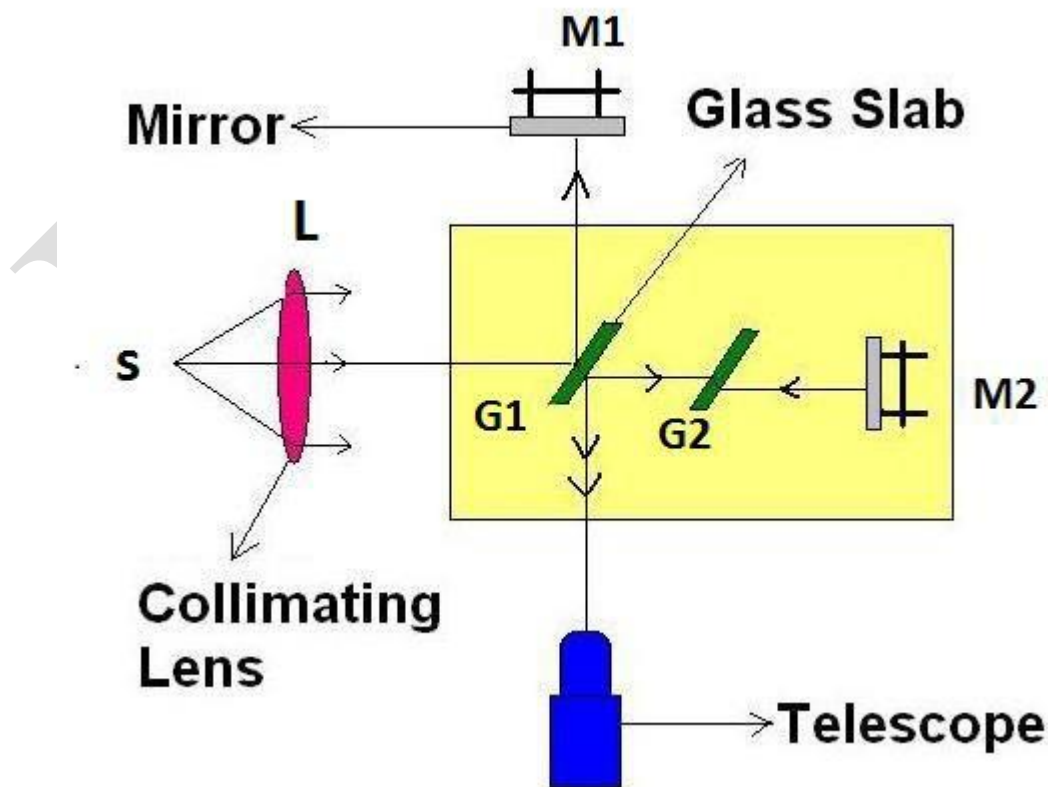
From the graph, $\frac{AB}{CD} = \frac{D_m^2 - D_n^2}{(m-n)}$

The radius of the planoconvex lens can be obtained with the help of a spherometer. λ can be calculated by substituting all the above values.

MICHELSON INTERFEROMETER

Construction:

- It consists of two excellent optically plane, highly polished plane mirrors M_1 and M_2 which are at right angle to each other. There are two optically flat glass plates G_1 and G_2 of same thickness and of the same material placed parallel to each other. These plates are inclined at an angle 45° with the mirrors M_1 and M_2 . T is a telescope which receives the reflected light from mirrors M_1 and M_2 .



Working:

- Monochromatic light from source S after being rendered parallel by a collimating lens L fallson the semi silvered glass plate G_1 .
- It is divided into two parts, the reflected ray travels towards mirror M_1 and the transmitted ray travels towards mirror M_2 .
- The two rays fall normally on mirrors M_1 and M_2 respectively and are reflected back along their original paths.
- The reflected ray again meets at the semisilvered surface of G_1 and enters a short focus telescope T.
- The two rays entering the telescope are originally derived from the same single beam, hence they are in a position to produce interference fringes in the field of view of the telescope.
- A desired path difference can be introduced between the two reflected rays by moving the mirror M_1 .
- It can be noted from figure that the reflected ray passes through G_1 twice where as the transmitted ray does not do so even once.
- Thus in the absence of glass plate G_2 the two paths are not equal.
- To equalise the two paths, glass plate G_2 of same thickness and material as that of G_1 is introduced in the path of transmitted ray.
- Because of this nature, the glass plate G_2 is called compensating plate.
- The interference fringes may be straight, circular, parabolic etc depending upon path difference and angle between mirrors.

Uses of Michelson's sinterferometer:

Michelson's sinterferometer has been used for a variety of purposes for example

1. In the determination of wavelength of monochromatic source of light.
2. To determine the difference between the two neighbouring wavelengths or resolution of the spectral lines.
3. In the determination of refractive index and thickness of various thin transparent materials.

DIFFRACTION:

Introduction:

The bending of light around the edges of an obstacle is called diffraction. It was first observed by Gremaldy.

When light falls on an obstacle then the corresponding geometrical shadow on the screen should be completely dark. Practically the geometrical shadow consists of bright and dark fringes. These fringes are due to the superposition of bent light waves around the corners of

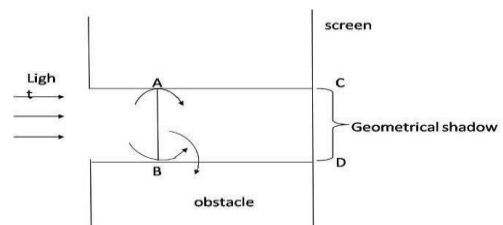


Fig: Diffraction

an obstacle. The amount of bending depends upon the size of an obstacle and wavelength of light.

When light falls on an obstacle whose size is comparable with the wavelength of light then light bends around the edges or corners of an obstacle and enters into the geometrical shadow. This bending of light is known as **diffraction**.

Types of Diffraction:

The diffraction phenomena are broadly classified into two types.

1. Fresnel's Diffraction:

In this type of diffraction, the source of light and the screen are placed at finite distance. In this, lenses are not necessary to study the diffraction. This diffraction can be studied in the direction of propagation of light. The incident wave fronts are either spherical or cylindrical.

2. Fraunhofer's Diffraction:

In this type of diffraction, the source and screen are placed at infinite distances. Here we need lenses to study the diffraction. This diffraction can be studied in any direction. In this, the incident wavefront is plane.

Difference between Fresnel diffraction and Fraunhofer diffraction:

S.No.	Fresnel diffraction	Fraunhofer Diffraction
1	Source and screen are placed at finite distances.	Source and screen are placed at infinite distances.
2	No lenses are used.	Lenses are used.
3.	The incident wavefront is spherical or cylindrical.	The incident wavefront is plane wavefront.
4.	The diffraction can be studied in the direction of propagation of light.	The diffraction can be studied in any direction of propagation of light.

Fraunhofer's diffraction at single slit:

Consider a slit AB of width ' e '. Let a plane wavefront WW' of monochromatic light of wavelength λ propagating normally towards the slit is incident on it. The diffracted light

through the slit be focused by means of a convex lens on a screen placed in the focal plane of the lens.

According to Huygens-Fresnel, every point on the wavefront in the plane of the slit is a source of secondary wavelets, which spread out to the right in all directions. These wavelets travelling normal to the slit i.e., along the direction OP_0 are brought to focus at P_0 by the lens. Thus, P_0 is a bright central image.

The secondary wavelets travelling at an angle θ with the normal are focused at a point P_1 on the screen. Depending on the path difference, the point P_1 may have maximum or minimum intensities. In order to find out intensity at P_1 , let us draw a parallel AC from A to the light ray at B .

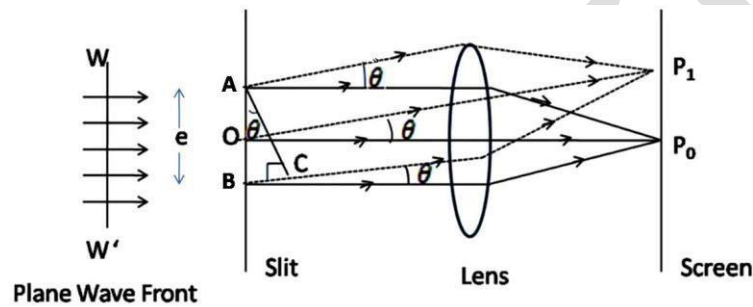


Fig: Fraunhofer Diffraction - Single Slit

The path difference between secondary wavelets from A and B in direction θ , is given by

$$\Delta = BC = AB \sin \theta = e \sin \theta$$

The relationship between phase difference and path difference is given by

$$\text{Phase difference} = \frac{2\pi}{\lambda} \times \text{path difference}$$

$$= \frac{2\pi}{\lambda} \times e \sin \theta$$

Let the width of the slit be divided into 'n' equal parts and the amplitude of the wave from each part is 'a'. Then the phase difference between any two successive waves is

$$\frac{1}{n} (\text{Total phase}) = \frac{1}{n} \left(\frac{2\pi}{\lambda} e \sin \theta \right) = d \quad (\text{say})$$

Using the method of vector addition of amplitudes the resultant amplitude R is given by

$$R = \frac{a \sin \left(\frac{nd}{2} \right)}{\sin \left(\frac{d}{2} \right)} = \frac{a \sin \left(\frac{\pi e \sin \theta}{\lambda} \right)}{\sin \left(\frac{\pi e \sin \theta}{n\lambda} \right)}$$

$$= \frac{a \sin \alpha \sin \left(\frac{\alpha}{n} \right)}{\alpha} \quad \text{Where } \alpha = \frac{\pi e \sin \theta}{\lambda}$$

$$= a \frac{\sin \alpha}{\alpha/n} \left(\because \frac{\alpha}{n} \text{ is very small} \right)$$

$$= \frac{n a \sin \alpha}{\alpha}$$

Let $na = A$,

$$R = \frac{A \sin \alpha}{\alpha}$$

Resultant intensity (I) is proportional to square of amplitude (R). Therefore, the intensity is given by

$$I = R^2 = A^2 \left(\frac{\sin \alpha}{\alpha} \right)^2$$

Intensity distribution:

The diffraction pattern consists of central principal maximum for $\alpha = 0$. There are subsidiary or secondary maxima of decreasing intensity on either side of it at positions $\alpha = \pm \frac{3\pi}{2}, \pm \frac{5\pi}{2}$. Between secondary maxima, there are minima at positions $\alpha = \pm \pi, \pm 2\pi$.

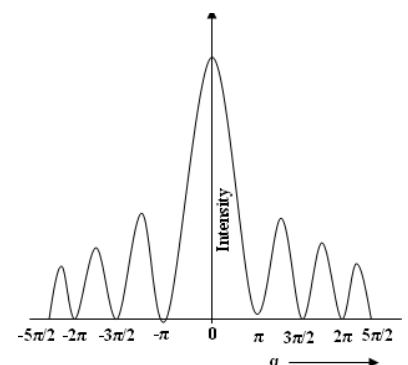


Fig: Intensity Distribution

Diffraction grating:

Diffraction grating is a closely placed multiple slits. It consists of a very large number of narrow slits side by side separated by opaque spaces. The incident light is transmitted through the slits and blocked by opaque spaces. Such a grating is called a transmission grating.

When light passes through the grating, each one of the slits diffracts the waves. All the diffracted waves combine one another producing sharp and intense maxima on the screen.

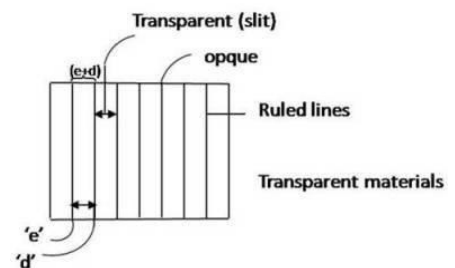


Fig: Diffraction grating

A plane transmission grating is a plane sheet of transparent material on which opaque lines are made with a diamond point. The space between the rulings is equal and transparent and consists of parallel slits. The combined width of a ruling and a slit is called grating element.

Let 'e' be the width of the line and 'd' be the width of the slit. Then (e + d) is known as grating element. If 'N' is the no. of lines per inch on the grating then

$$N(e+d)=1''=2.54\text{cm}$$

$$e+d=\frac{2.54}{N}\text{cm}$$

When light falls on the grating, the light gets diffracted through each slit. As a result, both diffraction and interference of diffracted light gets enhanced and forms a diffraction pattern. This pattern is known as diffraction spectrum.

Gratings spectrum:

The condition to form the principal maxima in a grating is given by

$$(e+d)\sin\theta=n\lambda \text{ is called grating equation.}$$

where (e+d) is the grating element and $n=1, 2, \dots$

Instead of monochromatic source of light such as sodium vapour lamp, if white light source such as mercury is used then each diffracted order will have different colours at different angles.

For $n=1$

$$(e+d)\sin\theta_v=\lambda_v \text{ (for violet ray)}$$

$$(e+d)\sin\theta_r=\lambda_r \text{ (for red ray)}$$

Thus in the spectrum for a grating, there is no overlapping or mixing of colours unlike the spectrum for a prism where different colours overlap. For a grating the angle of diffraction (i.e., angular dispersion) depends on λ and (e + d). Hence if two different gratings of same (e+d) values are chosen, they will produce same dispersion and hence they will be identical.

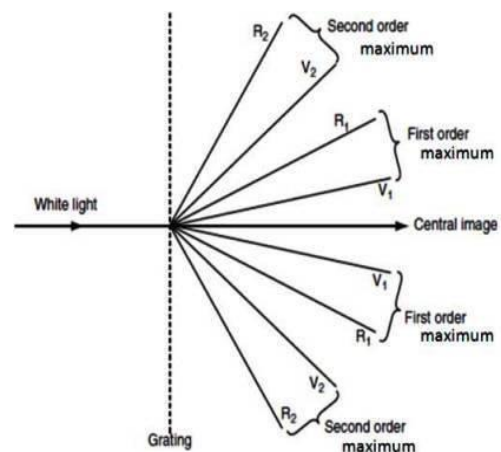


Fig: Grating spectrum for monochromatic source

Maximum no. of orders available with a grating:

The principal maxima in a grating

$$(e+d)\sin\theta=n\lambda$$

$$\frac{1}{e+d} = N$$

$$\text{Hence } r < r_2, \mu > \mu_e$$

For $\theta = 90^\circ$, the maximum possible value of $\sin \theta$ is 1.

$$nN\lambda \leq 1 \quad \text{or } n \leq \frac{1}{N\lambda}$$

This gives the maximum number of orders possible and n is an integer.

Resolving power of a grating:

The resolving power of a grating represents its ability to form separate spectral lines for wavelengths very close together. It is measured by $\frac{\lambda}{d\lambda}$, where $d\lambda$ is the smallest wavelength difference that can just be resolved at wavelength λ .

Analysis:

Let a parallel beam of light of two wavelengths λ and $\lambda + d\lambda$ be incident normally on the grating. If the n th principal maximum of λ is formed in the direction θ_n , we will have $(e + d)\sin \theta_n = n\lambda$ where $(e + d)$ is the grating element.

Now the grating equation for the minima is $N(e + d)\sin \theta = m\lambda$ where N is the total number of rulings on the grating and m can take all integral values except $0, N, 2N, \dots, Nn$, because these values of m give respective principal maxima. It is clear from figure that the first minimum of $\lambda + d\lambda$ adjacent to n th principal maximum of $(\lambda + d\lambda)$ in the direction of increasing θ will be obtained for $m = Nn + 1$. Therefore, for this minimum we have

$$(e + d)\sin(\theta_n + d\theta_n) = (nN + 1)\lambda$$

$$(e + d)\sin(\theta_n + d\theta_n) = \left(\frac{nN + 1}{N}\right)\lambda$$

According to Rayleigh's criterion, the wavelengths λ and $\lambda + d\lambda$ are just resolved by the grating when the n th maximum of $\lambda + d\lambda$ is also obtained in direction $\theta_n + d\theta_n$, i.e.

$$(e + d)\sin(\theta_n + d\theta_n) = n(\lambda + d\lambda) \dots \dots \dots (4)$$

Figure shows the overlapping of principal maxima of two patterns

Comparing eqs (3) and (4), we get

$$\left(\frac{nN + 1}{N}\right)\lambda = n(\lambda + d\lambda)$$

$$nN\lambda + \lambda = nN\lambda + nNd\lambda$$

$$\lambda / d\lambda = nN$$

But $\left(\frac{\lambda}{d\lambda}\right)$ is the resolving power R of the grating. Therefore,

$$R = nN$$

$$R = \frac{(e + d)\sin \theta_n}{\lambda}$$

As expected, the resolving power is zero for the central principal maximum ($n = 0$), all wavelengths being undiffracted in this order.

Unit III

INTRODUCTION TO SOLIDS

Free electron theory:

In solids, electrons in outer most orbits of atoms determine its electrical properties. Electron theory is applicable to all solids, both metals and non metals. In addition, it explains the electrical, thermal and magnetic properties of solids. The structure and properties of solids are explained employing their electronic structure by the electron theory of solids.

It has been developed in three main stages:

1. Classical free electron theory
 2. Quantum Free Electron Theory.
 3. Zone Theory.
- **Classical free electron theory:** The first theory was developed by Drude & Lorentz in 1900. According to this theory, metal contains free electrons which are responsible for the electrical conductivity and metals obey the laws of classical mechanics.
 - **Quantum Free Electron Theory:** In 1928 Sommerfield developed the quantum free electron theory. According to Sommerfield, the free electrons move with a constant potential. This theory obeys quantum laws.
 - **Zone Theory:** Bloch introduced the band theory in 1928. According to this theory, free electrons move in a periodic potential provided by the lattice. This theory is also called “Band Theory of Solids”. It gives complete informational study of electrons.

Classical free electron theory:

Even though the classical free electron theory is the first theory developed to explain the electrical conduction of metals, it has many practical applications. The advantages and disadvantages of the classical free electron theory are as follows:

Advantages:

1. It explains the electrical conductivity and thermal conductivity of metals.
2. It verifies Ohm's law.
3. It is used to explain the optical properties of metals.
4. Metal composed of atoms in which electrons revolve around the nucleus are many states available for occupation. If the density of states is zero, no states can be occupied at that energy level.
5. The valence electrons are freely moving about the whole volume of the metals like the molecules of perfect gas in a container
6. The free electrons move in random directions and collide with either positive ions or other free electrons. Collision is independent of charges and is elastic in nature
7. The movement of free electrons obeys the laws of classical kinetic theory of gases
8. Potential field remains constant throughout the lattice.
9. In metals, there are large numbers of free electrons moving freely within the metal i.e. the free electrons or valence electrons are free to move in the metal like gaseous

molecules, because nuclei occupy only 15% metal space and the remaining 85% space is available for the electrons to move.

Drawbacks:

1. It fails to explain the electric specific heat and the specific heat capacity of metals.
2. It fails to explain superconducting properties of metals.
3. It fails to explain new phenomena like photoelectric effect, Compton effect, black-body radiation, etc.
4. It fails to explain Electrical conductivity (perfectly) of semiconductors or insulators.
5. The classical free electron model predicts the incorrect temperature dependence of σ . According to the classical free electron theory, $\sigma \propto T^{-1}$.
6. It fails to give a correct mathematical expression for thermal conductivity.
7. Ferromagnetism couldn't be explained by this theory.
8. Susceptibility has greater theoretical value than the experimental value.

Quantum free electron theory of metals:

Advantages:

1. All the electrons are not present in the ground state at 0K, but the distribution obeys Pauli's exclusion principle. At 0K, the highest energy level filled is called Fermi-level.
2. The potential remains constant throughout the lattice.
3. Collision of electrons with positive ion cores or other free electrons is independent of charges and is elastic in nature.
4. Energy levels are discrete.
5. It was successful to explain not only conductivity, but also thermionic emission, paramagnetism, specific heat.

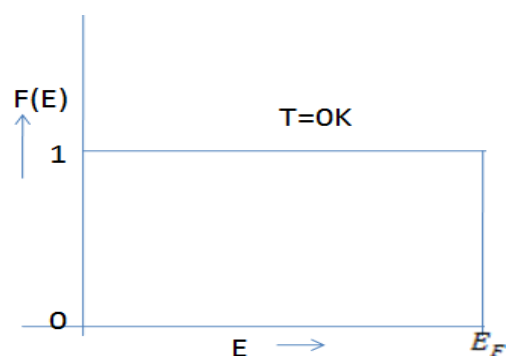
Drawbacks:

1. It fails to explain classification of solids as conductors, semiconductors and insulators.

Fermi level and Fermi energy:

The distribution of energy states in a semiconductor is explained by Fermi-Dirac statistics since it deals with the particles having half integral spin like electrons. Consider that the assembly of electrons as electron gas which behaves like a system of Fermi particles or fermions. The Fermions obeying Fermi-Dirac statistics i.e., Pauli's exclusion principle.

Fermi energy: It is the energy of state at which the probability of electron occupation is $\frac{1}{2}$ at any temperature above 0K. It separates filled energy states and unfilled energy states. The highest energy level that can be occupied by an electron at 0 K is called Fermi energy level.



Fermi level: It is a level at which the electron probability is $\frac{1}{2}$ at any temp above 0K (or) always it is 1 or 0 at 0K.

Therefore, the probability function $F(E)$ of an electron occupying an energy level E is given by,

$$F(E) = \frac{1}{1 + \exp\left(\frac{E - E_F}{KT}\right)} \dots\dots\dots (1)$$

Where E_F known as Fermi energy and it is constant for a system,

K is the Boltzmann constant and T is the absolute temperature.

Case I : Probability of occupation at $T = 0K$, and

$E < E_F$ Therefore $F(E) = 1$, as per above, clearly indicates that at $T = 0K$, the energy level below the Fermi energy level E_F is fully occupied by electrons leaving the upper level vacant.

Therefore, there is 100% probability that the electron occupies energy level below Fermi level.

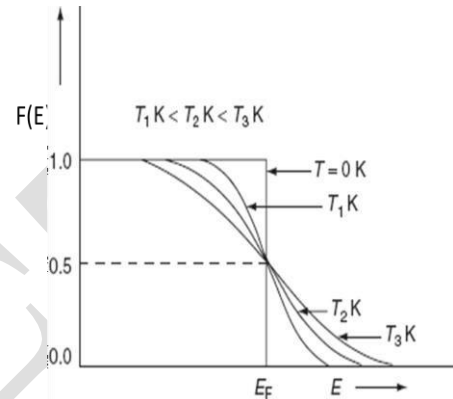
Case II: Probability of occupation at $T = 0K$, and $E > E_F$

Then

i.e., all levels below E_F are completely filled and all levels above E_F are completely empty. As the temperature rises $F(E)$.

Case III: Probability of occupation at $T = 0K$, and $E = E_F$

$$F(E) = \frac{1}{1 + e^0} = \frac{1}{1 + 1} = \frac{1}{2} = 0.5$$



The above condition states that, $T = 0K$, there is a 50% probability for the electron to occupy Fermi energy.

The probability function $F(E)$ lies between 0 and 1.

Hence there are three possible probabilities namely

$F(E) = 1$ 100% probability to occupy the energy level by electrons.

$F(E) = 0$ No probability to occupy the energy levels by electrons and hence, it is empty. $F(E)$

$= 0.5$ 50% probability of finding the electron in the energy level.

Density of States (DOS):

The number of electrons per unit volume in an energy level at a given temperature is equal to the product of density of states (number of energy levels per unit volume) and Fermi Dirac distribution function (the probability to find an electron).

$$n_c = \int g(E) \times f(E) dE \dots\dots\dots (1)$$

where n_c is the concentration of electrons, $g(E)$ is the density of states & $F(E)$ is the occupancy probability.

The number of energy states with a particular energy value E is depending on how many combinations of quantum numbers resulting in the same value n .

To calculate the number of energy states with all possible energies, we construct a sphere in 3D- space with 'n' as radius and every point (n_x , n_y and n_z) in the sphere represents an energy state.

As every integer represents one energy state, unit volume of this space contains exactly one state. Hence; the number of states in any volume is equal to the volume expressed in units of cubes of lattice parameters). Also

$$n^2 = n_x^2 + n_y^2 + n_z^2$$

Consider a sphere of radius n and another sphere of radius n+dn with the energy values are E and (E+dn) respectively.

Therefore, the number of energy states available in the sphere of radius 'n' is

by considering one octant of the sphere

(Here, the number of states in a shell of thickness dn at a

distance 'n' in coordinate system formed by

n_x, n_y and n_z and will take only positive values, in that

sphere $\frac{1}{8}$ of the volume will satisfy this condition).

The number of energy states within a sphere of radius (n+dn) is

$$\frac{1}{8} \left(\frac{4\pi}{3} \right) (n+dn)^3$$

Thus the number of energy states having energy values between E and E+dn is given by

$$\begin{aligned} g(E)dn &= \left(\frac{1}{8} \frac{4\pi}{3} \right) (n+dn)^3 - \left(\frac{1}{8} \frac{4\pi}{3} \right) n^3 \\ &= \frac{1}{8} \left(\frac{4\pi}{3} \right) (n+dn)^3 - n^3 \\ &= \frac{4\pi}{6} n^2 dn = \frac{\pi}{2} n^2 dn \end{aligned}$$

compared to 'dn', dn^2 and dn^3 are very small.

Neglecting higher powers of dn

$$g(E)dn = \frac{\pi}{2} n^2 dn \dots \dots (2)$$

The expression for nth energy level can be written as ,

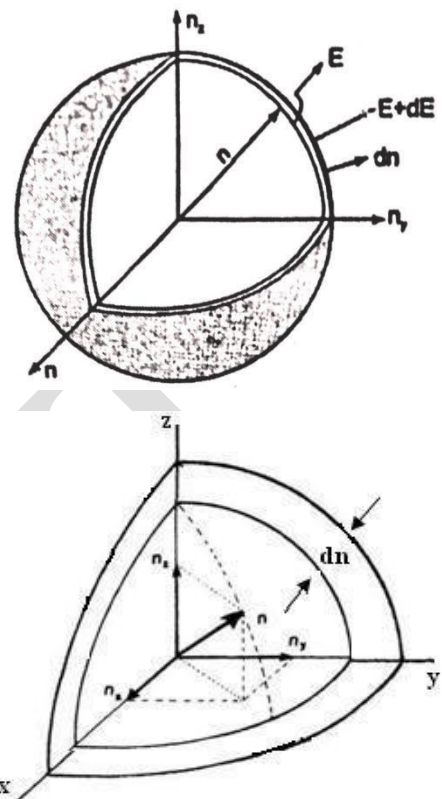
$$E = \frac{n^2 h^2}{8mL^2} \quad \text{or,} \quad n^2 = \frac{8mL^2 E}{h^2} \dots \dots (3)$$

$$\Rightarrow n = \left(\frac{8mL^2 E}{h^2} \right)^{1/2} \dots \dots (4)$$

Differentiating eq.(3):

$$2n dn = \frac{8mL^2}{h^2} dE \quad \Rightarrow dn = \frac{1}{2} \left(\frac{8mL^2}{h^2} \right)^{1/2} dE$$

∴ by substituting 1/n value in dn,



$$n = \frac{1}{2} \left(\frac{18mL^2}{h^2} \right)^{1/2} \frac{dE}{E^2} \dots\dots\dots(5)$$

Substituting n^2 and dn from eq. (3) and (5), we get

$$g(E)dE = \frac{\pi 8mL^2}{2} \left(\frac{1}{h^2} \right) dE \times 2 \left(\frac{18mL^2}{h^2} \right)^{1/2} \frac{dE}{E^{1/2}}$$

$$g(E)dE = \frac{\pi 8mL^2}{4} \left(\frac{1}{h^2} \right)^{3/2} E^{1/2} dE \dots\dots\dots(6)$$

According to Pauli's Exclusion Principle, two electrons of opposite spin can occupy each energy state

Equation (6) should be multiplied by 2

$$g(E)dE = 2 \times \frac{\pi 8mL^2}{4} \left(\frac{1}{h^2} \right)^{3/2} E^{1/2} dE$$

After mathematical simplification, we get

$$g(E)dE = \frac{4\pi}{h^3} (2m)^{3/2} E^{1/2} L dE$$

The density of energy states $g(E) dE$ per unit volume is given by,

$$g(E)dE = \frac{4\pi}{h^3} (2m)^{3/2} E^{1/2} dE \quad \because L=1$$

Bloch Theorem:

According to free electron model, a conduction electron in metal experiences constant potential. But in real crystal, there exists a periodic arrangement of positively charged ions through which the electrons move. As a consequence, the potential experienced by electrons is not constant but it varies with the periodicity of the lattice. In zone theory, as per Bloch, potential energy of electrons is considered as varying potential with respect to lattice 'a'.

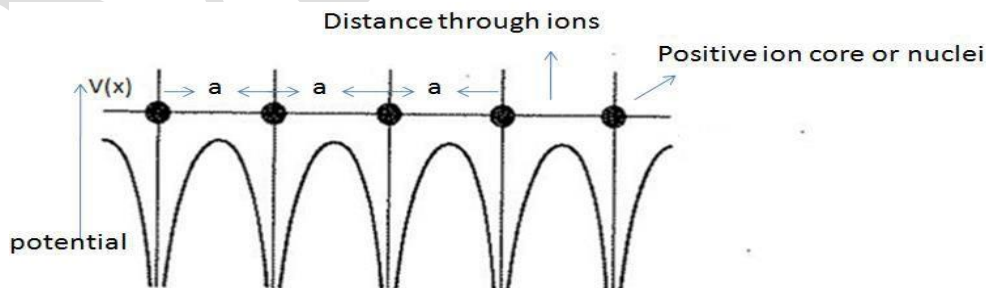


Fig: Variation of potential energy in a periodic lattice.

Let us examine one dimensional lattice as shown in figure. It consists of an array of ionic cores along the X-axis. A plot of potential V as a function of its position is shown in figure.

From graph:

At nuclei or positive ion cores, the potential energy of electron is minimum and in-between nuclei, the P.E. is considered as maximum w.r. to Lattice constant 'a'.

This periodic potential $V(x)$ changes with the help of lattice constant a , $V(x) = V(x+a)$ ('a' is the periodicity of the lattice)

To solve, by considering Schrodinger's time independent wave equation in one dimension,

$$\frac{d^2\psi}{dx^2} + \frac{8\pi^2m}{h^2} [E - V(x)]\psi = 0 \dots (1)$$

Bloch's 1D solution for Schrodinger wave equation (1) $\psi_k(x) = u_k(x) \exp(ikx) \dots (2)$

where $u_k(x) = u_k(x+a)$

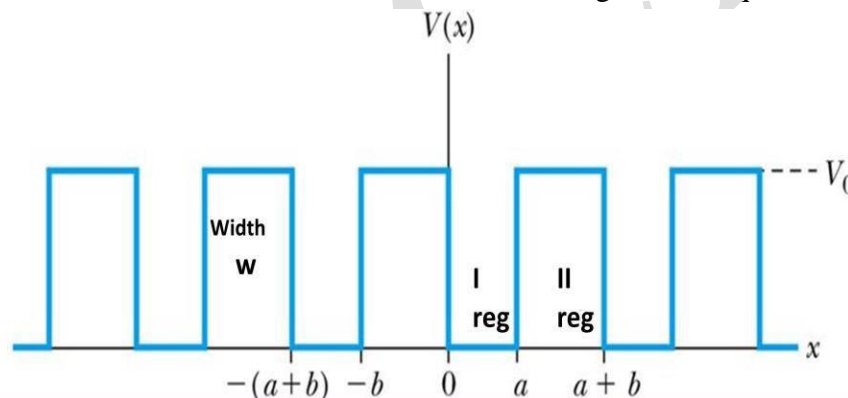
Here $u_k(x)$ - periodicity of crystal lattice, modulating function, k - propagation vector = $\frac{2\pi}{\lambda} e^{ikx}$ is plane wave.

By applying eq. (2) to eq. (1), it is not easy to solve Schrodinger wave equation and Bloch cannot explain complete physical information about an e^- in periodic potential field. Then Kronig Penny model was adopted to explain the electrical properties of an e^- .

Kronig-Penney model:

Kronig -penny approximated the potentials of an e^- s inside the crystal in terms of the shapes of rectangular steps as shown, i.e. square wells is known as Kronig Penny model.

i.e. The periodic potential is taken in the form of rectangular one dimensional array of square well potentials and it is the best suited to solve Schrodinger wave equation.



It is assumed that the potential energy is zero when x lies between 0 and a , and is considered as I region. Potential energy is V_0 , when x lies between $-b$ and 0. And considered as II region.

Boundary conditions:

$\psi(x) = 0$, where x lies between $0 < x < a$ - I region

$V(x) =$

V_0 , where x lies between $-b < x < 0$ - II region

This model explains many of the characteristic features of the behavior of electrons in a periodic lattice.

The wave function related to this model may be obtained by solving Schrodinger equations for the two regions,

$$\frac{d^2\psi}{dx^2} + \frac{2m}{\hbar^2} E \psi = 0, \text{ for } 0 < x < a \text{ with } V(x) = 0 \dots (1)$$

$$\frac{d^2\psi}{dx^2} + \frac{2m}{\hbar^2}(E - V_0)\psi = 0, \text{ for } -b < x < 0 \text{ with } V(x) = V_0 \dots \dots (2)$$

Again,

$$\frac{d^2\psi}{dx^2} + \alpha^2\psi = 0 \dots \dots (3) \quad \text{where } \alpha^2 = \frac{2mE}{\hbar^2} \text{ and } \alpha = \frac{2\pi}{h} \sqrt{2mE}$$

$$\frac{d^2\psi}{dx^2} - \beta^2\psi = 0 \dots \dots (4) \quad \text{where } \beta^2 = \frac{2m}{\hbar^2}(V_0 - E)$$

The solution of these equations from Bloch theorem, $\psi(x) = u_k(x) \exp(ikx)$. From figure, square well potentials, if V_0 increases, the width of barrier 'w' decreases, if V_0 decreases the width of barrier w increases. But the (product) barrier strength $V_0 w$ remains constant.

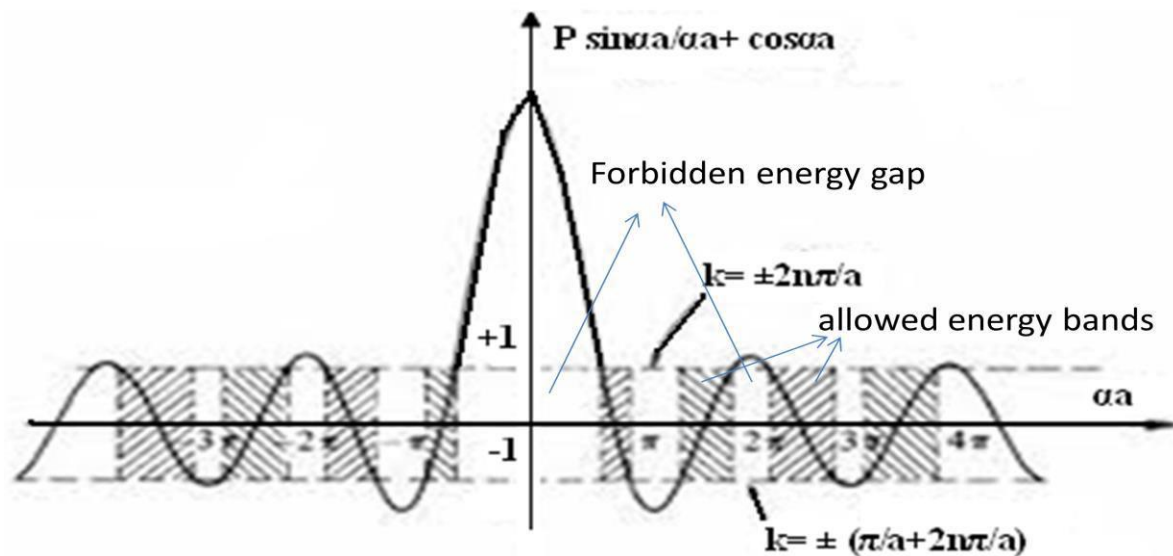
To get this, differentiating above Schrodinger wave equations 3 & 4 w.r. to x, and by applying boundary conditions of x (w.r. to their corresponding Ψ), to know the values of constants A, B of region-I, C, D for region-II, we get mathematical expression (by simplification)

$$\cos ka = P \frac{\sin \alpha a}{\alpha a} + \cos \alpha a$$

where,

$$P = \frac{4\pi^2 m a}{h^2} V_0 w \quad \text{and} \quad \alpha = \frac{2\pi}{h} \sqrt{2mE}$$

P-varying term, known as scattering power.



And 'v₀b' is known as barrier strength.

Conclusions:

1. The L.H.S. is a cosine term which varies between the limits -1 and +1, and hence the R.H.S. also varies between these limits. It means energy is restricted within -1 to +1 only.

2. If the energy of e^- lies between -1 to +1, are called **allowed energy bands** and it is shown by shaded portion in energy spectrum. This means that ' αa ' can take only certain range of values belonging to allowed energy band.
3. As the value of αa increases, the width of the allowed energy bands also increases.
4. If energy of e^- s not lies between -1 to +1 are known as **forbidden energy bands** and it decreases w.r.to increment of αa .
5. Thus, motion of e^- s. in a periodic lattice is characterized by the bands of allowed & forbidden energy levels.

Case1:

1. $P \rightarrow \infty$

If $P \rightarrow \infty$, the allowed band reduces to a single (line) energy level, gives us steeper lines.

We have

$$\cos ka \left(\frac{\alpha a}{p} \right) = \sin \alpha a + \cos \alpha a \left(\frac{\alpha a}{p} \right)$$

$$P \rightarrow \infty, \frac{1}{\infty} = 0 \quad \text{then} \quad \sin \alpha a = 0$$

$$\sin \alpha a = \sin n\pi$$

$$\alpha a = n\pi$$

$$\alpha^2 a^2 = n^2 \pi^2$$

$$\alpha^2 = \frac{n^2 \pi^2}{a^2}$$

$$\frac{2mE}{\hbar^2} = \frac{n^2 \pi^2}{a^2}$$

$$E = \frac{n^2 \pi^2 \hbar^2}{2ma^2} = \frac{n^2 \pi^2 \hbar^2}{2ma^2 4\pi^2}$$

$$E = \frac{n^2 \hbar^2}{8ma^2} \quad , \text{ here } a \text{ is lattice constant}$$

It means, it (zone theory) supports quantum free electron theory.

Case2:

$P \rightarrow 0$, We have

$$\cos ka = P \frac{\sin \alpha a}{\alpha a} + \cos \alpha a$$

$$\alpha a = ka$$

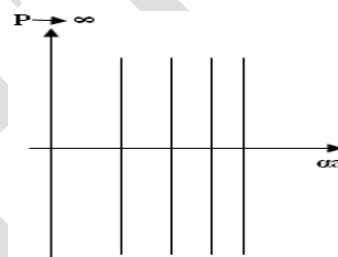
$$\alpha = k$$

$$\alpha^2 = k^2$$

$$\frac{2mE}{\hbar^2} = \left(\frac{2\pi}{\lambda} \right)^2 = \frac{4\pi^2}{\lambda^2}$$

$$E = \frac{4\pi^2 \hbar^2}{2m\lambda^2}$$

$$E = \frac{4\pi^2 \hbar^2}{2m\lambda^2 4\pi^2}$$



$$E = \frac{1}{2}mv^2$$

It gives us kinetic energy of an electron. It means zone theory supports classical free electron theory at this situation electron completely free electron not bounded with allowed and forbidden gaps (and no energy level exists).

Thus by varying P from 0 to ∞ , we find that the completely free electron(s) becomes completely bound to Brillouin Zone.

Brillouin Zone OR E-K diagram:

The Brillouin zone are the boundaries that are marked by the values of wave vector k , in which electrons can have allowed energy values. These represent the allowed values of k of the electrons in 1D, 2D, & 3D.

We have, the energy of the electron in a constant potential box is,

$$E = \frac{n^2 h^2}{8ma^2} \dots (1) \text{ where } a = \text{length of the box.}$$

But,

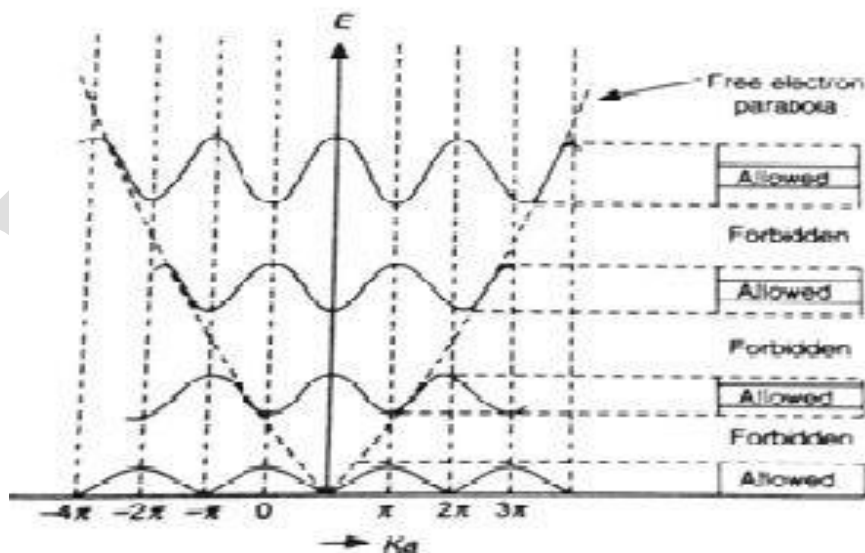
$$k = \frac{n\pi}{a} \Rightarrow k^2 = \frac{n^2 \pi^2}{a^2}$$

$$\frac{n^2}{a^2} = \frac{k^2}{\pi^2} \dots (2)$$

Substitute eqn (2) in (1), we get $E = \frac{k^2 h^2}{8m\pi^2} \propto k^2$. It represents parabolic equation.

A graph is drawn between the total energy (E) and the wave vector k , for various values of k .

i.e. $k = \frac{n\pi}{a} ; n = \pm 1, \pm 2, \pm 3, \dots$



It is the energy spectrum of an electron moving in the presence of a periodic potential field and is divided into allowed energy regions (allowed zones) or forbidden energy gaps (forbidden zones).

Allowed energy values lie in the region $k = -\pi/a$ to $+\pi/a$. This zone is called the first Brillouin zone. After a break in the energy values, called a forbidden energy band, we have

another allowed zone spread from $k = -\pi/a$ to $-\pi/a + 2\pi/a$ and $+\pi/a$ to $+\pi/a + 2\pi/a$. This zone is called the second Brillouin zone. Similarly, higher Brillouin zones are formed.

Concept of effective mass of electron:

When an electron in a periodic potential of lattice is accelerated by a known electric field or magnetic field, then the mass of the electron is called effective mass and is represented by m^* . To explain, let us consider an electron of charge 'e' and mass 'm' moving inside a crystal lattice of electric field E.

Then by taking known expression $F = ma$, can be considered here as $F = m^*a$ (1)

The acceleration $a = \frac{eE}{m}$ is not constant in the periodic lattice but varies due to the change in electronic mass.

If free electron under wave packet, the group velocity V_g corresponding to the particle's velocity can be written as

$$v_g = \frac{d\omega}{dk} = \frac{1}{\hbar} \frac{dE}{dk} \quad \text{where } E = \hbar\omega \quad \dots (2)$$

The rate of change of velocity is known as

$$\text{Acceleration, } a = \frac{dv_g}{dt} = \frac{1}{\hbar} \frac{d^2E}{dk^2} \frac{dk}{dt}$$

$$a = \frac{1}{\hbar} \frac{dE}{dk} \frac{dk}{dt} = \frac{1}{\hbar} \frac{d^2E}{dk^2} \frac{dk}{dt} \quad \dots (3)$$

From quantum mechanics relation, $p = \hbar k$ (4)

$$\text{and } F = \frac{dp}{dt} \quad \dots (5)$$

By differentiating eq (4) w.r.t. t, and by substituting eq (5)

$$\hbar \frac{dk}{dt} = \frac{dp}{dt} = F \Rightarrow \frac{dk}{dt} = \frac{F}{\hbar} \quad \dots (6)$$

by substituting eq. (6) in eq (3),

$$\therefore a = \frac{1}{\hbar} \frac{d^2E}{dk^2} \frac{F}{\hbar}$$

By rearranging the above term and by comparing with eq. (1)

$$F = m^* a \Rightarrow m^* = \frac{\hbar^2}{d^2E/dk^2}$$

Is known as expression for m^* and it is dependent on E and K.

Origin of energy band formation in Solids:

The band theory of solid explains the formation of energy bands and determines whether a solid is a conductor, semiconductor or insulator.

The existence of continuous bands of allowed energies can be understood starting with the atomic scale. The electrons of a single isolated atom occupy atomic orbitals, which form a discrete set of energy levels.

When two identical atoms are brought closer, the outermost orbits of these atoms overlap and interact. When the wave functions of the electrons of different atoms begin to overlap considerably, the energy levels corresponding to those wave functions split.

If more atoms are brought together more levels are formed and for a solid of N atoms, each of the energy levels of an atom splits into N energy levels. These energy levels are so close that they form an almost continuous band. The width of the band depends upon the degree of overlap of electrons of adjacent atoms and is largest for the outermost atomic electrons.

As a result of the finite width of the energy bands, gaps are essentially left over between the bands called forbidden energy gap.

The electrons first occupy the lower energy levels (and are of no importance) then the electrons in the higher energy levels are of important to explain electrical properties of solids and these are called valence band and conduction band.

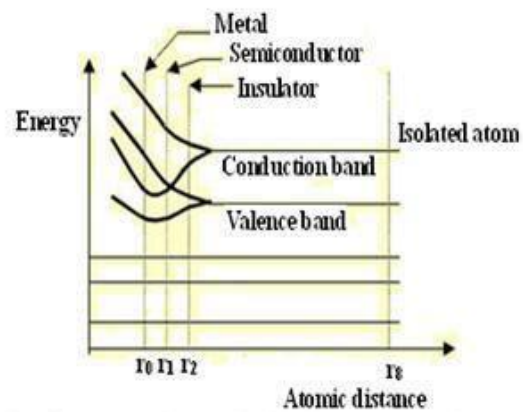
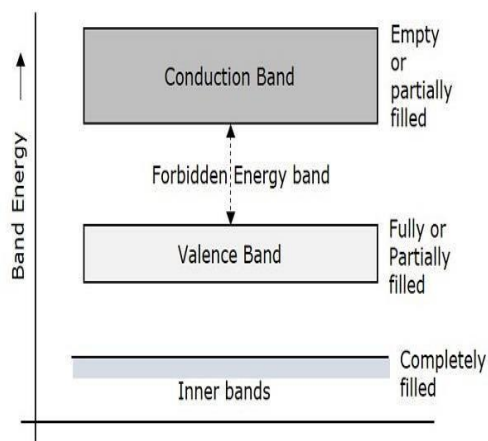


Fig: Formation of energy bands when atoms are closer

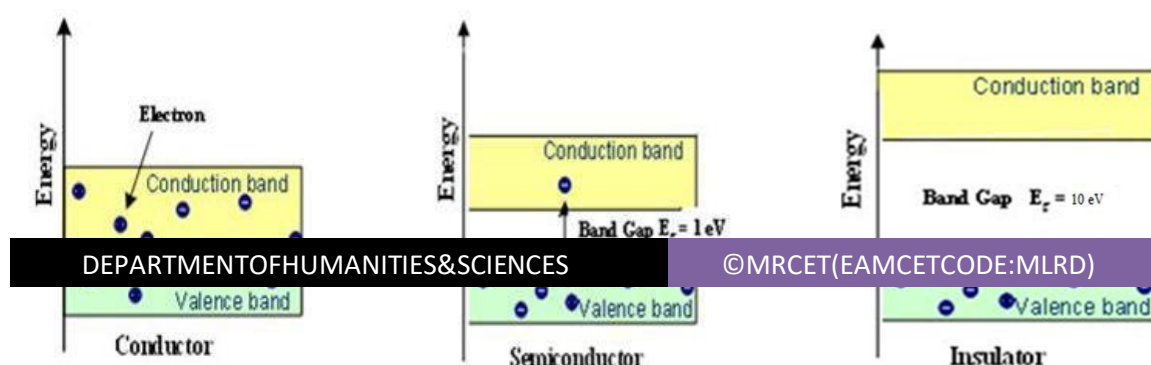
Valence band: A band occupied by valence electrons and is responsible for electrical, thermal and optical properties of solids and it is filled at 0K.

Conduction band: A band corresponding to outer most orbit is called conduction band and is the highest energy band and it is completely empty at 0K.

The forbidden energy gap between valence band conduction band is known as the energy band gap. By this solids are classified in to conductors, semiconductors and insulators.

Classification of solids into conductors, Semiconductors & Insulators:

Based on the energy band diagram materials or solids are classified as follows:



Conductors: In this kind of materials, there is no forbidden gap between the valence band and conduction band. It is observed that the valence band overlaps with the conduction band in metals as shown in figure. There are sufficient numbers of free electrons, available for electrical conduction and due to the overlapping of the two bands there is an easy transition of electrons from one band to another band takes place, and there is no chance for the presence of holes. Resistivity of conductors is very small and it is very few milli ohm meters ($\Omega \text{ m}$).

Examples: All metals (Na, Mg, Al, Cu, Ni, Ag, Li, A etc)

Semiconductors:

In semiconductors, there is a band gap exists between the valence band and conduction band and it is very less and it is the order of -1 to 2 eV are known as semiconductors. It will conduct electricity partially at normal conditions. The electrical resistivity values are 0.5 to 10^3 ohm meter . Due to thermal vibrations within the solid, some electrons gain enough energy to overcome the band gap (or barrier) and behave as conduction electrons. Conductivity exists here due to electrons and holes.

Examples: Silicon, Germanium, GaAs.

Insulators:

In insulators, the width of forbidden energy gap between the valence band and conduction band is very large. Due to large energy gap, electrons cannot jump from V.B to C.B. Energy gap is of the order of $\sim 10 \text{ eV}$ and higher than semiconductors. Resistivity values of insulators are 10^7 to 10^{12} ohm-m . Electrons are tightly bound to the nucleus, no valence electrons are available.

Examples: Wood, rubber, glass.

1) UNIT -4

2) LASER AND FIBER OPTICS

3) Explain i) Metastable state ii) optical pumping iii) population inversion

Metastable state: The excited state, which has a long life time, is known as metastable state.

Optical pumping: This process is required to achieve population inversion and used in Ruby laser.

Pumping process is defined as: "The process which excites the atoms from ground state to excited state to achieve population inversion".

Population Inversion:

Generally, number of atoms in the ground state is greater than the number of atoms in higher energy states.

But in order to produce a laser beam, the minimum requirement is stimulated emission.

Stimulated emission takes place only if the number of atoms in the higher energy level is greater than the number of atoms in the lower energy level.

Simply population inversion is nothing but number of atoms in higher energy level is greater than the number of atom in lower energy level.

4) Define spontaneous and stimulated emission of radiation?

Spontaneous Emission: When an atom in the excited state emits a photon of energy ' $h\nu$ ' coming down to ground state by itself without any external agency, such an emission is called spontaneous emission. $\text{Atom}^* \rightarrow \text{Atom} + h\nu$.

Photons released in spontaneous emission are not coherent. Hence spontaneous emission is not useful for producing lasers.

Stimulated Emission: When an atom in the excited state, emits two photons of same energy ' $h\nu$ ' while coming to down to ground state with the influence of an external agency, such an emission is called stimulated emission. $\text{Atom}^* \rightarrow \text{Atom} + 2h\nu$.

- In the two photons one photon induces the stimulated emission and the second one is released by the transition of atom from higher energy level to lower energy level.
- Both the photons are strictly coherent. Hence stimulated emission is responsible for laser production.

5) Explain the basic principle of optical fiber?

- Optical fibers are the waveguides through which electromagnetic waves of optical frequency range can be guided through them to travel long distances.
- An optical fiber works on the principle of total internal reflection (TIR).
- **Total Internal Reflection:** when a ray of light travels from a denser medium into a rarer medium and if the angle of incidence is greater than the critical angle then the light gets totally reflected into the denser medium

6) Explain i) Numerical Aperture ii) Acceptance angle

i) Numerical Aperture:

Numerical aperture of a fiber is a measure of its light gathering power.

"The Numerical Aperture (NA) is defined as the sine of the maximum acceptance angle" The light gathering ability of optical fiber depends on two factors i.e., (i) Core diameter (ii) NA

NA is defined as sine of the acceptance angle i.e.,

$$NA = \sin \theta_A \quad \text{i.e.} \quad NA = \sqrt{n_1^2 - n_2^2}$$

The efficiency of optical fiber is expressed in terms of NA; it is called as figure of merit of optical fiber.

ii) Acceptance Angle:

All right rays falling on optical fiber are not transmitted through the fiber.

Only those light rays making $\theta_i > \theta_c$ at the core-cladding interface are transmitted through the fiber by undergoing TIR. For which angle of incidence, the refraction angle is greater than 90° will be propagated through TIR.

There by Acceptance Angle is defined as: The maximum angle of incidence to the axis of optical fiber at which the light ray may enter the fiber so that it can be propagated through TIR.

5. What are the main sections of optical fiber? Describe the step index optical fiber?

- An optical fiber consists of three (3) co-axial regions.
- The innermost region is the light-guiding region known as "Core". It is surrounded by a middle co-axial regional known as "cladding". The outer most regions which completely covers the core & cladding regions is called "sheath or buffer jacket".
- Sheath protects the core & cladding regions from external contaminations, in addition to providing mechanical strength to the fiber.
- The refractive index of core (n_1) is always greater than the refractive index of cladding (n_2) i.e., $n_1 > n_2$ to observe the light propagation structure of optical fiber.

Step Index optical fiber:

- Based on variation in the core refractive index (n_1), optical fibers are divided into two types
 1. Step index fiber
 2. Graded index fiber
- Step index fibers have both single & multimode propagations.

6) Write a short note on attenuation in optical fibers.

Usually, the power of light at the output end of optical fiber is less than the power launched at the input end, then the signal is said to be attenuated.

Attenuation: It is the ratio of input optical power (P_i) into the fiber to the power of light coming out at the output end (P_o).

Attenuation coefficient is given as, $\alpha = 10/L \log_{10} P_i/P_o$ db/km.

Attenuation is mainly due to

1. Absorption.
2. Scattering.
3. Bending.

7) Write down advantages of fiber optics in communication system Or What are the Advantages of optical fibers over metallic cables?

- Optical fibers allow light signals of frequencies over a wide range and hence greater volume of information can be transmitted either in digital form or in analog form within a short time.

- In metallic cables only 48 conversations can be made at once without cross talks where as in optical fibers more than 15000 conversations can be made at once without cross talks.
- Light cannot enter through the surface of the optical fiber except at the entry interface i.e., interference between different communication channels is absent. Hence purity of light signal is protected.
- Optical signal do not produce sparks like electrical signals and hence it is safe to use optical fibers.
- External disturbances from TV or Radio Stations power electronic systems and lightning cannot damage the signals as in case of metallic cables.
- Materials used in the manufacture of optical fibers are SiO_2 , plastic, glasses which are cheaper & available in plenty.

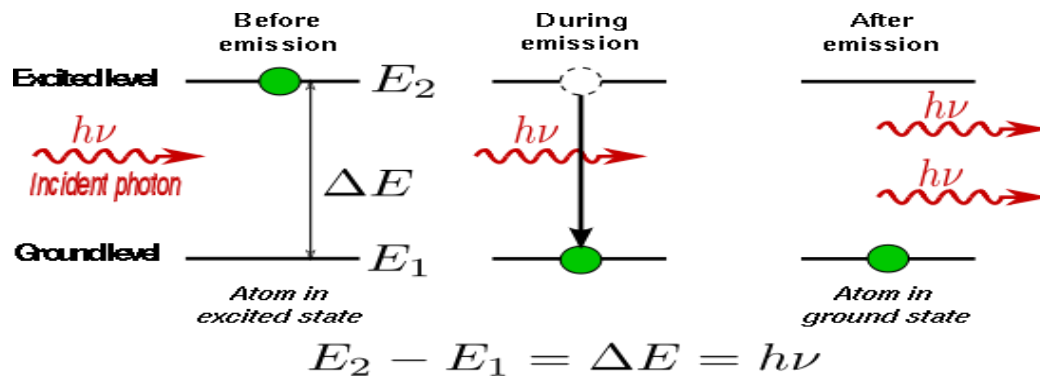
Part-B(Descriptive-10marks)

1) What is Acronym of a Laser, absorption, spontaneous and stimulated emissions?

- **Laser:** Laser means **Light Amplification by Stimulated Emission of Radiation**.
- **Absorption:** When an atom absorbs an amount of energy ' $h\nu$ ' in the form of photon from the external agency and excited into the higher energy levels from ground state, then this process is known as absorption. $\text{Atom} + h\nu \rightarrow \text{Atom}^*$
- **Spontaneous Emission:** When an atom in the excited state emits a photon of energy ' $h\nu$ ' coming down to ground state by itself without any external agency, such an emission is called spontaneous emission. $\text{Atom}^* \rightarrow \text{Atom} + h\nu$
- Photons released in spontaneous emission are not coherent. Hence spontaneous emission is not useful for producing lasers.
- **Stimulated Emission:** When an atom in the excited state, emits two photons of same energy ' $h\nu$ ' while coming down to ground state with the influence of an external agency, such an emission is called stimulated emission. $\text{Atom}^* \rightarrow \text{Atom} + 2h\nu$
- In the two photons one photon induces the stimulated emission and the second one is released by the transition of atom from higher energy level to lower energy level.
- Both the photons are strictly coherent. Hence stimulated emission is responsible for laser production.

2) Explain principle of laser/lasing action?

- **Laser Production Principle:**
- Two coherent photons produced in the stimulated emission, interacts with other two excited atoms, resulting in four coherent photons.
- Thus, coherent photons are multiplied in a lasing medium. The continuous successive emission of photons results for the production of laser beam.

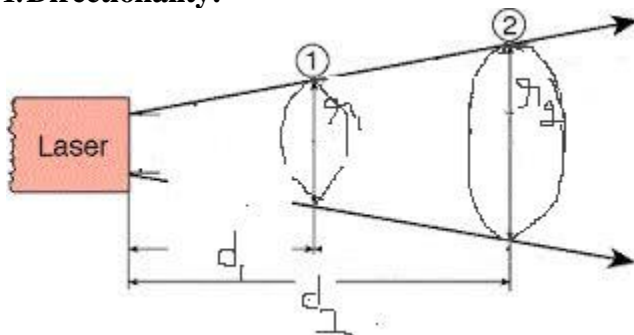


3) What are the characteristics/striking features/Properties of Laser Light?

Characteristic of Laser Beam: Some of the special properties which distinguish lasers from ordinary light sources are characterized by:

1. Directionality
2. High Intensity
3. Mono-chromaticity
4. Coherence

1. Directionality:



Laser emits radiation only in one direction. The directionality of laser beam is expressed in terms of angle of divergence (ϕ)

Divergence or Angular Spread is given by $\phi = \frac{r_2 - r_1}{d_2 - d_1}$

Where d_1, d_2 are any two distances from the laser source emitted and r_1, r_2 are the radii of beam spots at a distance d_1 and d_2 respectively as shown in above figure. Laser light having less divergence, it means that laser light having more directionality.

2. High Intensity: Generally, light from conventional sources spread uniformly in all directions. For example, take 100 watt bulb and look at a distance of 30 cm, the power enter into the eye is less than thousand of a watt. This is due to uniform distribution of light in all directions. But in case of lasers, light is a narrow beam and its energy is concentrated within the small region. The concentration of energy accounts for greater intensity of lasers.

3. Monochromaticity: The light emitted by laser is highly monochromatic than any of the other conventional monochromatic light. A comparison b/w normal light and laser beam, ordinary sodium (Na) light emits radiation at wavelength of 5893 \AA with the line width of 1 \AA . But He-Ne laser of wave length 6328 \AA with a narrow width of only 10^{-7} \AA i.e., Monochromaticity of laser is 10 million times better than normal light.

The degree of Monochromaticity of the light is estimated by line of width (spreading frequency of line).

4. Coherence: If any wave appears as pure sine wave for long time and infinite space, then it is said to be perfectly coherent.

Practically, no wave is perfectly coherent including lasers. But compared to other light sources, lasers have high degree of coherence because all the energy is concentrated within the small region. There are two independent concepts of coherence.

i) Temporal coherence (criteria of time)

ii) Spatial coherence (criteria of space)

4) Explain the concept of population inversion and pumping in lasers?

Population Inversion:

- Generally, number of atoms in the ground state is greater than the number of atoms in higher energy states.
- But in order to produce a laser beam, the minimum requirement is stimulated emission.
- Stimulated emission takes place only if the number of atoms in the higher energy level is greater than the number of atoms in the lower energy level.
- Simply population inversion is nothing but number of atoms in higher energy level is greater than the number of atoms in lower energy level.
- So, if there is a population inversion thereby only stimulated emission will be able to produce laser beam.
- Population inversion is associated with three phenomena.
 - Stimulated emission
 - Amplification
 - Pumping process
- Stimulated Emission: If majority of atoms are present in higher energy state than the process becomes very easy.
- Amplification: If 'N₁', represents number of atoms in the ground state and 'N₂' represents number of atoms in the excited state then the amplification of light takes place only when N₂ > N₁.
- If N₂ > N₁, there will be a population inversion so induced beam and induced emission are in the same directions and strictly coherent than the resultant laser is said to be amplified.
- Boltzmann's principle gives the information about the fraction of atoms found on average in any particular energy state at equilibrium temperature as
- $\frac{N_1}{N_2} = \exp\left(\frac{E_2 - E_1}{KT}\right) = \exp\left(\frac{\Delta E}{KT}\right)$

$$\frac{N_1}{N_2} = \exp\left(\frac{h\nu}{KT}\right)$$

Pumping Process:

- This process is required to achieve population inversion.
- Pumping process is defined as: "The process which excites the atoms from ground state to excited state to achieve population inversion".
- Pumping can be done by number of ways
 - i) Optical Pumping □ excitation by strong source of light (flashing of a camera)
 - ii) Electrical Pumping □ excitation by electron impact
 - iii) Chemical Pumping □ excitation by chemical reactions
 - iv) Direct Conversion □ Electrical energy is directly converted into

radiant energy in devices like LED's, population Inversion is achieved in forward bias.

5) What are Einstein's coefficients and explain the relation among them?

or

Derive the relation between the probabilities of spontaneous emission and stimulated emission in terms of Einstein's coefficient?

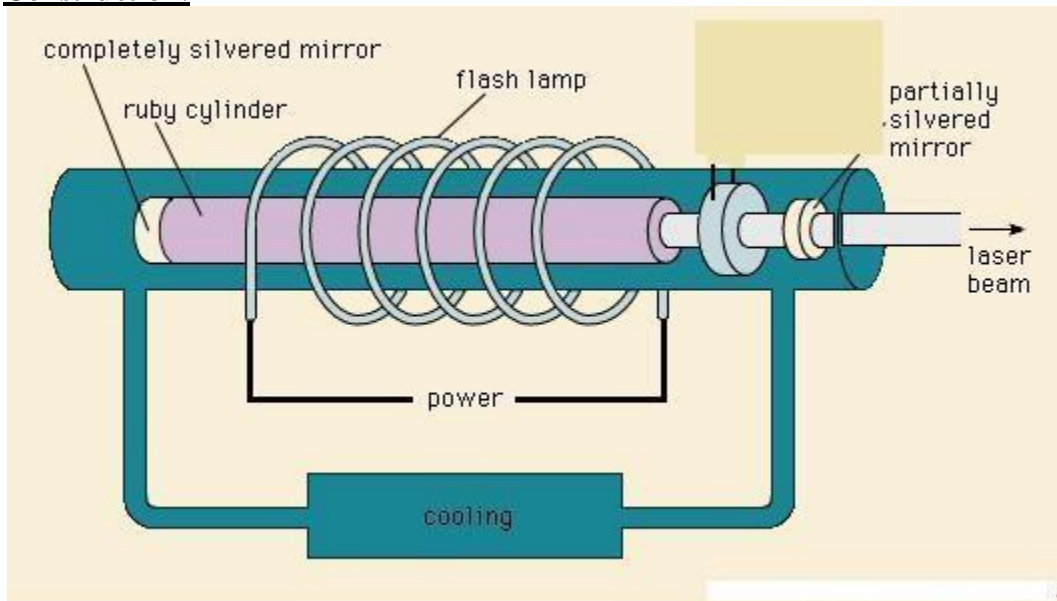
Einstein's Theory of Radiation:

- In 1917, Einstein predicted the existence of two different kinds of processes by which an atom emits radiation.
- Transition b/w the atomic energy states is statistical process. It is not possible to predict which particular atom will make a transition from one state to another state at a particular instant. For an assembly of very large number of atoms it is possible to calculate the rate of transitions b/w two states.
- Einstein was the first to calculate the probability of such transition, assuming the atomic system to be in equilibrium with electromagnetic radiation.
- The number of atoms excited during absorption in the time ' Δt ' is given by: $N_{ab} = Q N_1 B_{12} \Delta t$, Where N_1 = number of atoms in state ' E_1 ', Q = Energy density of induced beam and B_{12} = Probability of an absorption transition coefficient.
- The number of spontaneous transitions N_{sp} taking place in time ' Δt ' depends on only no. of atoms N_2 lying in excited state. $N_{sp} = A_{21} N_2 \Delta t$, Where A_{21} = probability of spontaneous transition.
- The number of stimulated transitions N_{st} occurring during the time Δt may be written as: $N_{st} = B_{21} N_2 \Delta t$, Where B_{21} = probability of stimulated emission.
- Under the thermal equilibrium number of upward transitions = number of downward transitions per unit volume per second.
- So, we can write: $A_{21} N_2 + B_{21} N_2 Q = B_{12} N_1 Q$ 1
- $Q = A_{21} N_2 / (B_{12} N_1 - B_{21} N_2)$ 2
- Dividing by $B_{21} N_2$ in all terms, $Q = (A_{21}/B_{21}) \times 1 / (B_{12} N_1 / B_{21} N_2 - 1)$ 3
- By substituting $N_1/N_2 = \exp(h\nu/kT)$ from Boltzmann Distribution law,
- $Q = (A_{21}/B_{21}) 1 / (B_{12}/B_{21} \exp(h\nu/kT) - 1)$ 4
- Above equation must agree with Planck's energy distribution – radiation formula. $Q = \frac{h\nu^3}{\pi^2 C^3} \frac{1}{\exp(h\nu/kT) - 1}$ 5
- From equations 4 & 5, $B_{12} = B_{21}$, we get $A_{21}/B_{21} = \frac{h\nu^3}{\pi^2 C^3}$
- The coefficients A_{21} , B_{12} , B_{21} are known as Einstein coefficients.
- Note: Since we are applying same amount of energy (Q) and observing in the same time (Δt), number of atoms excited into higher energy levels (absorption) = number of atoms that made transition into lower energy levels (stimulated emission)
 $B_{12} = B_{21}$ i.e. absorption = stimulated emission

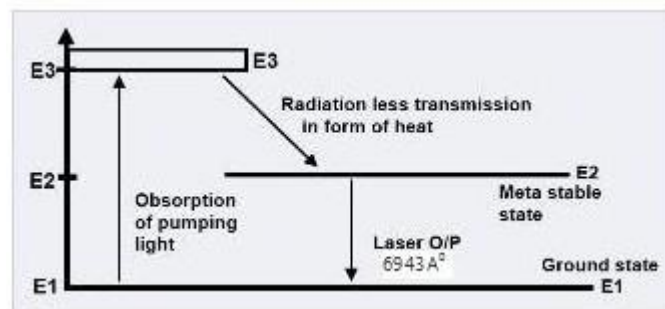
6) Describe the principle, construction and working of ruby laser with relevant energy level diagram?

- **Ruby Laser:** It is a 3 level solid state laser, discovered by Dr. T. Maiman in 1960.
- Principle:
- The chromium ions are raised to excited states by optical pumping using xenon flash lamp
- Then the atoms are accumulated at metastable state by non-radiative transition.

- Due to stimulated emission the transition of atoms take place from metastable state to ground state, there by emitting laser beam.
- **Construction:**



- Ruby is a crystal of aluminum oxide (Al_2O_3) in which some of the aluminum ions (Al^{3+}) is replaced by chromium ions (Cr^{3+}). This is done by doping small amount (0.05%) of chromium oxide (Cr_2O_3) in the melt of purified Al_2O_3 .
- These chromium ions give the pink color to the crystal. Laser rods are prepared from a single crystal of pink ruby. Al_2O_3 does not participate in the laser action. It only acts as the host.
- The ruby crystal is in the form of cylinder. Length of ruby crystal is usually 2 cm to 30 cm and diameter 0.5 cm to 2 cm.
- The ends of ruby crystal are polished, grounded and made flat.
- The one of the ends is completely silvered while the other one is partially silvered to get the efficient output. Thus the two polished ends act as optical resonator system.
- A helical flash lamp filled with xenon is used as a pumping source. The ruby crystal is placed inside a xenon flash lamp. Thus, optical pumping is used to achieve population inversion in ruby laser.
- As very high temperature is produced during the operation of the laser, the rod is surrounded by liquid nitrogen to cool the apparatus.
- **Working with Energy Level Diagram (ELD):**

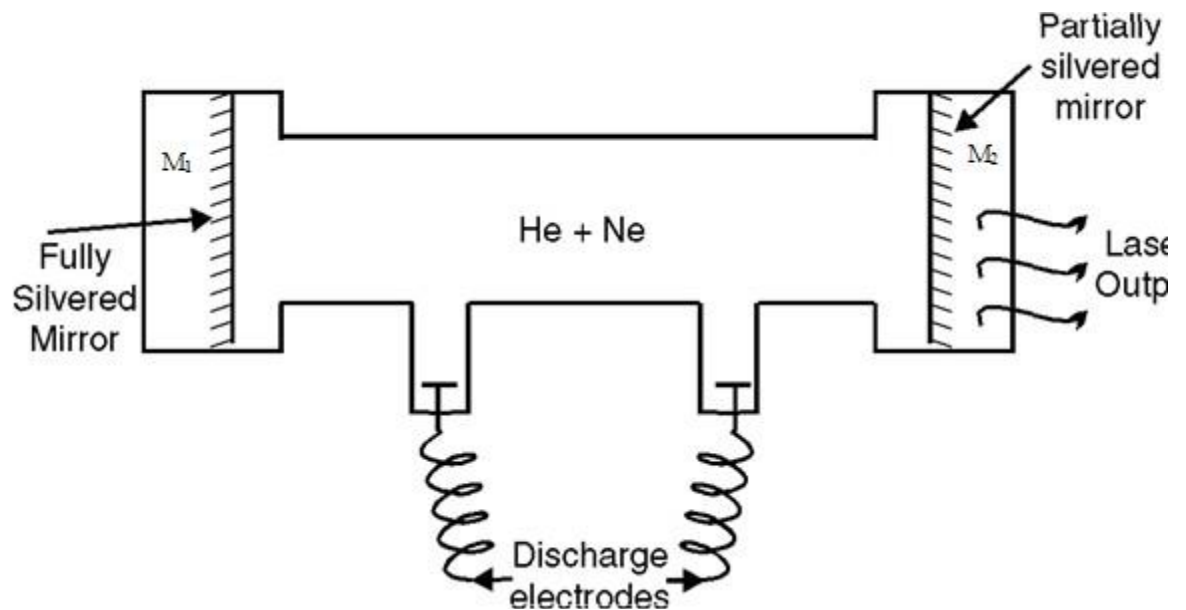


▪ Fig: Energy Level Diagram of Ruby Laser

- The flash lamp is switched on, a few thousand joules of energy is discharged in a few milliseconds.
- A part of this energy excites the Cr^{3+} Ions to excited state from their ground state and the rest heats up the apparatus can be cooled by the cooling arrangement by passing liquid nitrogen.
- The chromium ions respond to this flash light having wavelength 5600 \AA^0 (Green), $[4200 \text{ \AA}^0$ (Red) Also]
- When the Cr^{3+} Ions are excited to energy level E_3 from E_1 the population in E_3 increases.
- Cr^{3+} Ions stay here (E_3) for a very short time of the order of 10^{-8} sec, then they drop to the level E_2 which is metastable state of lifetime 10^{-3} sec. Here the transitions from E_3 to E_2 is non radiative in nature.
- As the lifetime of the state E_2 is much longer, the number of ions in this state goes on increasing while in the ground state (E_1) goes on decreasing. By this process population inversion is achieved between E_2 & E_1 .
- When an excited ion passes spontaneously from the metastable state E_2 to the ground state E_1 it emits a photon of wavelength 6943 \AA^0 .
- This photon travels through the ruby rod and if it is moving parallel to the axis of the crystal, is reflected back & forth by silvered ends until it stimulates an excited ion in E_2 and causes it to emit fresh photon in phase with the earlier photon. This stimulated transition triggers the laser Transition.
- The process is repeated again and again, because the photons repeatedly move along the crystal being reflected from ends. The photons thus get multiplied.
- When the photon beam becomes sufficiently intense, such that a part of it emerges through the partially silvered end of the crystal.

7) Describe the principle, construction and working of He-Ne laser with relevant energy level diagram?

- **He-Ne Laser:**
- **Principle:** This laser is based on the principle of stimulated emission, produced in the active medium of gas. Here, the population inversion achieved due to the interaction between the two gases which have closer higher energy levels.
- **Construction:**



▪ Fig:He-Ne laser

- The first gas laser to be operated successfully was the He-Ne laser in 1961 by the scientist A. Jawan.
- In this method, two gases helium & Neon were mixed in the ratio 10:1 in a discharge tube made of quartz crystal.
- The dimensions of the discharge tube are nearly 80 cm length and 1.5 cm diameter, with its windows slanted at Brewster's angle i.e., $\theta = \tan^{-1}(n)$, Where n = refractive index of the window substance.
- The purpose of placing Brewster windows on either side of the discharge tube is to get plane polarized laser output.
- Two concave mirrors M_1 & M_2 are made of dielectric material arranged on both sides of the discharge tube so that their foci lines within the interior of discharge tube.
- One of the two concave mirrors M_1 is thick so that all the incident photons are reflected back into lasing medium.
- The thin mirror M_2 allows part of the incident radiation to be transmitted to get laser output.
- Working:

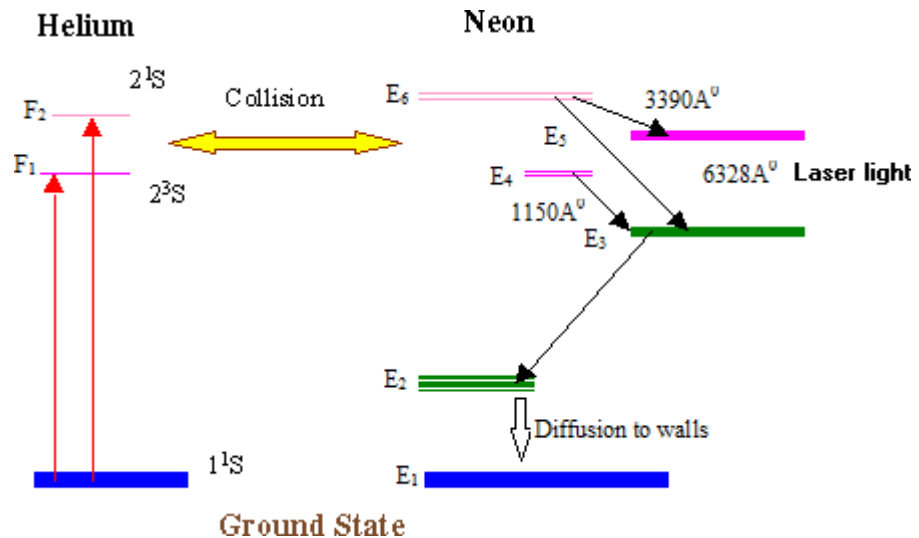


Fig: (E.L.D) Energy Level Diagram corresponding to He-Ne laser

- The discharge tube is filled with Helium at a pressure of 1 mm of Hg & Neon at 0.1 mm of Hg.
- When electric discharge is set-up in the tube, the electrons present in the electric field make collisions with the ground state He atoms.
- Hence ground state He atoms get excited to the higher energy levels $F_1(2S_1), F_2(2S_3)$.
- Here Ne atoms are active centers.
- The excited He atoms make collision with the ground state Ne atoms and bring the Ne atoms into the excited states E_4 & E_6 .
- The energy levels E_4 & E_6 of Ne are the metastable states and the Ne atoms are directly pumped into these energy levels.
- Since the Ne atoms are excited directly into the levels E_4 & E_6 , these energy levels are more populated than the lower energy levels E_3 & E_5 .
- Therefore, the population inversion is achieved between E_6 & E_5, E_6 & E_3, E_4 & E_3 .
- The transition between these levels produces wavelengths of $3390\text{\AA}^0, 6328\text{\AA}^0, 1150\text{\AA}^0$ respectively.
- Now the Ne atoms undergo transition from E_3 to E_2 and E_5 to E_2 in the form of fast decay giving photons by spontaneous emission. These photons are absorbed by optical elements placed inside the laser system.
- The Ne atoms are returned to the ground state (E_1) from E_2 by non-radiative diffusion and collision process, therefore there is no emission of radiation.
- Some optical elements placed inside the laser system are used to absorb the IR laser wavelengths $3390\text{\AA}^0, 1150\text{\AA}^0$.
- Hence the output of He-Ne laser contains only a single wavelength of 6328\AA^0 .
- The released photons are retransmitted through the concave mirror M_2 thereby producing laser.
- A continuous laser beam of red color at a wavelength of 6328\AA^0 .
- By the application of large potential difference, Ne atoms are pumped into higher energy levels continuously.
- A laser beam of power 0.5 to 50 MW comes out from He-Ne laser.

8) Describe the principle, construction and working of Semiconductor laser with relevant energy level diagram?

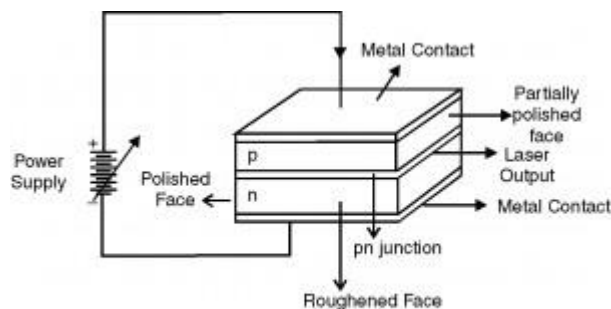
Semiconductor Laser:

- Semiconductor lasers are of two types, Except the Construction, Principle and working are same for both.
 1. Homojunction semiconductor Laser
 2. Heterojunction semiconductor Laser

Principle:

- After the invention of semiconductor laser in 1961, laser have become at common use.
- In conventional lasers, lasers are generated due to transition of electrons from high to lower energy level.
- But in semiconductor laser the transition takes place from conduction band to valence band.
- The basic mechanism responsible for light emission from a semiconductor laser is the recombination of e⁻ and holes at PN-junction when current is passed through the diode.
- Stimulated emission can occur when the incident radiation stimulates an electron in conduction band to make a transition into valence band in that process radiation will be emitted.
- When current is passed through PN – junction under forward bias, the injected e⁻s & holes will increase the density of e⁻ in CB & holes in VB. At some value of current the stimulated emission rate will exceed the absorption rate.
- As the current is further increased at some threshold value of current the amplification will take place and laser begin to emit coherent radiation.
- The properties of semiconductor laser depends upon the energy gap

Fabrication/construction:



• Fig: Homojunction Semiconductor Laser

Homojunction Semiconductor Laser:

- Ga-As is heavily doped with impurities in both P & N regions. N region is doped with tellurium & P – region by Germanium.
- The concentration of doping is of the order of 10^{17} to 10^{19} impure atoms per cm.
- The size of the diode is small i.e., 1 mm each side & the depletion layer's thickness varies from 1 to 100 μm .
- These values depend on diffusion condition and 40 mp at the time of fabrication.

Heterojunction Semiconductor Laser:

- Hetrojunction means the material on one side of the junction differs from that on the other side. i.e; Ga-As on one side and GaAlAs on other side.
- Generation and recombination takes place very fastly.

Working:

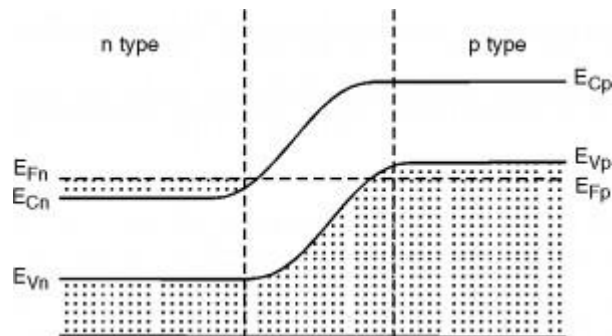


Fig a) when no biasing

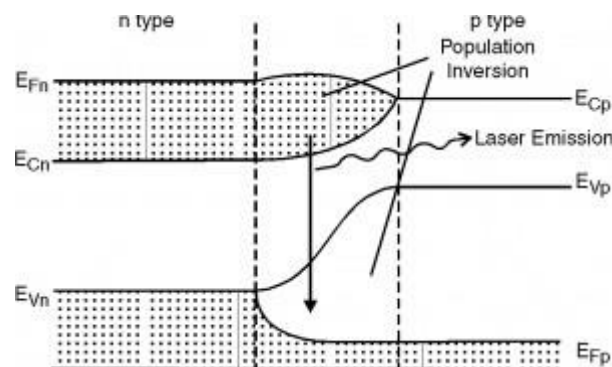


Fig b) with biasing

- When a forward bias with the source is applied to a semiconductor, e^- from N-region & holes from P-region move to cross the junction in opposite directions.
- In natural region e^- 's & holes combine recombination is possible due to transition of e^- from CB to VB.
- For low currents the population inversion does not take place hence only spontaneous emission takes place and photon released are not coherent.
- When forward current is further increased beyond the certain threshold value population inversion takes place and coherent photons are released.
- The energy gap of Gallium Arsenide (Ga-As) is 1.487 eV and corresponding wavelength of radiation is 6435 \AA which is responsible for laser emission.

9) Mention some important applications of Lasers in various fields?

Applications of Lasers: Lasers have wide applications in different branches of science and engineering because of the following.

- Very narrow bandwidth
- High directionality
- Extreme brightness

1. Communication:

- Lasers are used in optical communications, due to narrow bandwidth.
- The laser beam can be used for the communication b/w earth & moon (or) other satellites due to the narrow angular spread.
- Used to establish communication between submarines i.e; underwater communication.

2. Medical:

- Identification of tumors and curification.
- Used to detect and remove stones in kidneys.
- Used to detect tumors in brain.

3. Industry:

- Used to make holes in diamond and hard steel.
- Used to detect flaws on the surface of aeroplanes and submarines.

4. Chemical & Biological:

- Lasers have wide chemical applications. They can initiate or fasten chemical reactions.
- Used in the separation of isotopes.
- Lasers can be used to find the size & shape of biological cells such as erythrocytes.

10. with the help of a suitable diagram explain the principle, structure and working of an optical fiber as a wave guide?

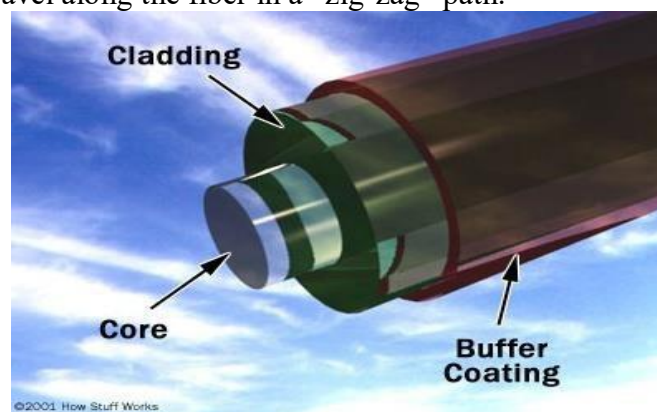
Principle: Optical fibers are the waveguides through which electromagnetic waves of optical frequency range can be guided through them to travel long distances.

- An optical fiber works on the principle of total internal reflection (TIR).

Total Internal Reflection: when a ray of light travels from a denser medium into a rarer medium and if the angle of incidence is greater than the critical angle then the light gets totally reflected into the denser medium.

Structure & Working:

- An optical fiber consists of three (3) co-axial regions.
- The innermost region is the light-guiding region known as "Core". It is surrounded by a middle co-axial region known as "cladding". The outermost region which completely covers the core & cladding regions is called "sheath or buffer jacket".
- Sheath protects the core & cladding regions from external contaminations, in addition to providing mechanical strength to the fiber.
- The refractive index of core (n_1) is always greater than the refractive index of cladding (n_2) i.e., $n_1 > n_2$ to observe the light propagation structure of optical fiber.
- When light enters through one end of optical fiber it undergoes successive total internal reflections and travel along the fiber in a "zig-zag" path.



11) Define and derive the expressions for acceptance angle and numerical Aperture?

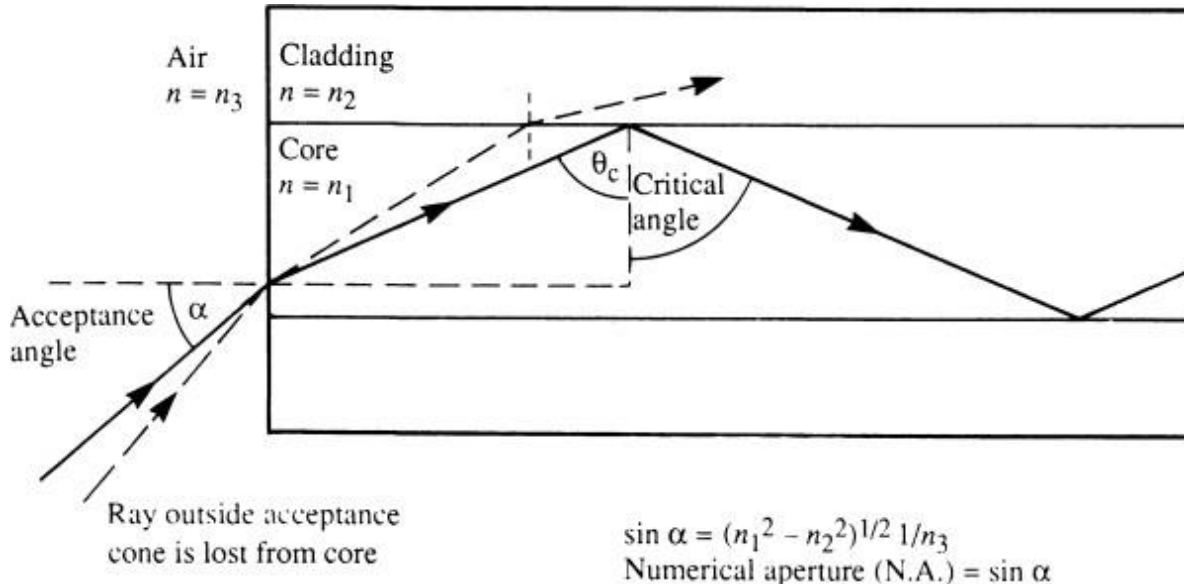
Expressions for acceptance angle & Numerical Aperture (NA):

Acceptance Angle:

- All light rays falling on optical fiber are not transmitted through the fiber. Only those light rays making $\theta_i > \theta_c$ at the core-cladding interface are transmitted through the

fiber by undergoing TIR. For which angle of incidence, the refraction angle is greater than 90° will be propagated through TIR.

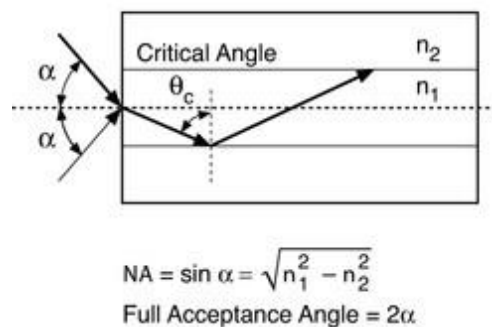
- Thereby Acceptance Angle is defined as: The maximum angle of incidence to the axis of optical fiber at which the light ray may enter the fiber so that it can be propagated through TIR.



- Consider the optical fiber with core refractive index n_1 and cladding refractive index n_2 . Light is incident on the axis of optical fiber at an angle θ_1 . It can produce an angle of refraction θ_2 .
- The relationship at the interface is given by Snell's law as:
At air-core interface (A), $n_0 \sin \theta_1 = n_1 \sin \theta_2$ 1
At core-clad interface (B), for TIR, $n_1 \sin(90^\circ - \theta_2) = n_2 \sin 90^\circ$
 $n_1 \cos \theta_2 = n_2$ 2
- Eq'n 1 can be written as, $n_0 \sin \theta_1 = n_1 \sqrt{1 - \cos^2 \theta_2}$ 3
- Substituting 2 in 3, $n_0 \sin \theta_1 = n_1 \sqrt{1 - (n_2/n_1)^2}$
- For air $n_0 = 1$, then $\sin \theta_1 = \sqrt{n_1^2 - n_2^2}$
- $\theta_1 = \theta_A = \sin^{-1} \sqrt{n_1^2 - n_2^2}$, Here θ_A is called Acceptance angle
- This gives max value of external incident angle for which light will propagate in the fiber.

Numerical Aperture (NA):

Numerical Aperture



- Numerical aperture of a fiber is a measure of its light gathering power.
- “The Numerical Aperture (NA) is defined as the sine of the maximum acceptance angle”
- The light gathering ability of optical fiber depends on two factors i.e., Core diameter & NA.
- NA is defined as sine of the acceptance angle i.e., $NA = \sin \theta_A$
 $NA = \sqrt{n_1^2 - n_2^2}$
- The efficiency of optical fiber is expressed in terms of NA, so it is called as figure of merit of optical fiber.
- # NA is also expressed like this: $NA = \sqrt{n_1^2 - n_2^2} = \sqrt{(n_1 - n_2)(n_1 + n_2)}$
- Fractional index change $\Delta = \frac{n_1 - n_2}{n_1} = n_1 - n_2 = \Delta n_1$
 Then $NA = \sqrt{\Delta n_1 (n_1 + n_2)}$
 Let $n_1 = n_2$, then $n_1 + n_2 = 2n_1$
- Then $NA = \sqrt{\Delta n_1 - 2n_1} = n_1 \sqrt{2\Delta} = n_1 \sqrt{2\Delta}$

12) How optical fibers are classified on the basis of refractive index profile?

Or

Describe the Step index and graded index optical fibers in detail and explain the transmission of signal through them?

Classification of Optical Fibers:

- Based on variation in the core refractive index (n_1), optical fibers are divided into two types
 1. Step index fiber
 2. Graded index fiber
- Based on mode of propagation, fibers are further classified into
 1. Single mode propagation
 2. Multimode propagation
- Step index fibers have both single & multimode propagations.
- Graded index fibers have multimode propagation only
- Altogether into total three (3) types of fibers
 1. Single mode step index fiber
 2. Multimode step index fiber
 3. Multimode graded index fiber

Transmission of Signal in Optical Fibers:

1. Step Index Fiber: The refractive index of core material is uniform throughout and undergoes a sudden change in the form of step at the core-clad interface.

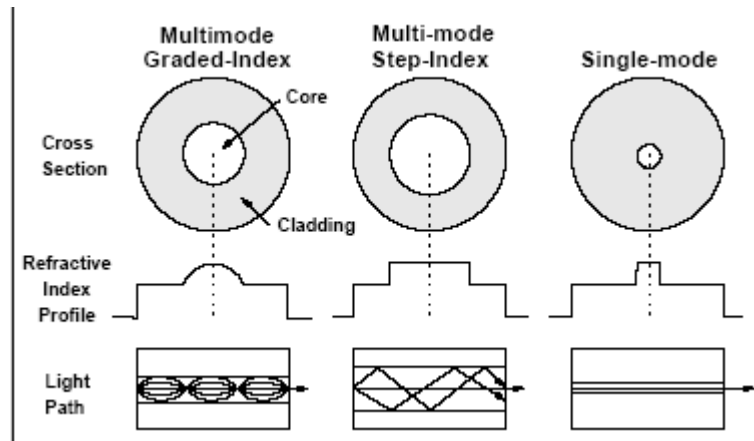


Fig:Refractiveindexprofile&propagationinsinglemode,stepindex&gradedindexfibers

a) SingleModeStepIndexFiber:

- The variation of the refractive index of a step index fiber as a function of distance can be mathematically represented as longitudinal cross section.

Note: Mode of propagation: It is defined as the number of paths available for the light ray to transfer through the optical fiber.

Structure:

- i) CoreDiameter: 8 to 12 μm , usually 8.5 μm
- ii) CladdingDiameter: Around 125 μm
- iii) SheathDiameter: 250 to 1000 μm
- iv) NA: 0.08 to 0.15 usually 0.10

PerformanceCharacteristics:

- i) BandWidth: Greater than 500 MHzKm.
- ii) Attenuation: 2 to 5 dB/ Km.
- iii) Applications: These fibers are ideally suited for high band width applications using single mode injection coherent (LASER) sources.

B) MultiModeStepIndex Fibers:

- These fibers have reasonably large core diameters and large NA to facilitate efficient transmission to incoherent or coherent light sources.
- These fibers allow finite number of modes.
- Normalized frequency (NF) is the cut off frequency, below which a particular mode cannot exist. This is related to NA, Radius of the core, and wave length of light as

$$NF = 2\pi/\lambda a(NA), \text{ Where } a = \text{radius of core}$$

Structure:

- i) CoreDiameter: 50 to 200 μm
- ii) CladdingDiameter: 125 to 400 μm
- iii) SheathDiameter: 250 to 1000 μm
- iv) NA: 0.16 to 0.5

PerformanceCharacteristics:

- i) BandWidth: 6 to 50 MHzKm.
- ii) Attenuation: 2.6 to 50 db/km.
- iii) Applications: These fibers are ideally suited for limited bandwidth and relatively low cost applications.

c) MultiModeGradedIndexFibers:

- In case of graded index fibers, the refractive index of core is made to vary as a function of radial distance from the centre of the optical fiber.

- Refractive index increases from one end of core diameter to center and attains maximum value at the centre. Again refractive index decreases as moving away from center to towards the other end of the core diameter.
- Therefractiveindexvariationisrepresentedas $n(r) = n_1(1 - 2\Delta)^{1/2} = n_2$ Here $\Delta = \text{fractional change in refractive index} = (n_1 - n_2)/n_1$
- Thenumberofmodesisgivenbytheexpression $N = 4.9[d(\text{NA})/\lambda]^2$
Where d = core diameter, λ = wavelength of radiation

Structure:

- i) CoreDiameter: 30 to 100 μm
- ii) CladdingDiameter: 105 to 150 μm
- iii) SheathDiameter: 250 to 1000 μm
- iv) NA: 0.2 to 0.3

PerformanceCharacteristics:

- i) BandWidth: 300 MHzK to 3 GHzK.
- ii) Attenuation: 2 to 10 dB/km.
- iii) Applications: These are ideally suited for medium to high band width applications using incoherent and coherent multimode sources.

13) Distinguish Step index & Graded index fibers And Single mode & Multimode fibers?

Step Index	Graded Index
1. RI of core is uniform throughout except at one stage. 2. Single & multimode propagation exist. 3. Used for short distance applications. 4. Attenuation losses are of the order 100 dB/km. 5. Mer 4 in rays propagation takes place. 6. Easy to manufacture.	1. Refractive index varies gradually with radial distance. 2. It is a multimode fiber. 3. Used for long distance applications. 4. Attenuation losses are of the order 10 dB/km. 5. Skew rays propagation takes place. 6. Difficult to manufacture.
Single Mode	Multi Mode
1. Core diameter is small. 2. Signal entry is difficult. 3. Exists in step index fiber. 4. Light must be coherent.	1. Core diameter is large. 2. Signal entry is easy. 3. Exists in both step & graded index fibers. 4. Light source may be coherent or incoherent source.

14) What are the Advantages of optical fibers over metallic cables?

- Optical fibers allow light signals of frequencies over a wide range and hence greater volume of information can be transmitted either in digital form or in analog form within a short time.
- In metallic cables only 48 conversations can be made at once without cross talks where as in optical fibers more than 15000 conversations can be made at once without cross talks.
- Light cannot enter through the surface of the optical fiber except at the entry interface i.e., interference b/w different communication channels is absent. Hence purity of light signal is protected.

- Optical signal do not produce sparks like electrical signals and hence it is safe to use optical fibers.
- External disturbances from TV or Radio Stations power electronics systems and lightning cannot damage the signals as in case of metallic cables.
- Materials used in the manufacture of optical fibers are SiO_2 , plastic, glasses which are cheaper & available in plenty.

15) How optical fibers are used in communication field? Or Explain optical fiber communication link with help of block diagram.

Optical Fiber Communication Link:

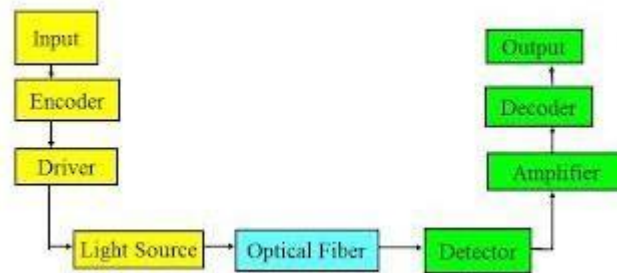


Fig: Block Diagram of Optical fiber communication link

Optical fiber is an ideal communication medium by systems that require high data capacity, fast operation and to travel long distances with a minimum number of repeaters.

Encoder: It is an electronics system that converts the analog information signals, such as voice of telephone user, in to binary data. The binary data consists of series of electrical pulses.

Transmitter: Transmitter consists of a driver which is a powerful amplifier along with light source. The o/p of amplifier feeds to light source, which converts electrical pulses in to light pulses.

Source to Fiber Connector: It is a special connector that sends the light from source to fiber. The connector acts as temporary joint b/w the fiber and light source, misalignment of this joint, leads to loss of signal.,

Fiber to Detector Connector: It is also temporary joint, which collects the source from fiber.

Receiver: Receiver consists of a detector followed by amplifier. This combination converts light pulses in to electrical pulses.

Decoder: Electrical pulses containing information are fed to the electronic circuit called decoder. Decoder converts binary data of electrical pulses in to analog information signals.

16) Write a short note on attenuation in optical fibers.

Usually, the power of light at the output end of optical fiber is less than the power launched at the input end, then the signal is said to be attenuated.

Attenuation: It is the ratio of input optical power (P_i) into the fiber to the power of light coming out at the output end (P_o).

Attenuation coefficient is given as, $\alpha = 10/L \log_{10} P_i/P_o$ db/km.

Attenuation is mainly due to

1. Absorption.
2. Scattering.
3. Bending.

1. Absorption Losses: In glass fibers, three different absorption take place.

Ultra violet absorption: Absorption of UV radiation around $0.14\mu\text{m}$ results in the ionization of valence electrons.

Infrared absorption: Absorption of IR photons by atoms within the glass molecules causes heating. This produces absorption peak at $8\mu\text{m}$, also minor peaks at 3.2 , 3.8 and $4.4\mu\text{m}$.

Ion resonance/OH⁻ absorption: The OH⁻ ions of water, trapped during manufacturing causes absorption at 0.95 , 1.25 and $1.39\mu\text{m}$.

2. Scattering Losses:

The molten glass, when it is converted into thin fiber under proper tension creates sub microscopic variations in the density of glass leads to losses.

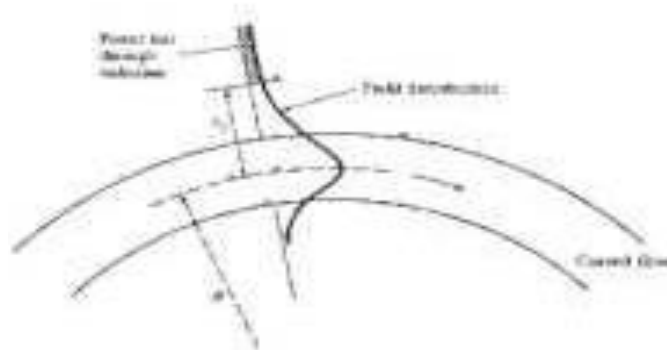
The dopants added to the glass to vary the refractive index also lead to the inhomogeneities in the fiber. As a result losses occur.

Scattering losses are inversely proportional to fourth power of λ . (λ^4)

3. Bending Losses:

In a bent fiber, there is a loss in power of the transmitted signal called as Bending Loss. According to the theory of light, the part of the wavefront travelling in cladding (rarer medium) should travel with more velocity than the wavefront in the core (denser medium). But according to Einstein's theory of relativity, in a single wave front one part should not travel with higher velocity than the other part.

So the part of wave front travelling in cladding medium lost in the form of radiation leads to bending losses.



Faraday's law of electromagnetic induction:

Faraday stated two laws from the observations of Oersted study is called Faraday's law of electromagnetic induction

1) Whenever the magnetic flux linked with an electric circuit (coil) changes, an e.m.f is induced in the circuit (coil). The induced e.m.f exists as long as the change in magnetic flux continues.

2) The magnitude of induced e.m.f is directly proportional to the negative rate of variation of the magnetic flux linked with the circuit.

If Φ_B be the magnetic flux linked with circuit at any instant and e be the induced e.m.f then

$$e = - \frac{d\Phi_B}{dt} \rightarrow (1)$$

Your text here 1

If there are N turns in the coil, then

$$e = - N \frac{d\Phi_B}{dt} \rightarrow (2)$$

The negative sign is in accordance with Lenz's law, this is also called **Neumann's Law**

Lenz's Law:

The Lenz's law is based on the principle of conservation of energy. Thus it helps to explain the direction (polarity) of induced e.m.f or induced current in a coil.

The polarity of the induced e.m.f is always such that it tends to produce a current which opposes the change in magnetic flux that produced it.

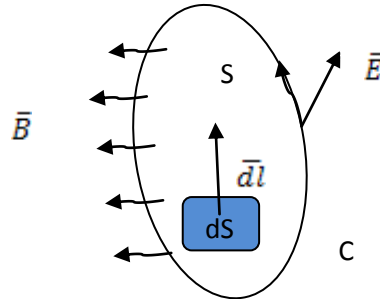
The changing magnetic field and magnetic flux induces an electric current in a coil. This induced current itself creates magnetic field and hence magnetic flux is induced around the coil. Therefore the change in the external magnetic field and flux is always opposed.

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Vector form of Faraday's law(Integral and Differential forms):

Integral form:

Consider that magnetic field is produced by a stationary magnet or current carrying coil. Suppose there is a closed circuit C of any shape which encloses a surface S in the field as shown in fig., Let \vec{B} be the magnetic flux density in the neighborhood of the circuit.



The magnetic flux through a small area “dS” will be \vec{B} . Now the flux through the entire circuit is

$$\Phi_B = \int_S \vec{B} \cdot d\vec{S} \rightarrow (1)$$

When the magnetic flux is changed, an electric field \vec{E} induced around the circuit. The line integral of the electric field gives the induced e.m.f in the closed circuit. Thus

$$e = \oint \vec{E} \cdot d\vec{l} \rightarrow (2)$$

Where E is the electric field at an element of dl of the circuit. Substituting the values of e and Φ_B from equations (2) and (1) we have

$$\oint \vec{E} \cdot d\vec{l} = - \frac{d\Phi_B}{dt} = - \frac{d}{dt} \int_S \vec{B} \cdot d\vec{S} \rightarrow (3)$$

This is the integral form of Faraday's law.

Differential form:

According to equation (3) the line integral of the electric field around any closed circuit is equal to the negative rate of change of magnetic flux through the circuit.

Further by stokes theorem we have

$$\oint \vec{E} \cdot d\vec{l} = \int_S (\nabla \times \vec{E}) \cdot d\vec{S} \rightarrow (4)$$

From equation (3) and (4) we get

$$\int_S (\nabla \times \mathbf{E}) \cdot d\mathbf{S} = - \frac{d}{dt} \int_S \mathbf{B} \cdot d\mathbf{s}$$

$$\text{Hence } \nabla \times \mathbf{E} = - \frac{d\mathbf{B}}{dt}$$

This is the differential form of Faraday's law.

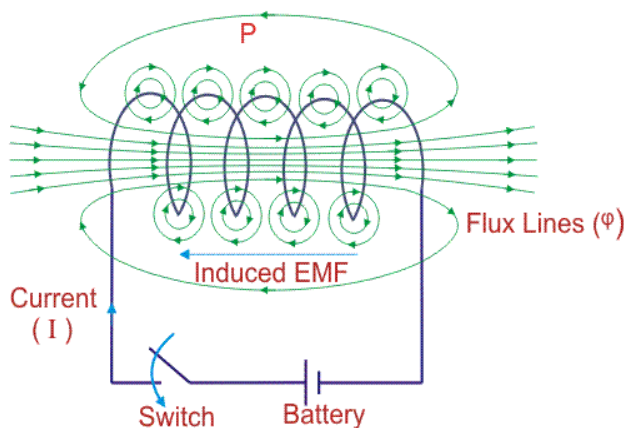
Self Induction:

The phenomenon of self induction was discovered by J. Henry in 1832. When a transient current passes through a coil it produces magnetic field around it, flux due to this field is linked with the coil itself. Due to its own flux change, an e.m.f is induced in the coil and it is called induced e.m.f or back e.m.f.

Definition: Self induction is the property of a coil by virtue of which opposes the growth or decay of the current flowing through it.

(OR)

Self inductance is the phenomenon of inducing e.m.f in a coil due to flow of current which changes with time in the same coil



Consider a coil connected to a battery through key (k). When the key is closed, due to increasing current in the coil, the magnetic field and hence flux linkage with the coil also increases. As a result of this, induced emf is set up in the coil. According to Lenz's law, the direction of induced emf is such that it opposes the growth of current in the coil. This delays the current to acquire the maximum value.

When the key is released, the current in the coil starts decreasing, so the magnetic flux linked with the coil decreases. As a result of this change in the magnetic flux, induced emf is set up in the coil itself. According to Lenz's law, the direction of induced emf is such that it opposes the decay of the current in the coil. This delays the current to acquire minimum (or) zero value.

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NOTE: Self induction is also known as inertia of electricity as it opposes the growth or decay of the current in the circuit

Coefficient of self induction (or) self inductance:

The total magnetic flux Φ_B linked with the coil is proportional to the current I flowing in it i.e.,

$$\Phi_B \propto i$$

Or

$\Phi_B = Li$

$\rightarrow(1)$

Where L is a constant called the coefficient of self induction or self inductance of the coil. When $i=1$, $\Phi_B =L$. Hence the coefficient of self induction is numerically equal to the magnetic flux linked with the coil when unit current flows through it.

The e.m.f. induced in the coil is given by

$$\begin{aligned} e &= - \frac{d\Phi_B}{dt} \\ &= - \frac{d(Li)}{dt} \\ &= - L \frac{di}{dt} \end{aligned} \quad \rightarrow(2)$$

The negative sign indicates that the induced e.m.f is in such a direction as to oppose the change

$$\text{When } \frac{di}{dt} =1 ; e= -L$$

Therefore the coefficient of self inductance is numerically equal to the induced e.m.f in the coil, when the rate of change of current is unity.

Unit: The unit of self inductance is henry which is inductance of a coil in which an e.m.f of 1 volt is set up by the change of current at 1 ampere per second

$$1 \text{ henry} = \frac{1 \text{ volt}}{1 \text{ amp/sec}}$$

Energy stored in electric fields:

Consider a capacitor of capacitance C and carrying a charge q at any instant. Let the potential difference between the plates V , then

$$V = q/c \quad \rightarrow (1)$$

If an additional charge dq is to be given to this capacitor then some work must be done against the potential difference.

So the work done increasing charge by dq is given by

$$dW = V dq = (q/C) dq \quad \rightarrow (2)$$

\therefore Total work to charge a capacitor to a charge q_0

$$W = \int dW = \int_0^{q_0} \frac{q}{c} dq = \frac{q_0^2}{2c} \quad \rightarrow (3)$$

Now the energy stored by a charged capacitor

$$U = W = \frac{1}{2} \frac{q^2}{c} = \frac{1}{2} (CV^2) \quad (\because q = CV) \rightarrow (4)$$

For a parallel plate capacitor of area A and plate separation d , the capacitance C is given by

$$C = \frac{\epsilon_0 A}{d} \quad \text{and} \quad V = Ed$$

$$\therefore \text{Energy stored } U = \frac{1}{2} \times \frac{\epsilon_0 A}{d} \times E^2 d^2$$

$$U = \frac{1}{2} (\epsilon_0 E^2 A d) \text{ joules}$$

Energy stored in magnetic fields:

When the current in a coil is switched on, self induction opposes the growth of current i.e. the current flows against back e.m.f and does work against it.

$$dW = - e i dt$$

$$dW = +L \frac{di}{dt} i dt \quad (\because e = -L \frac{di}{dt})$$

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Hence, total work done in bringing the current from zero to a steady maximum value i_0 is

$$W = L \int_0^{i_0} i \, di$$

$$W = \frac{1}{2} (L i_0^2)$$

Consider a very long solenoid of length 'l' and cross-sectional area 'A'. When current flows in it, a magnetic field is established and the work done is stored as energy in the magnetic field given by

$$U = \frac{1}{2} (L i_0^2)$$

But inductance in coil is given by $L = \mu_0 n^2 A l$

Where n is number of turns in solenoid per meter.

$$\therefore U = \frac{1}{2} (\mu_0 n^2 A l i_0^2) = \frac{1}{2} \times \frac{(\mu_0 n i_0)^2}{\mu_0} \times A l$$

But the magnetic field inside a coil $B = \mu_0 n i_0$

$$U = \frac{B^2}{2\mu_0} \times A l$$

OR

$$U = \frac{\mu_0 H^2}{2} \times A l$$

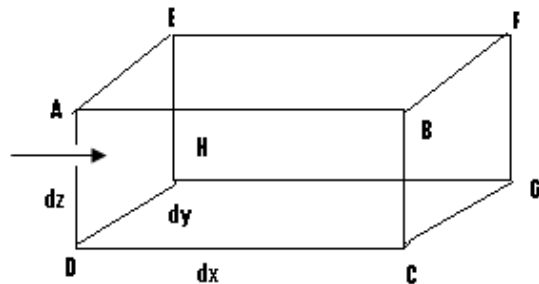
Poynting Vector:

One important characteristic of electromagnetic waves is that they transport energy from one point to another point.

The amount of field energy passing through unit area of the surface perpendicular to the direction of propagation of energy is called as Poynting vector.

The Poynting vector is denoted by \vec{P} given by

$$\vec{P} = \frac{1}{\mu_0} (\vec{E} \times \vec{B}) \quad \text{or} \quad (\vec{E} \times \vec{H})$$



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Proof: Consider an elementary volume in the form of a rectangular parallopiped of sides dx , dy and dz as shown in fig. The volume of parallopiped is “ $dx dy dz$ ”. Suppose the energy is propagated along X-axis, now the area perpendicular to the direction of propagation of energy is $dy dz$. Let the electromagnetic energy in this volume is U . Then the rate of change of energy is $\frac{\partial U}{\partial t}$.

$$\frac{\partial U}{\partial t} = - \oint_s \vec{P} \cdot d\vec{s} \rightarrow (1)$$

Negative sign is used to show that energy is entering in the volume.

Since the energy per unit volume in electric magnetic fields

$$U = (1/2 \epsilon_0 E^2 + 1/2 \mu_0 H^2) \rightarrow (2)$$

The rate of decrease of energy in volume dV is given b

$$-\frac{\partial U}{\partial t} = - \frac{\partial}{\partial t} \int_V (1/2 \epsilon_0 E^2 + 1/2 \mu_0 H^2) dV \rightarrow (3)$$

$$= \int_V \left[\epsilon_0 E \left(\frac{\partial E}{\partial t} \right) + \mu_0 H \left(\frac{\partial H}{\partial t} \right) \right] dV \rightarrow (4)$$

From Maxwell's equation

$$\frac{\partial E}{\partial t} = \frac{\nabla \times H}{\epsilon_0} \text{ and } \frac{\partial H}{\partial t} = - \frac{\nabla \times E}{\mu_0} \rightarrow (5)$$

From equation (5) and (4)

$$-\frac{\partial U}{\partial t} = \int_V [H \cdot (\nabla \times E) - E \cdot (\nabla \times H)] dV \rightarrow (6)$$

$$= \int_V \nabla \cdot (E \times H) dV \rightarrow (7)$$

From Gauss theorem of divergence (7) can be written as

$$= \oint_s (E \times H) \cdot \hat{n} ds \rightarrow (8)$$

Comparing (8) and (1)

$$\vec{P} = \vec{E} \times \vec{H}$$

The vector shows that energy flow takes place in a direction perpendicular to the plane containing E and H

Characteristics of Poynting Vector

- (1) The poynting vector is perpendicular to both electric field vector E and magnetic field vector H
- (2) The quantity $\nabla \cdot P$ represents the net energy flow in electromagnetic field
- (3) In alternating fields

$$P_{avg} = E_{rms} \times H_{rms}$$

Fundamental laws of electromagnetism:

(1) Gauss law in electrostatics:

The electric flux Φ_E through a closed surface is equal to $(1/\epsilon)$ times the net charge q enclosed by the surface. i.e. The surface integral of the normal component of the electric field \vec{E} over any closed surface equals $(1/\epsilon)$ times the net charge with that volume

$$\Phi_E = \oint \vec{E} \cdot d\vec{s} = q/\epsilon$$

(2) Gauss law in magnetostatics:

The magnetic lines of force leaving from the closed surface will be equal to the number of lines of force entering the surface. i.e. The net outward magnetic flux from any closed surface is zero.

$$\Phi_B = \oint \vec{B} \cdot d\vec{s} = 0$$

(3) Faraday's law of electromagnetism

The magnitude of induced e.m.f is directly proportional to the negative rate of variation of the magnetic flux linked with the circuit.

$$\text{e.m.f. (e)} = -\frac{d\Phi_B}{dt}$$

(OR)

$$\oint \vec{E} \cdot d\vec{l} = - \int \frac{\partial \vec{B}}{\partial t} \cdot d\vec{s}$$

(4) Ampere's circuit law(or) Biot-Savart's law:

The steady current carrying conductor generates magnetic field around it

$$\oint \vec{B} \cdot d\vec{l} = \mu_0 i$$

Basic equations of electromagnetism- Maxwell's equation

Maxwell's Law	Integral form	Differential form
First law (based on Gauss law of electrostatics)	$\oint \nabla \cdot \bar{D} \, ds = \int_V \rho \, dV$	$\nabla \cdot \bar{D} = \rho$
Second law (based on Gauss law of magnetostatics)	$\oint \bar{B} \, ds = 0$	$\nabla \cdot \bar{B} = 0$
Third law (based on Faraday's law of electromagnetic induction)	$\oint \bar{E} \cdot d\bar{l} = - \int_s \frac{\partial \bar{B}}{\partial t} \cdot d\bar{s}$	$\nabla \times \bar{E} = - \frac{\partial \bar{B}}{\partial t}$
Fourth law (based on Ampere's circuital law)	$\oint \bar{H} \cdot d\bar{l} = \int_s (J + \frac{\partial D}{\partial t}) \cdot d\bar{s}$	$\nabla \times \bar{B} = \mu_0 (J + \frac{\partial D}{\partial t})$

Displacement current:

According to Maxwell Fourth law, steady current carrying conductor generates magnetic field around it.

$$\oint \bar{B} \cdot d\bar{l} = \mu_0 i \quad \rightarrow (1)$$

Hence

$$\oint \bar{B} \cdot d\bar{l} = \mu_0 \int_s \bar{j} \cdot d\bar{s} \quad \rightarrow (2)$$

$$\text{Using Stoke's theorem } \oint \bar{B} \cdot d\bar{l} = \int_s \nabla \times \bar{B} \cdot d\bar{s} \quad \rightarrow (3)$$

From (2) and (3) $\int_S \nabla \times \vec{B} \cdot d\vec{S} = \mu_0 \int_S \vec{j} \cdot d\vec{S} \rightarrow (4)$

Or $\nabla \times \vec{B} = \mu_0 \vec{j} \rightarrow (5)$

Taking divergence of this equation we get

$$\nabla \cdot (\nabla \times \vec{B}) = \text{div}(\text{curl } \vec{B}) = \text{div} \mu_0 \vec{j} = \mu_0 \text{div } \vec{j} \rightarrow (6)$$

However that divergence of curl of a vector is always zero and hence

$$\text{div } \vec{j} = 0 \rightarrow (7)$$

This shows that the total flux of current out of any closed surface is zero. However from equation of continuity

$$\text{div } \vec{j} + \frac{\partial \rho}{\partial t} = 0 \rightarrow (8)$$

The equation (8) contradicts equation(7)

From Maxwell I equation $\nabla \cdot \vec{D} = \rho$

$$\text{Also } \text{div } \vec{j} + \frac{\partial \rho}{\partial t} = \text{div } \left(\vec{j} + \frac{\partial \vec{D}}{\partial t} \right)$$

Thus $\nabla \cdot \vec{j} = 0$ for steady current

And $\nabla \cdot \left(\vec{j} + \frac{\partial \vec{D}}{\partial t} \right) = 0$ every where

The term $\frac{\partial \vec{D}}{\partial t}$ is called as displacement current density.

Thus a changing electric field is equivalent to a current which flows as long as the electric field is changing and produced the same magnetic effect as an ordinary conduction current. This is known as displacement current.

Electromagnetic wave equation:-

According to Maxwell's electromagnetic equations in a homogeneous medium

- (i) It has infinite resistance to the current and hence its conductivity is zero i.e. $\vec{j} = 0$
- (ii) It has no volume distribution of charge, thus the charge density $\rho = 0$
- (iii) Also $\vec{D} = \epsilon \vec{E}$ and $\vec{B} = \mu \vec{H}$

Hence Maxwell equations

$$\nabla \cdot \vec{E} = 0 \rightarrow (1)$$

$$\nabla \cdot \vec{B} = 0 \rightarrow (2)$$

$$\nabla \times \vec{E} = -\frac{\partial \vec{B}}{\partial t} \rightarrow (3)$$

$$\nabla \times \vec{B} = \mu\epsilon \frac{\partial \vec{E}}{\partial t} \rightarrow (4)$$

The wave equation of propagation of a wave can be obtained by taking curl of eq.(4) as

$$\begin{aligned} \nabla \times \nabla \times \vec{B} &= \nabla \times \mu\epsilon \frac{\partial \vec{E}}{\partial t} = \mu\epsilon (\nabla \times \frac{\partial \vec{E}}{\partial t}) \\ &= \mu\epsilon \frac{\partial}{\partial t} (\nabla \times \vec{E}) \quad (\because \text{from eq.(3)}) \\ &= \mu\epsilon \frac{\partial}{\partial t} \left(-\frac{\partial \vec{B}}{\partial t} \right) \\ &= -\mu\epsilon \frac{\partial^2 \vec{B}}{\partial t^2} \rightarrow (5) \end{aligned}$$

$$\text{Thus } \nabla \times \nabla \times \vec{B} = \nabla(\nabla \cdot \vec{B}) - \nabla^2 \vec{B} = -\nabla^2 \vec{B} \rightarrow (6)$$

From (5) and (6)

$$-\nabla^2 \vec{B} = -\mu\epsilon \frac{\partial^2 \vec{B}}{\partial t^2}$$

$$\text{Or } \nabla^2 \vec{B} = \mu\epsilon \frac{\partial^2 \vec{B}}{\partial t^2} \rightarrow (7)$$

Similarly, from eq.(3) we can show that

$$\nabla^2 \vec{E} = \mu\epsilon \frac{\partial^2 \vec{E}}{\partial t^2} \rightarrow (8)$$

From eqs.(7) and (8) represents the relation between the space and time variation of magnetic field \vec{B} and electric field \vec{E} . Hence the general wave equation is represented by

$$\boxed{\nabla^2 y = \frac{1}{v^2} \frac{\partial^2 y}{\partial t^2}} \rightarrow (9)$$

$$\frac{1}{v^2} = \mu\epsilon \rightarrow (10)$$

$$v = \frac{1}{\sqrt{\mu\epsilon}}$$

Where μ and ϵ are permeability and permittivity of the medium

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When the electromagnetic wave propagating in free space

$$v = \frac{1}{\sqrt{\mu_0 \epsilon_0}}$$

$$\mu_0 = 4\pi \times 10^{-7} \text{ and } \epsilon_0 = 1/(4\pi \times 9 \times 10^9)$$

$$v = \frac{1}{\sqrt{4\pi \times 10^{-7} \times \frac{1}{4\pi \times 9 \times 10^9}}}$$

$$v = \sqrt{9 \times 10^{16}} = 3 \times 10^8 \text{ m/s}$$

Thus the velocity of propagation of variation of **E** and **B** is the same as the velocity of light.

Equation (9) indicate the wave propagation in 3-D free space. These waves involve periodic variations of electric and magnetic field. So they are called electromagnetic waves.

Ultrasonic's: Ultrasonic sound waves have frequencies above the human ear's audible range that is greater than 20kHz and often into megahertz range.

PRODUCTION OF ULTRASONIC WAVES:

The ultrasonic waves cannot be produced by our usual method of loudspeaker fed with alternating current. This is due to the fact that at very high frequencies. There are two methods namely magnetostriction and piezoelectric are used to produce the ultrasonic's

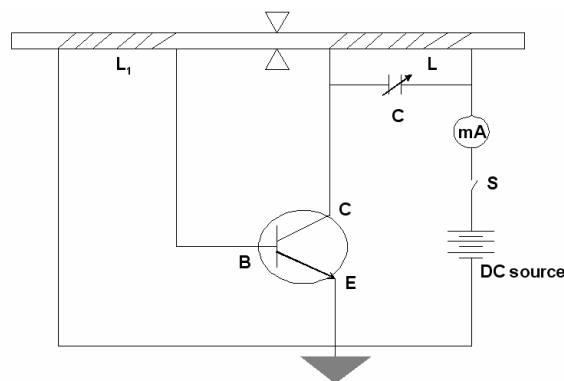
Magnetostriction method

Principle:

When an alternating magnetic field is applied parallel to the length of a ferromagnetic rod such as iron or nickel, a small elongation or contraction occurs in its length. This phenomenon is known as magnetostriction.

Construction:

An experimental arrangement due to production of ultrasonic waves is shown in fig. There is a short nickel rod which is clamped at the centre. A simple tuned oscillator constructed with a NPN transistor with L-C circuit is connected in the collector



Working:

When the supply is switched on, collector current starts rising and oscillations start in the L-C circuit, the changes of current in inductor L is feedback to the base emitter circuit through mutual inductance between L and L_1 . The transistor merely ensures that energy is feedback at

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the correct phase from the source. The frequency of oscillation of L-C circuit is given by $f = 1/2\pi\sqrt{LC}$. By varying C, the frequency can be adjusted to be in tune with the natural frequency of the rod. Under the resonance condition sustained oscillations and hence ultrasonic waves are produced by the rod.

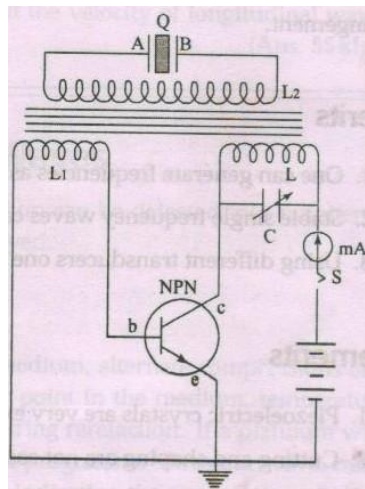
Piezoelectric method:

Principle:

In certain crystals such as quartz is subjected to pressure on one pair of opposite faces (mechanical faces) then in the other pair of opposite faces (electric faces), an opposite electric charges are developed. Similarly when the electric faces are subjected to the alternating voltage then the mechanical faces produces frequencies of order more than 100 kHz

Construction:

The circuit diagram used for generating ultrasonic waves using piezoelectric effect is shown in fig. Q is a thin slice of quartz crystal cut with its opposite faces perpendicular to optic axes. The crystal Q is placed between two metal plates A and B which act as electrodes. The plates are connected to the coil L_2 Coil L is connected to the collector circuit.



Working :

When the supply is switched on, collector current starts raising and oscillations start in the L-C circuit. The frequency of oscillation of the L-C circuit is given by the expression $f = 1/2\pi\sqrt{LC}$. By varying C the frequency of oscillation of the circuit can be adjusted such that electrodes A and B connected to the coil L_2 are induced with alternating e.m.f. Hence the crystal

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Q placed between the electrodes A and B experience oscillating electric force and due to inverse piezoelectric effect, high frequency ultrasonic waves are produced.

APPLICATIONS:

- 1) Ultrasonic inspection is used for quality control and material applications
- 2) Used to measure thickness of metal sections
- 3) Thickness measurements are made on refinery and chemical procession equipments, submarine hulls, aircraft sections and pressure vessels.
- 4) Used to detect internal corrosion.