

BE8255-BASIC ELECTRICAL ,ELECTRONICS AND MEASUREMENT ENGINEERING

UNIT I ELECTRICAL CIRCUIT ANALYSIS

PART A

1. State Ohm's Law.may 2017

The potential difference across any two ends of a conductor is directly proportional to the current flowing between the two ends provided the temperature of the conductor remains constant.

2. State Krichoff's Law ,may2015

KVL states that the algebraic sum of voltages in a closed path is zero.

KCL states that the algebraic sum of currents in a node is zero.

3. Distinguish between a mesh and a loop of a circuit.

A mesh is a loop that does not contain other loops. All meshes are loop, but all loops are not meshes. A loop is any closed path of branches.

4. Write down the formula for a star connected network is converted into a delta network?

$$R_A = \frac{R_1 R_2}{R_1 + R_2 + R_3} \quad R_B = \frac{R_1 R_3}{R_1 + R_2 + R_3} \quad R_C = \frac{R_2 R_3}{R_1 + R_2 + R_3}$$

5. What is reactive power?

If we consider the circuit as purely inductive the output power is reactive power. Its unit is VAR.

6. Define Form factor and Crest factor.

Form factor= RMS value / Average Value

Crest(peak) factor=Maximum Value / RMS value

7. What are the three types of power used in a a.c circuit?

- i) Real power or active power $P=EI \cos \phi$
- ii) Reactive power $Q=EI \sin \phi$
- iii) Apparent power, $S=EI$

8. Define RMS value.

The effective value of an alternating current is that value of steady ,direct current which produces the same heat as that produced by the alternating current when passed which produces the same heat as that produced by the alternating current when passed through the same resistance for the same interval of time.

9. What is meant by eddy current damping?

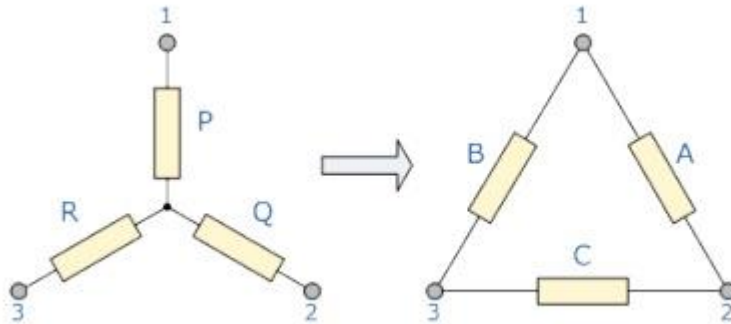
When the conductor moves in a magnetic field an emf is induced in it and if a closed path is provided ,a current flows known as eddy current. This current intersect with the magnetic field to produce an electromagnetic torque , which opposes the deflecting torque.

10. Which type of instrument is called as universal instrument?

The moving iron instrument are known as universal instruments, because these instruments can be used for AC and DC.

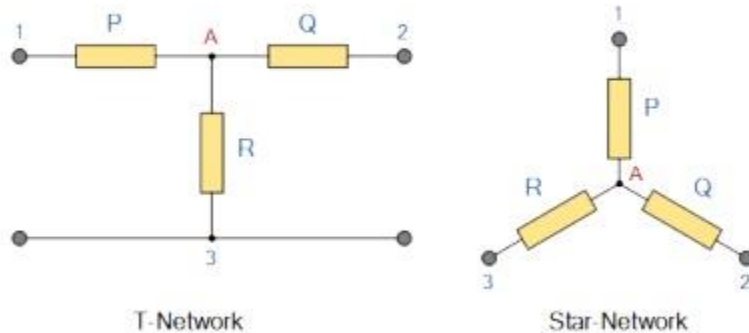
PART B

1.Explain the Star to Delta transformation in detailed manner?

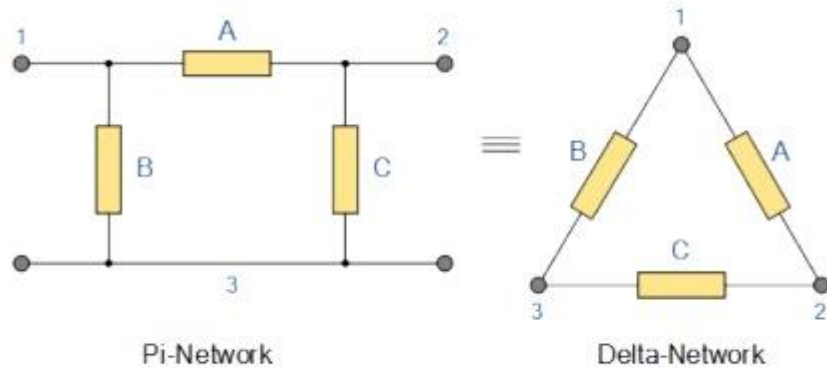


Star Delta Transformations allow us to convert impedances connected together from one type of connection to another. Standard 3-phase circuits or networks take on two major forms with names that represent the way in which the resistances are connected, a **Star** connected network which has the symbol of the letter, Y (wye) and a **Delta** connected network which has the symbol of a triangle, (delta).

If a 3-phase, 3-wire supply or even a 3-phase load is connected in one type of configuration, it can be easily transformed or changed it into an equivalent configuration of the other type by using either the **Star Delta Transformation** or **Delta Star Transformation** process.



Pi-connected and Equivalent Delta Network.



These Circuit Transformations allow us to change the three connected resistances (or impedances) by their equivalents measured between the terminals 1-2, 1-3 or 2-3 for either a star or delta connected circuit.

For example, resistor A is given as:

$A = (PQ + QR + RP) / R$ with respect to terminal 3 and resistor B is given as:

$B = (PQ + QR + RP) / Q$ with respect to terminal 2 and resistor C given as:

$C = (PQ + QR + RP) / R$ with respect to terminal 1.

Star Delta Transformation allows us to convert one type of circuit connection into another type in order for us to easily analyze the circuit and star delta transformation techniques can be used for either resistances or impedances.

2. Three resistances of values 2Ω , 3Ω and 5Ω are connected in series across 20 V, D.C supply Calculate (a) equivalent resistance of the circuit (b) the total current of the circuit (c) the voltage drop across each resistor and (d) the power dissipated in each resistor?

Given data:

$$R_1 = 2\Omega$$

$$R_2 = 3\Omega$$

$$R_3 = 5\Omega$$

$$V = 20V$$

To find:

$$R_T = ?$$

$$I_T = ?$$

$$V_1, V_2, V_3 = ?$$

$$P_1, P_2, P_3 = ?$$

Formula used:

$$R_T = R_1 + R_2 + R_3 \text{ (series connection)}$$

$$I_T = V_T / R_T$$

$$V_1 = R_1 * I_1$$

$$V_2 = R_2 * I_2$$

$$V_3 = R_3 * I_3$$

$$P_1 = V_1 * I_1$$

$$P_2 = V_2 * I_2$$

$$P_3 = V_3 * I_3$$

Solution:

$$R_T = R_1 + R_2 + R_3 = 2 + 3 + 5$$

$$R_T = 10\Omega$$

$$I_T = V_T / R_T = 20 / 10$$

$$I_T = 2 \text{ A}$$

$$\text{In series connection } I_1 = I_2 = I_3 = I_T = 2 \text{ A}$$

$$V_1 = I_1 * R_1 = 2 * 2$$

$$V_1 = 4 \text{ V}$$

$$V_2 = I_2 * R_2 = 2 * 3$$

$$V_2 = 6 \text{ V}$$

$$V_3 = I_3 * R_3 = 2 * 5$$

$$V_3 = 10 \text{ V}$$

$$P_1 = V_1 * I_1 \\ = 4 * 2$$

$$P_1 = 8W$$

$$P_2 = V_2 * I_2 \\ = 6 * 2$$

$$P_2 = 12W$$

$$P_3 = V_3 * I_3 = 10 * 2$$

$$P_3 = 20W$$

Result:

- (a). Equivalent resistance of the circuit $R_T = 10\Omega$
- (b). The total current of the circuit $I_T = 2A$
- (c). Voltage drop across each resistor $V_1 = 4 V$, $V_2 = 6 V$, $V_3 = 10V$
- (d). The power dissipated in each resistor $P_1 = 8W$, $P_2 = 12W$, $P_3 = 20W$

3.Explain the Kirchhoff's Law in detailed manner?

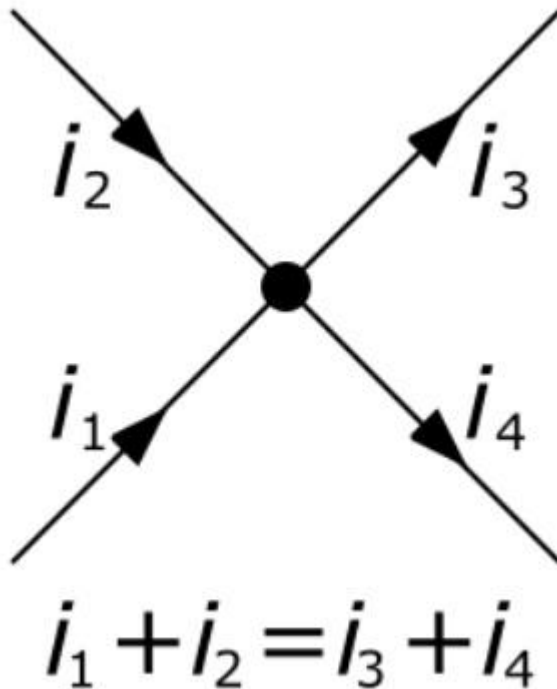
Kirchhoff's Current Law:

First law (Current law or Point law): Statement:

The sum of the currents flowing towards any junction in an electric circuit equal to the sum of currents flowing away from the junction.

Kirchhoff's Current law can be stated in words as the sum of all currents flowing into a node is zero. Or conversely, the sum of all currents leaving a node must be zero. As the image below demonstrates, the sum of currents I_b , I_c , and I_d , must equal the total current in I_a . Current flows through wires much like water flows through pipes. If you have a definite amount of water entering a closed pipe system, the amount of water that enters the system must equal the amount of water that exists the system. The number of branching pipes does not change the net volume of water (or current in our case) in the system.

Kirchhoff's law:



Kirchhoff's Voltage Law:

Second law (voltage law or Mesh law): Statement:

In any closed circuit or mesh, the algebraic sum of all the electromotive forces and the voltage drops is equal to zero.

Kirchhoff's voltage law can be stated in words as the sum of all voltage drops and rises in a closed loop equals zero. As the image below demonstrates, loop 1 and loop 2 are both closed loops within the circuit. The sum of all voltage drops and rises around loop 1 equals zero, and the sum of all voltage drops and rises in loop 2 must also equal zero. A closed loop can be defined as any path in which the originating point in the loop is also the ending point for the loop. No matter how the loop is defined or drawn, the sum of the voltages in the loop must be zero.

4.Explain the superposition theorem in detailed manner?

The superposition theorem is unquestionably one of the most powerful in this field. It has such widespread application that people often apply it without recognizing that their maneuvers are valid only because of this theorem.

In general, the theorem can be used to do the following:

- Analyze networks such as introduced in the last chapter that have two or more sources that are not in series or parallel.
- Reveal the effect of each source on a particular quantity of interest.
- For sources of different types (such as dc and ac, which affect the parameters of the network in a different manner) and apply a separate analysis for each type, with the total result simply the algebraic sum of the results.
- Become familiar with the superposition theorem and its unique ability to separate the impact of each source on the quantity of interest.
- Be able to apply Thévenin's theorem to reduce any two-terminal, series-parallel network with any number of sources to a single voltage source and series resistor.
- Become familiar with Norton's theorem and how it can be used to reduce any two-terminal, seriesparallel network with any number of sources to a single current source and a parallel resistor.

The superposition theorem states the following: The current through, or voltage across, any element of a network is equal to the algebraic sum of the currents or voltages produced independently by each source. In other words, this theorem allows us to find a solution for a current or voltage using only one source at a time. Once we have the solution for each source, we can combine the results to obtain the total solution. The term algebraic appears in the above theorem statement because the currents resulting from the sources of the network can have different directions, just as the resulting voltages can have opposite polarities. If we are to consider the effects of each source, the other sources obviously must be removed. Setting a voltage source to zero volts is like placing a short circuit across its terminals. Therefore, when removing a voltage source from a network schematic, replace it with a direct connection (short circuit) of zero ohms. Any internal resistance associated with the source must remain in the network. Setting a current source to zero amperes is like replacing it with an open circuit. Therefore, when removing a current source from a network schematic, replace it by an open

circuit of infinite ohms. Any internal resistance associated with the source must remain in the network.

5. Explain the thevenin's theorem in detailed manner?

Thévenin's theorem The next theorem to be introduced, Thévenin's theorem, is probably one of the most interesting in that it permits the reduction of complex networks to a simpler form for analysis and design.

In general, the theorem can be used to do the following:

- Analyze networks with sources that are not in series or parallel.
- Reduce the number of components required to establish the same characteristics at the output terminals.

Thévenin's Theorem Procedure Preliminary:

1. Remove that portion of the network where the Thévenin equivalent circuit is found. This requires that the load resistor R_L be temporarily removed from the network.

2. Mark the terminals of the remaining two-terminal network. (The importance of this step will become obvious as we progress through some complex networks.) R_{Th}

3. Calculate R_{Th} by first setting all sources to zero (voltage sources are replaced by short circuits and current sources by open circuits) and then finding the resultant resistance between the two marked terminals. (If the internal resistance of the voltage and/or current sources is included in the original network, it must remain when the sources are set to zero.) E_{Th}

4. Calculate E_{Th} by first returning all sources to their original position and finding the open-circuit voltage between the marked terminals. (This step is invariably the one that causes most confusion and errors. In all cases, keep in mind that it is the open-circuit

5. Draw the Thévenin equivalent circuit with the portion of the circuit previously removed replaced between the terminals of the equivalent circuit.

6. Explain the norton's theorem in detailed manner?

The theorem states the following: Any two-terminal linear bilateral dc network can be replaced by an equivalent circuit consisting of a current source and a parallel resistor

Norton's Theorem Procedure Preliminary:

1. Remove that portion of the network across which the Norton equivalent circuit is found.

2. Mark the terminals of the remaining two-terminal network. R_N :

3. Calculate R_N by first setting all sources to zero (voltage sources are replaced with short circuits and current sources with open circuits) and then finding the resultant resistance between the two marked terminals. (If the internal resistance of the voltage and/or current sources is included in the original network, it must remain when the sources are set to zero.) Since $R_N = R_{Th}$, the procedure and value obtained using the approach described for Thévenin's theorem will determine the proper value of R_N . I_N :

4. Calculate I_N by first returning all sources to their original position and then finding the short-circuit current between the marked terminals. It is the same current that would be measured by an ammeter placed between the marked terminals. Conclusion:

5. Draw the Norton equivalent circuit with the portion of the circuit previously removed replaced between the terminals of the equivalent circuit.

7. Explain the maximum power transfer theorem in detailed manner?

maximum power transfer theorem When designing a circuit, it is often important to be able to answer one of the following questions:

What load should be applied to a system to ensure that the load is receiving maximum power from the system

Conversely: For a particular load, what conditions should be imposed on the source to ensure that it will deliver the maximum power available

Even if a load cannot be set at the value that would result in maximum power transfer, it is often helpful to have some idea of the value that will draw maximum power so that you can compare it to the load at hand.

For instance, if a design calls for a load of $100\ \Omega$, to ensure that the load receives maximum power, using a resistor of $1\ \Omega$ or $1\ \text{k}\Omega$ results in a power transfer that is much less than the maximum possible. However, using a load of $82\ \Omega$ or $120\ \Omega$ probably results in a fairly good level of power transfer.

Fortunately, the process of finding the load that will receive maximum power from a particular system is quite straightforward due to the maximum power transfer theorem, which states the following: A load will receive maximum power from a network when its resistance is exactly equal to the Thévenin resistance of the network, for the Thévenin equivalent circuit in which the load is set equal to the Thévenin resistance, the load will receive maximum power from the network. with $R_L = R_{Th}$,

we can determine the maximum power delivered to the load by first finding the current: $I_L = \frac{E_{Th}}{R_{Th} + R_L} = \frac{E_{Th}}{R_{Th} + R_{Th}} = \frac{E_{Th}}{2R_{Th}}$ Then we substitute into the power equation: $P_L = I^2 R_L = \left(\frac{E_{Th}}{2R_{Th}}\right)^2 R_L = \frac{E_{Th}^2 R_L}{4R_{Th}^2}$ and $P_{Lmax} = \frac{E_{Th}^2}{4R_{Th}}$ To demonstrate that maximum power is indeed transferred to the load under the conditions defined above, consider the Thévenin equivalent circuit in Fig. 9.85. Before getting into detail, however, if you were to guess what value of R_L would result in maximum power transfer to R_L , you might think that the smaller the value of R_L , the better it is because the current reaches a maximum when it is squared in the power equation. The problem is, however, that in the equation $P_L = I^2 R_L$, the load resistance is a multiplier. As it gets smaller, it forms a smaller product. Then again, you might suggest larger values of R_L because the output voltage increases, and power is determined by $P_L = V^2 / R_L$. This time, however, the load resistance is in the denominator of the equation and causes the resulting power to decrease. A balance must obviously be made between the load resistance and the resulting current or voltage. The following discussion shows that maximum power transfer occurs when the load voltage and current are one-half their maximum possible values. For the circuit in Fig. 9.85, the current through the load is determined by $I_L = \frac{E_{Th}}{R_{Th} + R_L} = \frac{60\ \text{V}}{9\ \Omega + R_L}$ The voltage is determined by $V_L = \frac{R_L E_{Th}}{R_L + R_{Th}} = \frac{R_L(60\ \text{V})}{R_L + 9\ \Omega}$ and the power by $P_L = I^2 R_L = \left(\frac{60\ \text{V}}{9\ \Omega + R_L}\right)^2 R_L = \frac{3600R_L}{(9\ \Omega + R_L)^2}$ If we tabulate the three quantities versus a range of values for R_L from $0.1\ \Omega$ to $30\ \Omega$, If the load applied is less than the Thévenin resistance, the power to the load will drop off rapidly as it gets smaller.

However, if the applied load is greater than the Thévenin resistance, the power to the load will not drop off as rapidly as it increases.

8. Explain the parallel circuit analysis in detailed manner?

Resistors in Parallel Consider a circuit with 3 resistors in parallel (such as the circuit below, if $N = 3$). Since the voltages across all the parallel elements in a circuit are the same ($E = V_1 = V_2 = V_3$),

This result can be generalized to provide the total resistance of any number of resistors in parallel:

Special Case: Two Resistors in Parallel For only two resistors connected in parallel, the equivalent resistance may be found by the product of the two values divided by the sum:

Current through resistors in parallel. The total current I is shared by the resistors in inverse proportion to their resistances. Stated another way: "More current follows the path of least resistance."

Extreme cases for current division: **Current Divider Rule** The Current Divider Rule (CDR) allows us to determine how the current flowing into a node is split between the various parallel resistors.

Compare the formulas for the voltage divider rule and the current divider rule. **Special Case: Two resistors in parallel.** For only two resistors in parallel: open circuit short circuit

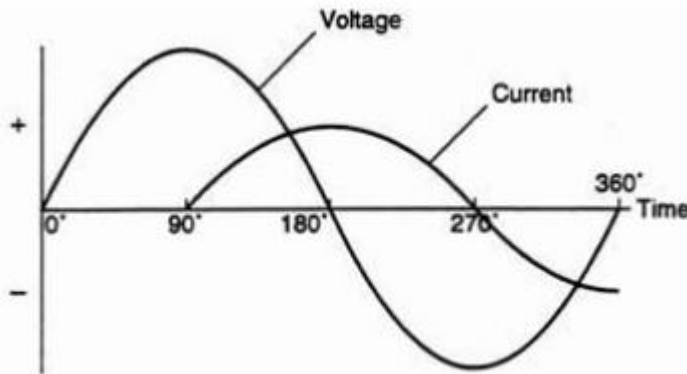
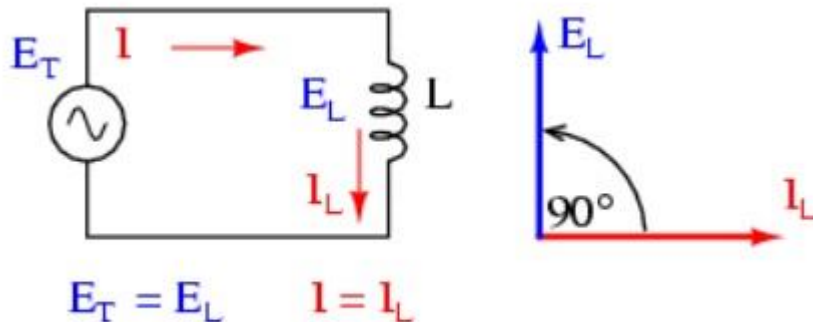
Special Case: If current enters a parallel network with a number of equal resistors, the current will split equally between the resistors. Note:

In a parallel network, the smallest value resistor will have the largest current. **Analysis of Parallel Circuits** To analyze parallel circuits we should use the following guidelines:

1. Voltage across all branches is the same as the source voltage
2. Determine current through each branch using Ohm's Law
3. Find the total current using Kirchhoff's Current Law

9. Explain the inductive networks in detailed manner?

Pure Inductive circuits:



This simple circuit above consists of a pure inductance of L Henrys (H), connected across a sinusoidal voltage given by the expression: $V(t) = V_{\max} \sin \omega t$. When the switch is closed this sinusoidal voltage will cause a current to flow and rise from zero to its maximum value. This rise or change in the current will induce a magnetic field within the coil which in turn will oppose or restrict this change in the current.

But before the current has had time to reach its maximum value as it would in a DC circuit, the voltage changes polarity causing the current to change direction. This change in the other direction once again being delayed by the self-induced back emf in the coil, and in a circuit containing a pure inductance only, the current is delayed by 90°.

The applied voltage reaches its maximum positive value a quarter ($1/4f$) of a cycle earlier than the current reaches its maximum positive value, in other words, a voltage applied to a purely inductive circuit “LEADS” the current by a quarter of a cycle or 90° as shown below.

The instantaneous voltage across the resistor, V_R is equal to the supply voltage, V_t and is given as:

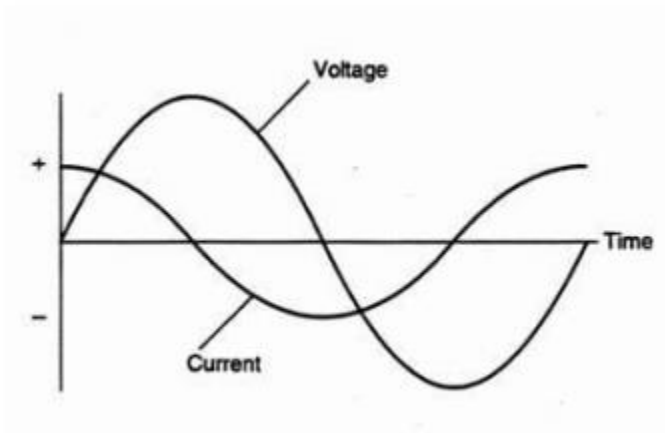
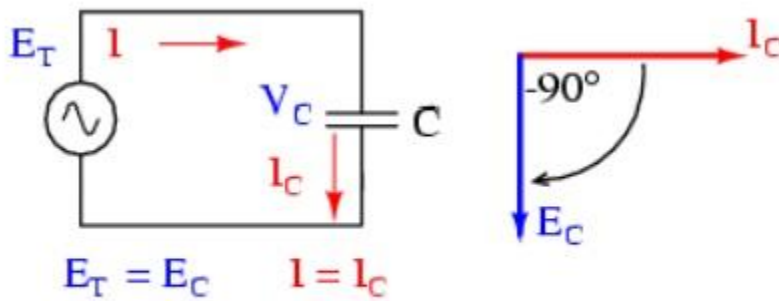
$$V_L = V_{\max} \sin(\omega t + 90)$$

$$I_L = V / X_L$$

$$X_L = 2\pi fL$$

10. Explain the capacitive networks in detailed manner?

Pure Capacitive circuits:



When the switch is closed in the circuit above, a high current will start to flow into the capacitor as there is no charge on the plates at $t = 0$. The sinusoidal supply voltage, V is increasing in a positive direction at its maximum rate as it crosses the zero reference axis at an instant in time given as 0° . Since the rate of change of the potential difference across the plates is now at its maximum value,

the flow of current into the capacitor will also be at its maximum rate as the maximum amount of electrons are moving from one plate to the other.

As the sinusoidal supply voltage reaches its 90° point on the waveform it begins to slow down and for a very brief instant in time the potential difference across the plates is neither increasing nor decreasing therefore the current decreases to zero as there is no rate of voltage change. At this 90° point the potential difference across the capacitor is at its maximum (V_{\max}), no current flows into the capacitor as the capacitor is now fully charged and its plates saturated with electrons.

At the end of this instant in time the supply voltage begins to decrease in a negative direction down towards the zero reference line at 180°. Although the supply voltage is still positive in nature the capacitor starts to discharge some of its excess electrons on its plates in an effort to maintain a constant voltage. These results in the capacitor current flowing in the opposite or negative direction.

When the supply voltage waveform crosses the zero reference axis point at instant 180°, the rate of change or slope of the sinusoidal supply voltage is at its maximum but in a negative direction, consequently the current flowing into the capacitor is also at its maximum rate at that instant. Also at this 180° point the potential difference across the plates is zero as the amount of charge is equally distributed between the two plates.

Then during this first half cycle 0° to 180°, the applied voltage reaches its maximum positive value a quarter ($1/4f$) of a cycle after the current reaches its maximum positive value, in other words, a voltage applied to a purely capacitive circuit “LAGS” the current by a quarter of a cycle or 90°

$$I_C = I_{\max} \sin (\omega t + 90)$$

$$I_L = V / X_C$$

$$X_C = 1 / 2\pi fC$$

UNIT II ELECTRICAL MACHINES

PART A

1. What is an electric generator?

An electrical machine, which converts mechanical energy into electrical Energy, is called as electric generator.

2. What is an electric motor?

An electrical machine, which converts electrical energy into mechanical Energy, is called as electric motor.

3. What is the function yoke?

It serves the purpose of outermost cover of the dc machine. So that the insulating material get protected from harmful atmospheric elements like moisture, dust and various gases like SO₂, acidic fumes etc. It provides mechanical support to the poles.

4. Define voltage transformation ratio?dec2015

The ratio of secondary induced emf to primary induced emf is called as voltage regulation ratio devoted by K.

$$\frac{E_2}{E_1} = \frac{N_2}{N_1} = K$$

5. Write down the emf equation for d.c generator.

$$E = (\phi NZ / 60)(P/A) \quad \text{V}$$

Where

P = number of poles

Z = Total number of conductors

A = number of parallel paths

Φ = flux per pole

N = speed in rpm

6. what is the principle of DC motor?

Whenever a current carrying conductor placed in a magnetic field, it experiences a mechanical force.

7. List the different types of DC motor.

- DC series motor
- DC Shunt motor
- DC Compound motor
 - Long shunt compound motor
 - Short shunt compound motor

8. Give the torque equation of a DC motor.

$$T_a = 0.159 f I_a \frac{PZ}{A} \text{ N-m}$$

I_a - Armature current

P - Number of poles

Z - Total number of conductors

A - Number of parallel paths

9. List out the characteristics of DC motor.

- i. Torque-Armature current characteristics (T VS I_a)
- ii. Speed-Armature current characteristics (N VS I_a)

10. What are all the applications of DC motor?

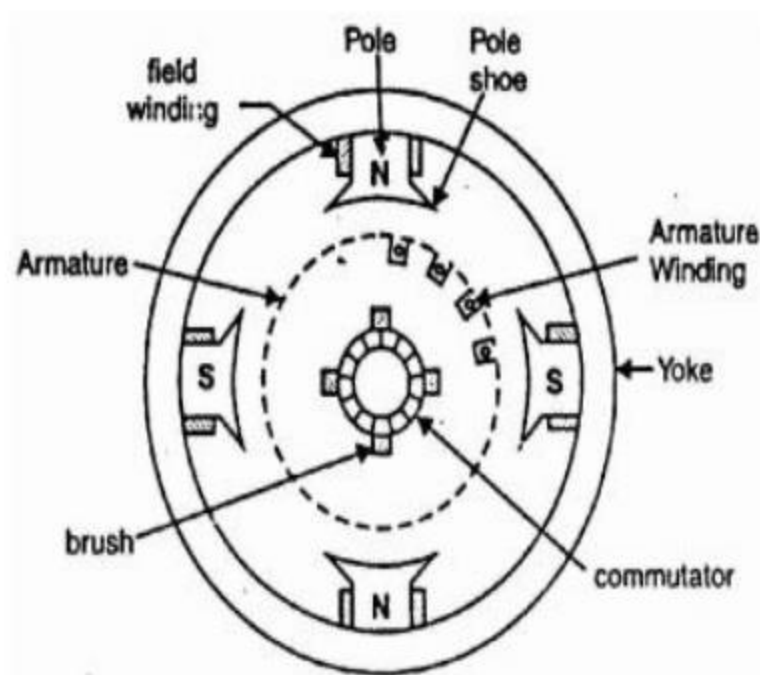
Blowers and fans
Centrifugal and reciprocating pumps
Lathe machines
Machine tools
Milling machines
Drilling machines

PART B

1. write the notes on construction of d.c. machines in clear manner?

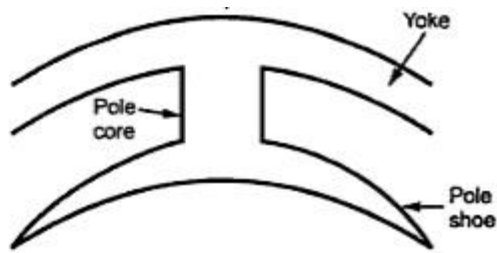
Frame

Frame is the stationary part of a machine on which the main poles and commutator poles are bolted and it forms the supporting structure by connecting the frame to the bed plate.



poles:

Inter-poles are small additional poles located in between the main poles.



These can be solid, or laminated just as the main poles. These are also fastened to the yoke by bolts.

Armature

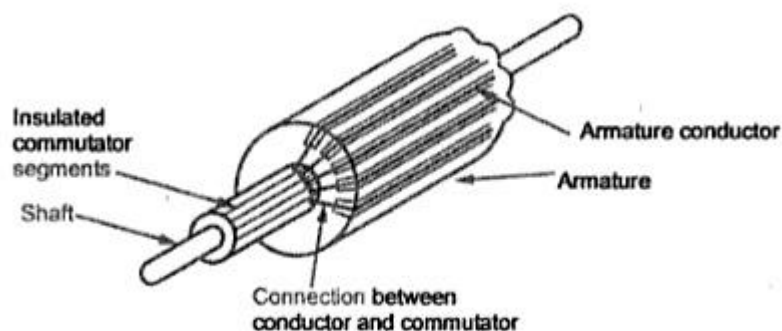
The armature is where the moving conductors are located. The armature is constructed by stacking laminated sheets of silicon steel. Thickness of these lamination is kept low to reduce eddy current losses.

Field windings:

Hence the armature windings are in general pre-formed, taped and lowered into the open slots on the armature. In the case of small machines, they can be hand wound. The coils are prevented from flying out due to the centrifugal forces by means of bands of steel wire on the surface of the rotor in small groves cut into it. In the case of large machines slot wedges are additionally used to restrain the coils from flying away.

Commutator:

Commutator is the key element which made the d.c. machine of the present day possible. It consists of copper segments tightly fastened together with mica/micanite insulating separators



on an insulated base. The whole commutator forms a rigid and solid assembly of insulated copper strips and can rotate at high speeds.

Brush and brush holders:

The brushes are kept pressed on the commutator with the help of springs. This is to ensure proper contact between the brushes and the commutator even under high speeds of operation. Jumping of brushes must be avoided to ensure arc free current collection and to keep the brush contact drop low. Other mechanical parts End covers, fan and shaft bearings form other important mechanical parts.

End Shields or Bearings

If the armature diameter does not exceed 35 to 45 cm then in addition to poles end shields or frame head with bearing are attached to the frame. If the armature diameter is greater than 1m **pedestal type bearings** are mounted on the machine bed plate outside the frame.

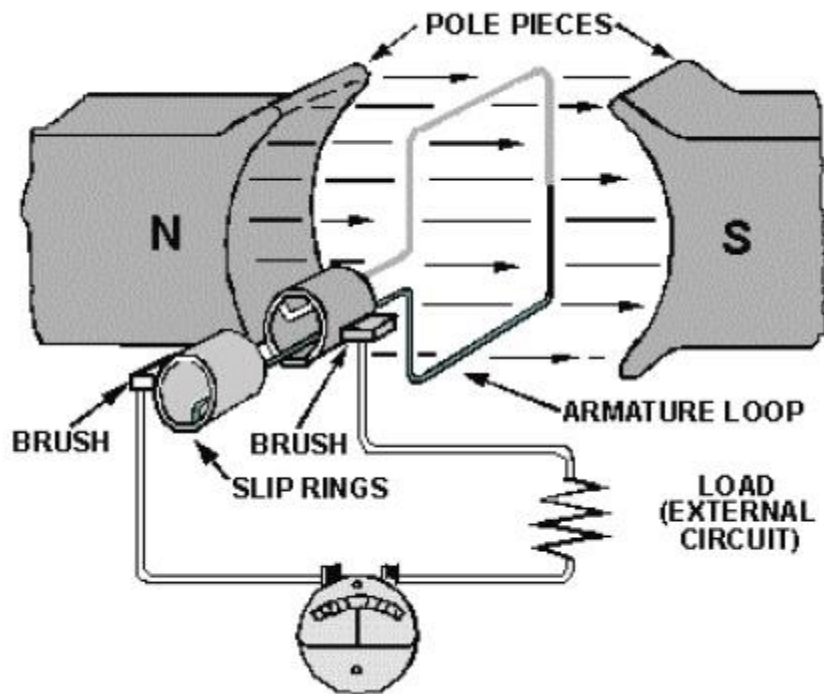
2. write the notes on principle of d.c. machines in clear manner?

PRINCIPLE OF OPERATION

DC generator converts mechanical energy into electrical energy. when a conductor move in a magnetic field in such a way conductors cuts across a magnetic flux of lines and emf produces in a generator and it is defined by faradays law of electromagnetic induction emf causes current to flow if the conductor circuit is closed. The pole pieces (marked N and S) provide the magnetic field. The pole pieces are shaped and positioned as shown to concentrate the magnetic field as close as possible to the wire loop. The loop of wire that rotates through the field is called the **ARMATURE**. The ends of the armature loop are connected to rings called **SLIP RINGS**. They rotate with the armature.

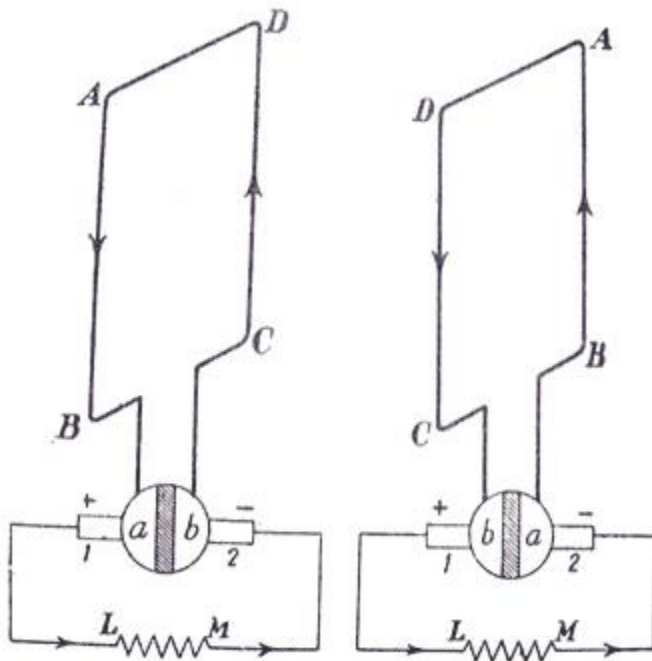
The brushes, usually made of carbon, with wires attached to them, ride against the rings. The generated voltage appears across these brushes. The elementary generator produces a voltage in the following manner (fig. 1-3). The armature loop is rotated in a clockwise direction. The initial or starting point is shown at position A. (This will be considered the zero-degree position.) At 0° the armature loop is perpendicular to the magnetic field. The black and white conductors of the loop are moving parallel to the field. The instant the conductors are moving parallel to the magnetic field, they do not cut any lines of flux. Therefore, no emf is induced in the conductors, and the meter at position A indicates zero. This position is called

the NEUTRAL PLANE. As the armature loop rotates from position A (0°) to position B (90°), the conductors cut through more and more lines of flux, at a continually increasing angle. At 90° they are cutting through a maximum number of lines of flux and at maximum angle. The result is that between 0° and 90° , the induced emf in the conductors builds up from zero to a maximum value. Observe that from 0° to 90° , the black conductor cuts DOWN through the field. At the same time the white conductor cuts UP through the field.



The induced emfs in the conductors are series-adding. This means the resultant voltage across the brushes (the terminal voltage) is the sum of the two induced voltages. The meter at position B reads maximum value. As the armature loop continues rotating from 90° (position B) to 180° (position C), the conductors which were cutting through a maximum number of lines of flux at position B now cut through fewer lines. They are again moving parallel to the magnetic field at position C. They no longer cut through any lines of flux. As the armature rotates from 90° to 180° , the induced voltage will decrease to zero in the same manner that it increased during the rotation from 0° to 90° . The meter again reads zero. From 0° to 180° the conductors of the armature loop have been moving in the same direction through the magnetic field. Therefore, the polarity of the induced

voltage has remained the same. This is shown by points A through C on the graph.



3. Explain the E.M.F equations of d.c. machines in clear manner?

E.M.F EQUATION

Let

Φ = flux/pole in weber

Z = total number of armature conductors = No. of slots x No. of conductors/slot

P = No. of generator poles

A = No. of parallel paths in armature

N = armature rotation in revolutions per minute (r.p.m) E = e.m.f induced in any parallel path in armature

Generated e.m.f E_g = e.m.f generated in any one of the parallel paths i.e E .
Average e.m.f generated /conductor = $d\Phi/dt$ volt ($n=1$)

Now, flux cut/conductor in one revolution $d\Phi = \Phi P$ Wb No. of revolutions/second = $N/60$

Time for one revolution, $dt = 60/N$ second

Hence, according to Faraday's Laws of Electroagnetic Induction,

E.M.F generated/conductor is

$$\frac{d\phi}{dt} = \frac{\phi PN}{60}$$

For a simplex wave-wound generator

No.of parallel paths = 2

No.of conductors (in series) in one path = $Z/2$

E.M.F. generated/path is

$$\frac{\phi PN}{60} \times \frac{Z}{2} = \frac{\phi ZPN}{120} \text{ volt}$$

For a simplex lap-wound generator

No.of parallel paths = P

No.of conductors (in series) in one path = Z/P

E.M.F.generated/path

$$\frac{\phi PN}{60} \times \frac{Z}{P} = \frac{\phi ZN}{60} \text{ volt}$$

In general generated e.m.f

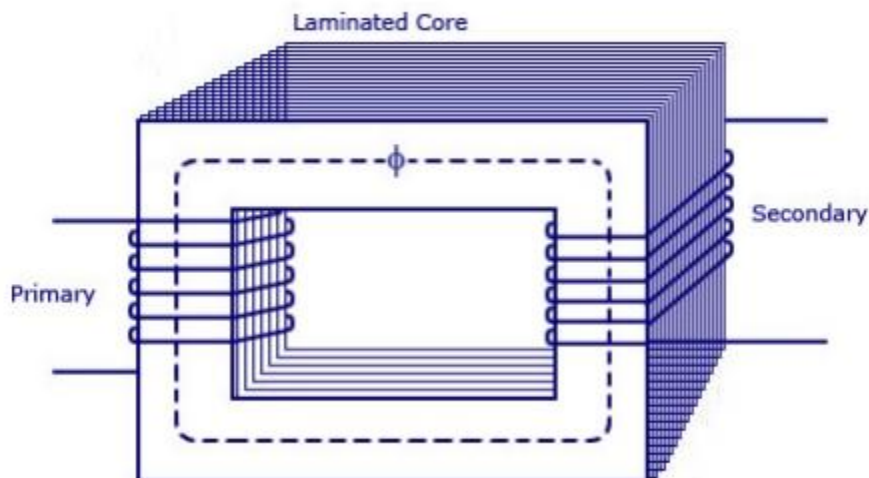
$$E_g = \frac{\phi ZN}{60} \times \left(\frac{P}{A}\right) \text{ volt}$$

where $A = 2$ for simplex wave-winding $A = P$ for simplex lap-winding

4. Explain the working principle of transformer in clear manner?

BASIC WORKING PRINCIPLE OF TRANSFORMER

A transformer can be defined as a static device which helps in the transformation of electric power in one circuit to electric power of the same frequency in another circuit. The voltage can be raised or lowered in a circuit, but with a proportional increase or decrease in the current ratings. The main principle of operation of a transformer is mutual inductance between two circuits which is linked by a common magnetic flux. A basic transformer consists of two coils that are electrically separate and inductive, but are magnetically linked through a path of reluctance. The working principle of the transformer can be understood from the figure below.



The core laminations are joined in the form of strips in between the strips there are some narrow gaps right through the cross-section of the core. These staggered joints are said to be 'imbricated'. Both the coils have high mutual inductance. A mutual electro-motive force is induced in the transformer from the alternating flux that is set up in the laminated core, due to the coil that is connected to a source of alternating voltage. Most of the alternating flux developed by this coil is linked with the other coil and thus produces the mutual induced electro-motive force. The so produced electro-motive force can be explained with the help of Faraday's laws of Electromagnetic Induction as

$$e = M \cdot dI/dt$$

If the second coil circuit is closed, a current flows in it and thus electrical energy is transferred magnetically from the first to the second coil.

The alternating current supply is given to the first coil and hence it can be called as the primary winding. The energy is drawn out from the second coil and thus can be called as the secondary winding.

In short, a transformer carries the operations shown below:

Transfer of electric power from one circuit to another.

Transfer of electric power without any change in frequency.

Transfer with the principle of electromagnetic induction.

The two electrical circuits are linked by mutual induction

5. Explain the construction of transformer in clear manner?

TRANSFORMER CONSTRUCTION

Two coils of wire (called windings) are wound on some type of core material. In some cases the coils of wire are wound on a cylindrical or rectangular cardboard form. In effect, the core material is air and the transformer is called an AIR-CORE TRANSFORMER. Transformers used at low frequencies, such as 60 hertz and 400 hertz, require a core of low-reluctance magnetic material, usually iron. This type of transformer is called an IRON-CORE TRANSFORMER. Most power transformers are of the iron-core type.

The principle parts of a transformer and their functions are:

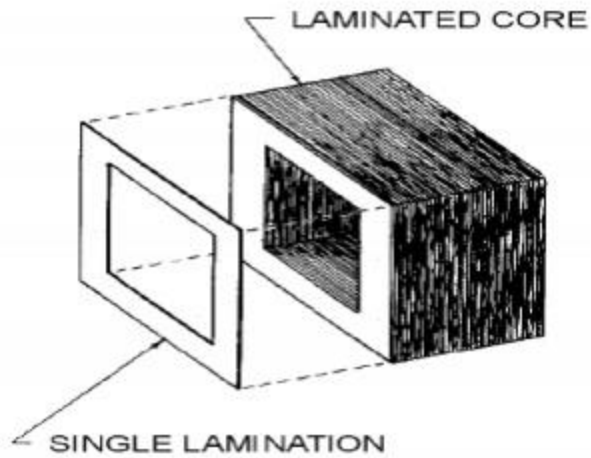
The CORE, which provides a path for the magnetic lines of flux.

The PRIMARY WINDING, which receives energy from the ac source.

The SECONDARY WINDING, which receives energy from the primary winding and delivers it to the load.

The ENCLOSURE, which protects the above components from dirt, moisture, and mechanical damage.

(i) CORE



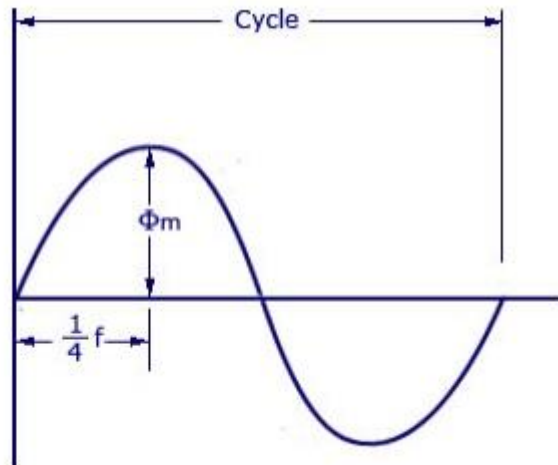
There are two main shapes of cores used in laminated-steel-core transformers. One is the **HOLLOWCORE**, so named because the core is shaped with a hollow square through the center. This shape of core. Notice that the core is made up of many laminations of steel it shows how the transformer windings are wrapped around both sides of the core.

(ii) WINDINGS

the transformer consists of two coils called **WINDINGS** which are wrapped around a core. The transformer operates when a source of ac voltage is connected to one of the windings and a load device is connected to the other. The winding that is connected to the source is called the **PRIMARY WINDING**. The winding that is connected to the load is called the **SECONDARY WINDING**. The primary is wound in layers directly on a rectangular cardboard form.

6. Write the E.M.F equation of transformer in clear manner?

The applied voltage V_1 applied to the primary of a transformer, with secondary open-circuited, be sinusoidal (or sine wave). Then the current I_1 , due to applied voltage V_1 , will also be a sine wave. The mmf $N_1 I_1$ and core flux Φ will follow the variations of I_1 closely. That is the flux is in time phase with the current I_1 and varies sinusoidally.



N_A = Number of turns in primary

N_B = Number of turns in secondary

Φ_{\max} = Maximum flux in the core in webers = $B_{\max} \times A$ f = Frequency of alternating current input in hertz (Hz)

the core flux increases from its zero value to maximum value Φ_{\max} in one quarter of the cycle, that is in $\frac{1}{4}$ frequency second.

Therefore, average rate of change of flux = $\Phi_{\max} / \frac{1}{4} f = 4f \Phi_{\max}$ Wb/s

Now, rate of change of flux per turn means induced electro motive force in volts. Therefore,

average electro-motive force induced/turn = $4f \Phi_{\max}$ volt

If flux Φ varies sinusoidally, then r.m.s value of induced e.m.f is obtained by multiplying the average value with form factor.

$$E_A = 4.44f N_A \Phi_{\max} = 4.44f N_A B_m A$$

Similarly, r.m.s value of induced e.m.f in secondary is

$$E_B = 4.44f N_B \Phi_{\max} = 4.44f N_B B_m A$$

In an ideal transformer on no load, $V_A = E_A$ and $V_B = E_B$, where V_B is the terminal voltage

Voltage Transformation Ratio.

The ratio of secondary voltage to primary voltage is known as the voltage transformation ratio and is designated by letter K . i.e.

$$\text{Voltage transformation ratio, } K = V_2/V_1 = E_2/E_1 = N_2/N_1$$

Current Ratio

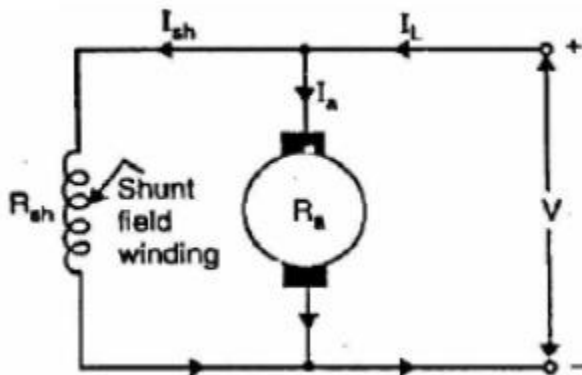
The ratio of secondary current to primary current is known as current ratio and is reciprocal of voltage transformation ratio in an ideal transformer.

7. Explain the types of DC motors in clear manner?

DC MOTOR TYPES

1. Shunt Wound
2. Series Wound
3. Compound wound

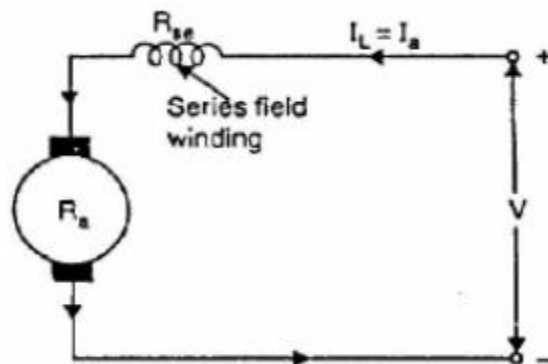
1. Shunt Motor



Shunt field windings are designed to produce the necessary m.m.f. by means of a relatively large number of turns of wire having high resistance. Therefore, shunt field current is relatively small compared with the armature current

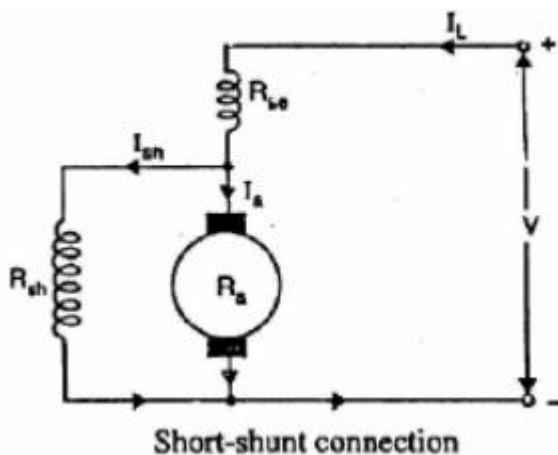
2. Series Motor

The current passing through a series field winding is the same as the armature current, series field windings must be designed with much fewer turns than shunt field windings for the same mmf. Therefore, a series field winding has a relatively small number of turns of thick wire and, therefore, will possess a low resistance.



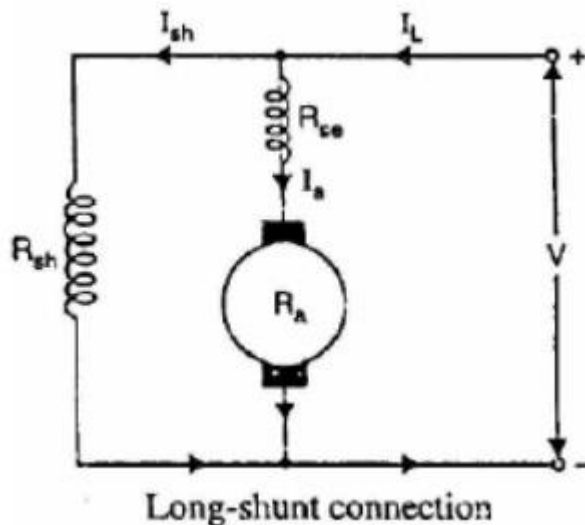
3. Compound Wound Motor

- 1) Short-shunt connection
- 2) Long shunt connection



When the shunt field winding is directly connected across the armature terminals it is called short-shunt connection.

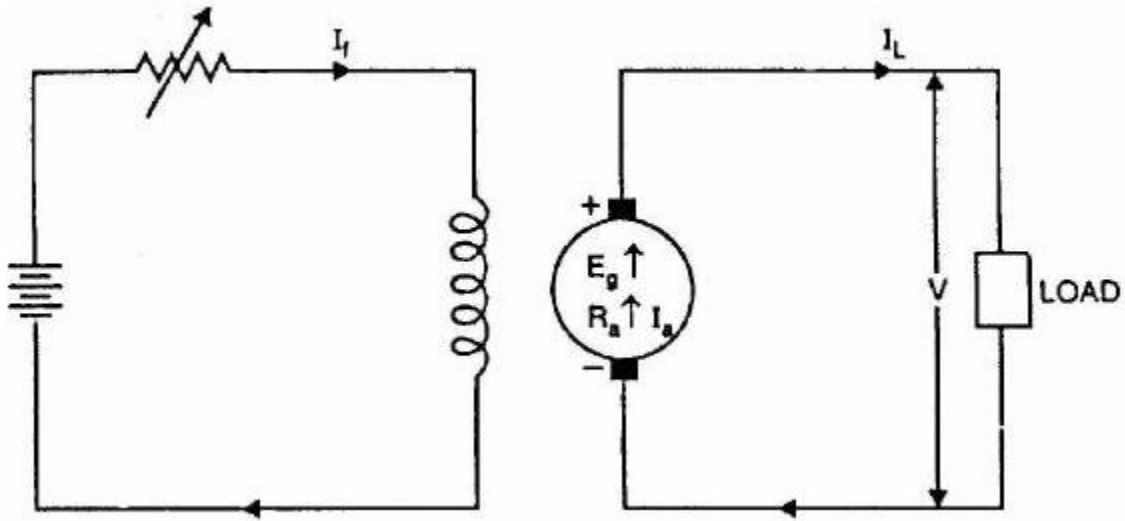
When the shunt winding is so connected that it shunts the series combination of armature and series field it is called long-shunt connection.



8. Explain the types of DC generators in clear manner?

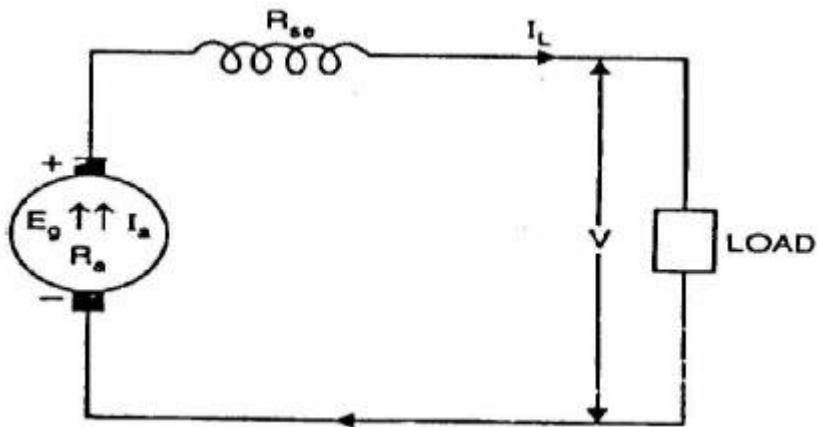
(i) Separately Excited D.C. Generators

A d.c. generator whose field magnet winding is supplied from an independent external d.c. source (e.g., a battery etc.) is called a separately excited generator. Fig shows the connections of a separately excited generator. The voltage output depends upon the speed of rotation of armature and the field current ($E_g = P\phi ZN/60$ A). The greater the speed and field current, greater is the generated e.m.f. It may be noted that separately excited d.c. generators are rarely used in practice. The d.c. generators are normally of self-excited type.



(ii) Self-Excited D.C. Generators

A d.c. generator whose field magnet winding is supplied current from the output of the generator itself is called a self-excited generator. There are three types of self-excited generators depending upon the manner in which the field winding is connected to the armature, namely;



Armature current, $I_a = I_{se} = I_L = I$ (say)

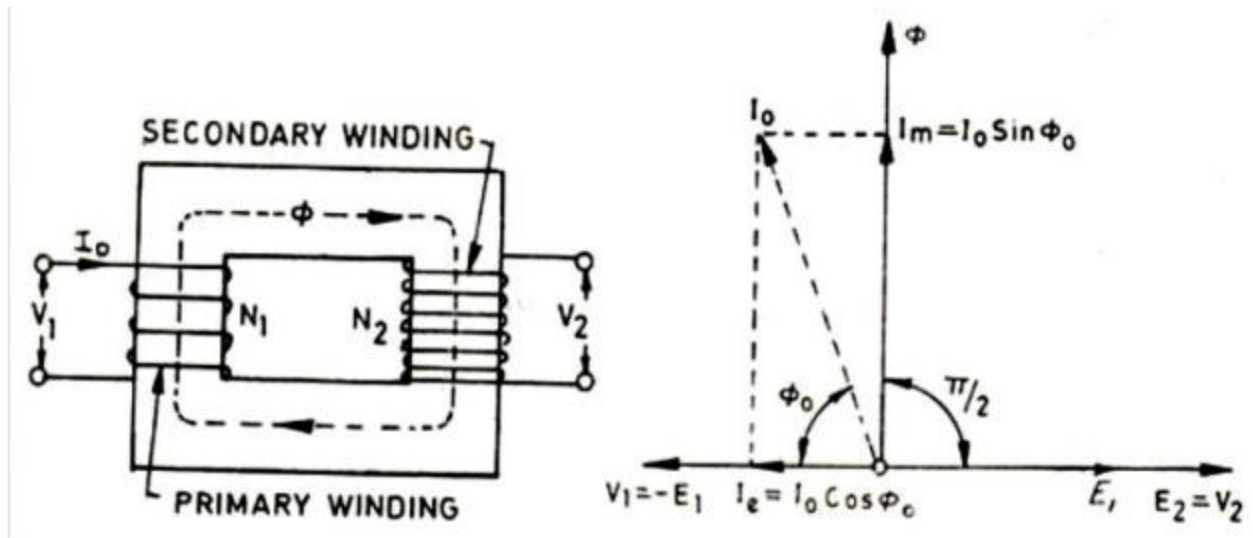
Terminal voltage, $V = E_g - I(R_a + R_{se})$

Power developed in armature = $E_g I_a$

Power delivered to load

9. Explain the transformer with no load condition in clear manner?

When the primary of a transformer is connected to the source of an ac supply and the secondary is open circuited, the transformer is said to be on no load. The Transformer on No Load alternating applied voltage will cause flow of an alternating current I_0 in the primary



winding, which will create alternating flux ϕ . No-load current I_0 , also known as excitation or exciting current, has two components the magnetizing component I_m and the energy component I_e . I_m is used to create the flux in the core and I_e is used to overcome the hysteresis and eddy current losses occurring in the core in addition to small amount of copper losses occurring in the primary only (no copper loss occurs in the secondary, because it carries no current, being open circuited.)

From vector diagram shown in above it is obvious that

1. Induced emfs in primary and secondary windings, E_1 and E_2 lag the main flux ϕ by $\pi/2$ and are in phase with each other.

2. Applied voltage to primary V_1 and leads the main flux ϕ by $\pi/2$ and is in phase opposition to E_1 .

3. Secondary voltage V_2 is in phase and equal to E_2 since there is no voltage drop in secondary.

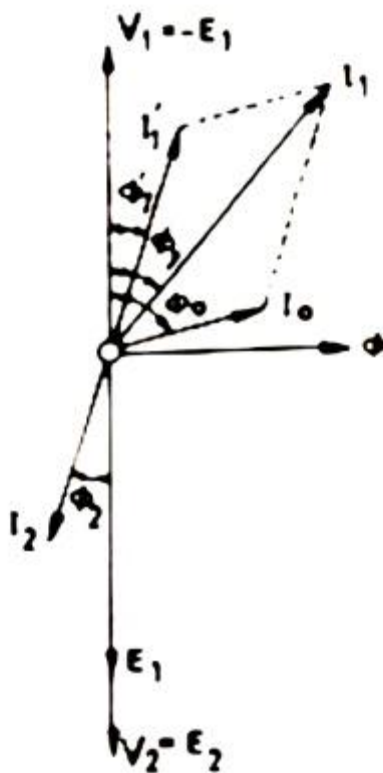
4. I_m is in phase with ϕ and so lags V_1 by $\pi/2$

5. I_e is in phase with the applied voltage V_1 .

6. Input power on no load = $V_1 I_e = V_1 I_0 \cos \phi_0$ where $\phi_0 = \tan^{-1} \frac{I_m}{I_e}$

10. Explain the transformer with load condition in clear manner?

The transformer is said to be loaded, when its secondary circuit is completed through an impedance or load. The magnitude and phase of secondary current (i.e. current flowing through secondary) I_2 with respect to secondary terminals depends upon the characteristic of the load i.e. current I_2 will be in phase, lag behind and lead the terminal voltage V_2 respectively when the load is non-inductive, inductive and capacitive. The net flux passing through the core remains almost constant from no-load to full load irrespective of load conditions and so core losses remain almost constant from no-load to full load. Vector diagram for an ideal transformer supplying inductive load is shown



Resistance and Leakage Reactance In actual practice, both of the primary and secondary windings have got some ohmic resistance causing voltage drops and copper losses in the windings. In actual practice, the total flux created does not link both of the primary and secondary windings but is divided into three components namely the main or mutual flux ϕ linking both of the primary and secondary windings, primary leakage flux ϕ_{L1} linking with primary winding only and secondary leakage flux ϕ_{L2} linking with secondary winding only. The primary leakage flux ϕ_{L1} is produced by primary ampere-turns and is proportional to primary current, number of primary turns being fixed.

The primary leakage flux ϕ_{L_1} is in phase with I_1 and produces self induced emf E_{L_1} is in phase with I_1 and produces self induced emf E_{L_1} given as $2\pi f L_1 I_1$ in the primary winding.

The self induced emf divided by the primary current gives the reactance of primary and is denoted by X_1 .

$$X_1 = E_{L_1}/I_1 = 2\pi f L_1 I_1/I_1 = 2\pi f L_1,$$

$$\text{leakage reactance of secondary } X_2 = E_{L_2}/I_2 = 2\pi f L_2 I_2/I_2 = 2\pi f L_2$$

Equivalent Resistance and Reactance. The equivalent resistances and reactance's of transformer windings referred to primary and secondary sides are given as below Referred to primary side Equivalent resistance,

Equivalent resistance, = X'_1 = Referred to secondary side Equivalent resistance,

Equivalent resistance, = $X_2 + K^2 X_1$ Where K is the transformation ratio.

UNIT III UTILIZATION OF ELECTRICAL POWER

PART A

1. What are the mechanisms for producing forces from wind?

There are two primary mechanism for producing forces from the winds. They are

- i. Lift force, and
- ii. Drag force

When lift force

2. Define Airfoil

Lift forces are produced by changing the velocity of the air stream flowing over either side of the lifting surface. Speeding up the air flow causes the pressure to drop, while slowing the air stream down leads to increase in pressure.

This pressure difference produces a force that begins to act on the high pressure side and moves towards the low pressure side of the lifting surface which is called an **airfoil**.

3. Define Magnus Effect

Magnus Effect, caused by spinning a cylinder in an air stream at a high-speed of rotation. The spin slows down the air speed on the side where the cylinder is moving into wind and increases it on the other side; the result is similar to an airfoil. This principle has been put to practical use in one or two cases but is not generally employed.

4. Define Stalling.

When lift decreases and the drag increases quite substantially; this phenomenon is known as **Stalling**. For efficient operation, a wind turbine blade needs to function with as much lift and as little drag as possible because drag dissipates energy.

5. What is the function of back-up in small producers?

For small producers, back-up can take the form of

- (1) Battery storage

- (2) Connection with the local electricity distribution systems, or
- (3) A stand by generator powered by liquid or gaseous fuels

Drag force

6. Define Power Co-efficient.

The fraction of the free-flow wind power that can be extracted by a rotor is called the power co-efficient. Thus

Power Coefficient = Power of wind rotor / Power available in the wind

Where, power available is calculated from the air density, rotor diameter and free wind speed as discussed earlier. The maximum theoretical power coefficient is equal to $16/27$ or 0.593 . This value cannot be exceeded by a rotor in free-flow wind-stream.

7. Why utilization of wind is considered as part of solar technology?

The major forcing function causing surface winds from the poles toward the equator is convective circulation. Solar radiation heats the air near the equator, and this low density heated air is buoyed up. At the surface it is displaced by cooler more dense higher pressure air flowing from the poles. In the upper atmosphere near the equator the air thus tends to flow back toward the poles and away from the equator. The net result is a global convective with surface winds from north to south in the northern hemisphere.

It is clear from the above over simplified model that the wind is basically caused by the Solar Energy irradiating the Earth. This is why wind utilization is considered part solar technology.

8. Define Wind.

Wind results from air motion. Air in motion arises from a pressure gradient. The circulation of air in the atmosphere is caused by the non-uniform heating of the earth's surface by the Sun.

9. What are the different types of forces acting on propeller type wind turbine.

There are two types of forces which are acting on the blades. They are

- (1) Circumferential force acting in the direction of wheel rotation that provides the torque, and

(2) Axial force acting in the wind stream that provides an axial thrust that must be countered by proper mechanical design.

10. What are the different types of vertical axis wind turbines.

i. Savonius Rotor type machines

ii. Darrieus type machines

PART B

1. Write short notes about Wind Power?

All renewable energy (except tidal and geothermal power), ultimately comes from the sun. The earth receives 1.74×10^{17} watts of power (per hour) from the sun. About one or 2 percent of this energy is converted to wind energy (which is about 50-100 times) more than the energy converted to biomass by all plants on earth.

Differential heating of the earth's surface and atmosphere induces vertical and horizontal air currents that are affected by the earth's rotation and contours of the land - > WIND. E.g.: Land Sea Breeze Cycle. Winds are influenced by the ground surface at altitudes up to 100 meters. Wind is slowed by the surface roughness and obstacles.

When dealing with wind energy, we are concerned with surface winds. A wind turbine obtains its power input by converting the force of the wind into a torque (turning force) acting on the rotor blades. The amount of energy which the wind transfers to the rotor depends on the density of the air, the rotor area, and the wind speed.

The kinetic energy of a moving body is proportional to its mass (or weight). The kinetic energy in the wind thus depends on the density of the air, i.e. its mass per unit of volume. In other words, the "heavier" the air, the more energy is received by the turbine at 15°C air weight about 1.225 kg per cubic meter, but the density decreases slightly with increasing humidity.

A typical 600 kW wind turbine has a rotor diameter of 43-44 meters, i.e., a rotor area of some 1,500 square meters. Fig 5.1 shows the power generated by the wind mill with respect to the height.

The rotor area determines how much energy a wind turbine is able to harvest from the wind. Since the rotor area increases with the square of the rotor diameter, a turbine which is twice as large will receive $2^2 = 2 \times 2 =$ Four times as much energy.

To be considered a good location for wind energy, an area needs to have average annual wind speeds of at least 12 miles per hour.

They have traditionally been measured at a standard height of ten meters where they are found to be 20-25% greater than close to the surface. At a height of 60 metre they may be 30-60% higher because of the reduction in the drag effect of the earth surface.

2. Write short notes about the Wind Power calculation?

Calculation of Power in the Wind

The power in the wind can be computed by using of Kinetics (Kinetic means relating to or resulting from motion). The wind mill works on the principle of converting Kinetic energy of the wind to mechanical energy.

We know that power is equal to energy per unit time. The energy available is the kinetic energy of the wind. The kinetic energy of any particle is equal to one half its mass times the square of its velocity.

Kinetic Energy of particle = $\frac{1}{2} mv^2$

Where

M : Mass of particle (kg)

V : Velocity of particle (m/s)

The amount of air passing in unit time, through an area 'A', with velocity 'V' is $A \times V$, and its mass 'm' is equal to its volume multiplied by its density 'ρ' of air.

$$m = \rho AV$$

Where, m is the mass of air transversing the area 'A' swept by the rotating blades of a wind mill type generator.

Power Coefficient

The fraction of the free-flow wind power that can be extracted by the rotor is called the power co-efficient; Thus,

Power Coefficient = Power of wind rotor / Power available in the wind

Where, power available is calculated from the air density, rotor diameter and free wind speed as discussed earlier. The maximum theoretical power coefficient is

equal to $16/27$ or 0.593. This value cannot be exceeded by a rotor in free-flow wind-stream.

An ideal rotor, with propeller-type blades of proper aerodynamic design, would have a power co-efficient approaching 0.59. But such a rotor would not be strong enough to withstand the stresses to which it is subjected when rotating at a high rate in a high-speed wind stream.

3. Write short notes about the circuit breakers?

Braking is very frequent in electric drives to stop a motor in a reasonably short time.

- 1) Reliable and quick in its action.
- 2) The braking force must be capable of being controlled.
- 3) Adequate means be provided for dissipating the stored energy that is kinetic energy of the rotating parts.
- 4) In case of a fault in any part of the braking system the whole system must come to instantaneous rest or result in the application of the brakes.

There are two types of braking:

1) Mechanical braking:

The motor in this case is stopped due to friction between the moving part of the motor and the brake shoe that is stored energy is dissipated as heat by a brake shoe or brake lining which rubs against a brake shoe or brake lining which rubs against a brake drum.

2) Electric braking:

the kinetic energy of the moving parts that is motor is converted into electrical energy which is consumed in a resistance as heat or alternatively it is returned to the supply source. During braking operation a motor has to function as a generator. The motor can be held at stand still. In other words the electric braking cannot hold the motor at rest. Thus it becomes essential to provide mechanical brakes in addition to electric braking.

Electric braking:

Various types of electrical braking are:

a) Plugging

b) Rheostatic braking

c) Regenerative braking

4. Write short notes about the illumination in detailed manner?

Luminous intensity

Luminous intensity is an expression of the amount of light power emanating from a point source within a solid angle of one steradian.

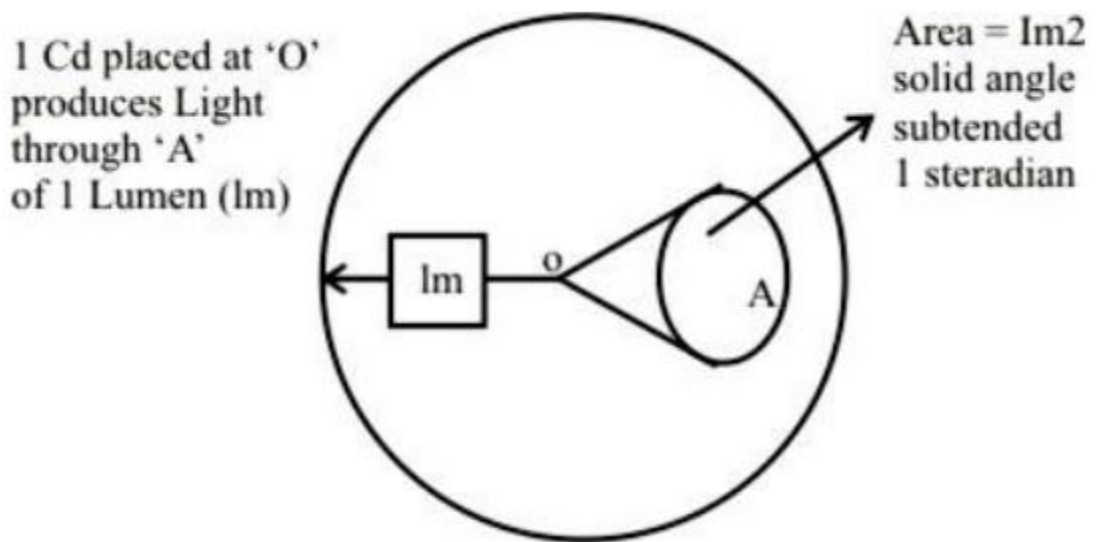
Laws of illumination:

The original standard of light was Wax Candle, which is highly unreliable. It was replaced by a Vaporized Pentane Lamp. This is equal to 10 original Candles. In the year 1909, Incandescent Lamp was taken as standard by comparison with a Pentane Lamp. Thing to be kept in mind is Primary Standard should be reproducible. It was in 1948, Luminous Intensity; based on Luminance (objective brightness) of a small aperture due to Light from a Radiator maintained at 1773°C i.e. Solidification temperature of platinum was adopted as Standard. It consists of:

1. Radiator – Fused Thoria – Thorium Oxide. 45mm long internal dia of 2.5mm. Packed with Fused Thoria Powder at the bottom.
2. Supported Vertically Pure Platinum in a Fused thoria crucible with a small aperture of 1.5mm in a large refractory container.
3. Platinum melted by a High Frequency Eddy current. Luminance = 589000 Candles /m² ≈ 600 000 units.

Transparent (Law of Inverse Squares):

Common unit of light intensity is candela. It is Luminous intensity in the Perpendicular direction of a surface, $1 / 600,000$ of a black body at temperature of solidification or Freezing of Platinum under Standard Atmospheric pressure. It is abbreviated as Cd. It is indicative of Light Radiating Capacity of a source of Lamp.



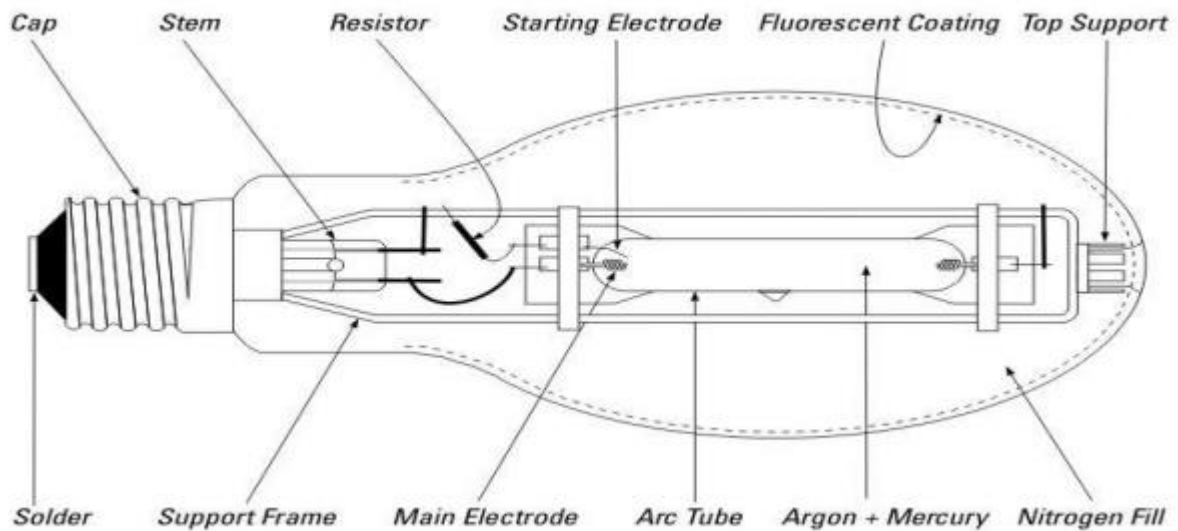
Consider a transparent sphere of radius 1m shown in Fig. If we place a 1 Cd source at the centre then light flux coming out through an area of 1m² over 1 steradian solid angle will be 1 lumen. Thus Luminous Intensity over 1 Str. by 1, Cd, we call it 1 lumen \approx 1 lm. Basic unit of Light Flux.

$$\phi = I 4\pi \text{lumens}$$

5. Write short notes about the Sodium Vapour Lamp in detailed manner?

Principally sodium vapour lamp consists of a bulb containing a small amount of metallic sodium, neon gas and two sets of electrodes connected to a pin type base. The lamp operates at a temperature of about 300°C and in order to conserve the heat generated and assure the lamp operating at normal air temperatures the discharge envelope is enclosed in special vacuum envelope designed for this purpose.

The efficiency of a sodium vapour lamp under practical conditions is about 40-50 lumens/watt. Such lamps are manufactured in 45, 60, 85 and 140 W ratings. The average life is about 3000 hours and is not affected by voltage variations. The major application of this type of lamp is for highway and general outdoor lighting where colour discrimination is not required, such as street lighting, parks, rail yards, storage yards.



6. Write short notes about the fluorescent Lamp in detailed manner?

Employs transformation of UV radiation due to low pressure mercury vapor. Luminescent Powder in tubular vapor Lamps Enhances brilliancy of light. Radiation from Low Pressure Mercury Vapor (which is in UV region) is impinged on Luminescent Materials and re – radiated at longer wavelengths of visible spectrum. In a Glass Tube small drop of Mercury and small amount of Argon gas are placed for initiation of discharge. Pressure, voltage and current are so adjusted that 253.7 nm line is excited. This re-radiates at longer wavelength. Typically a 40W lamp requires 2-3g of phosphors. Maximum sensitivity is around 250 – 260 nm. Various types of Fluorescent Lamps are:

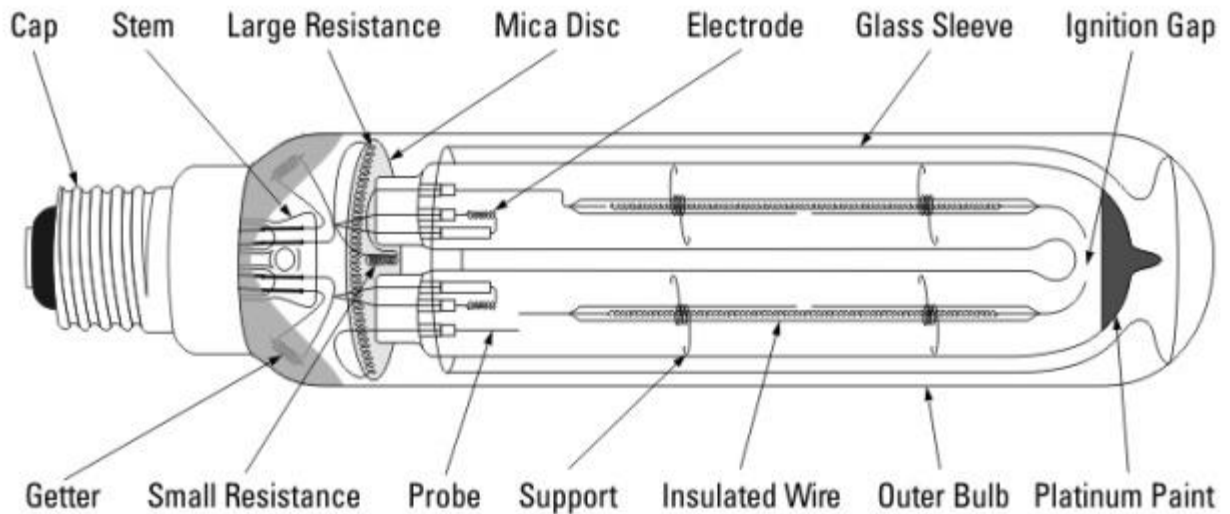
1. Day Light Fluorescent Lamps- Average Noon Day Light. 6500°k suitable where demands are not exacting
2. Standard white Light - 3500°k general Lighting.
3. 4500°k white Lamp – between std. white Light & Day Light Lamp.
4. Soft white Lamp – Pinker Light. 25% lower light output than Std. white Lamp suitable for Residential lighting and Restaurants.

Dimension and Voltage depend on Luminous Efficacy, Brightness, Lumen Output and Lumen Maintenance. Reliable Starting is achieved by having preheated cathodes / hot cathode. Half the open circuit voltage should be used by the Lamp and the other half by the ballast. Lamp Voltage decides the arc

length, bulb diameter and lamp current. Hot Cathode lamps operate at lower voltage < cold Cathode lamps. Typically cold cathodes have 70-100V drop at the cathode.

7. Write short notes about the mercury vapour Lamp in detailed manner?

High Pressure Mercury Vapour Lamp:



The mercury vapour lamp in construction is similar to sodium vapour lamp. It gives greenish blue colour light, which causes colour distortion. The efficiency is about 30-40 lumens per watt.

These lamps (MA type) are manufactured in 250 and 400 W ratings for use on 200-250 Vac supply. Lamps of this type are used for general industrial lighting, railway yards, ports, work areas; shopping centers etc where greenish-blue colour light is not objectionable.

Another type, which is manufactured in 300 and 500 W ratings for use on ac as well as dc supply mains is MAT type. This is similar to MA type except that it does not use choke as ballast.

Lower wattage lamps, such as 80 and 125 W, are manufactured in a different design and using high vapour pressure of about 5-10 atmospheres. These are known as MB type lamps.

8. Write short notes about the charge and discharge characteristics of LI ION in detailed manner?

nickel and cobalt have been replaced by cobalt has been replaced by nickel and manganese right, and here it is only manganese and this is of course, is a different, the same structure the Mn_2O_4 , the Mn_2O_4 will have a spinal structure, but this is also a same structure ok, so these are some of the layers. Next oxides and here this is the characteristics charge, charge discharge characteristics of the particular compound, which is spinal structure lithium manganese 4, manganese 4 spinal structural, we have seen earlier.

This is the charging characteristics of that, and these are the discharge characteristics under different conditions of discharging, now discharging is done normally, done at a rate of designated like C by 2 or C by 5 or 20 C the charge rate or discharge rate, the charge rate is denoted as C or C rate, and signifies a charge or discharge rate, charge or discharge rate equal to the capacity of the battery in in 1 hour

the capacity of a battery is normally given as ampere hour, that is the total amount of energy which it can accommodate. So, if you have a capacity of a particular batteries 1.6 ampere hour, the C 1 C it becomes 1.6 ampere, that means in 1 hour, 1 hour you charge it fully, or of discharge it fully at the rate of 1.6 amp. So, if you draw current of 1.6 ampere over a period of 1 hour, then this is the total capacity of the battery, and this becomes then the C, a charge rate of C by 2 means actually 0.8 ampere, 0.8 ampere, half of 1.6 will be needed for 2 hours, and a charge rate of 2 C that means the double the value of the capacity, we have to charge 3.2 ampere would need 30 minutes half a minute to fully charged. So, it is basically there rate at which the current is being drawn from the, from the battery or it is being introduced the charges, introduced battery or to recharge the battery, so that is what is known as C rate, it assumes that the battery of course, whenever we are doing it assumed the battery is 100 percent efficient at absorbing the charge, and can support the rate that means, if you charges it at a very fast rate. For example, 2 C or 10 C, 20 C that means your charging or discharging at a very fast rate, and the kinetics may not support it. So, there may be a breakdown of the structure, this can be done in fact, that is one of the criteria when it battery is tested.

charging is done at C by 2 rate, half the total capacity, and discharging has been tried out, discharge have been tried out at C by 5 and in this increasing 20 C, 20 times the overall the capacity of the battery, and even then one can although the voltage will drop obviously, and charging rate increase, and then you can get a fairly high a specific charges. So, MnO_4 , so $LiMnO_4$ an alternative cathode

material replacing LiCO_2 , so these are the some of the typical experimental data, and which is been taken from literature of course, and the these are the possibilities, so one can charge it was very fast rate, as well as discharge fairly at a very fast rate.

9. Write short notes about the domestic refrigerator in detailed manner?

refrigeration is the cooling of air/liquids, thus providing lower temperatures to preserve food, cool beverages, make ice and for many other applications.

refrigeration is the process of removing heat from a place where it is not wanted and transferring that heat to a place where it makes little difference.

In the average household, the room temperature from summer to winter is normally between 70°F and 90°F.

The temperature inside the refrigerator fresh food section should be about 35°F

The rating of a refrigeration machine is obtained by refrigerating effect or the amount of heat extracted in a given time from a body or space.

One tonne of refrigeration is defined as the refrigerating effect (amount of heat extracted) produced by melting 1 tonne of ice from and at 0°C in 24 hours. We know that latent heat of fusion of ice = 336 kJ/kg.

Refrigerating effect of this heat in terms of tonne in 24 hours is rated as tonne of refrigeration

$$1\text{T.R} = 336 \times 1000 / 24 = 14000\text{kJ/hr.}$$

$$= 210\text{kJ/min.} = 3.85\text{ kJ/se}$$

Performance of the refrigerator is determined by using co-efficient of performance which is defined as required input. In desired output from the conservation of energy principle.

10. Write short notes about the earthing system in detailed manner?

Need of Earthing or Grounding.

The primary purpose of earthing is to avoid or minimize the danger of electrocution, fire due to earth leakage of current through undesired path and to ensure that the potential of a current carrying conductor does not rise with respect to the earth than its designed insulation.

When the metallic part of electrical appliances (parts that can conduct or allow passage of electric current) comes in contact with a live wire, maybe due to failure of installations or metal become charged and static charge accumulates on it. To avoid such instances, the power supply systems and parts of appliances have to be earthed so as to transfer the charge directly to the earth.

Below are the basic needs of Earthing.

- To protect human lives as well as provide safety to electrical devices and appliances from leakage current.
- To keep voltage as constant in the healthy phase (If fault occurs on any one phase).
- To Protect Electric system and buildings form lightning.
- To serve as a return conductor in electric traction system and communication.
- To avoid the risk of fire in electrical installation systems.

Different Terms used in Electrical Earthing

- **Earth:** The proper connection between electrical installation systems via conductor to the buried plate in the earth is known as Earth.
- **Earthed:** When an electrical device, appliance or wiring system connected to the earth through earth electrode, it is known as earthed device or simple “Earthed”.
- **Solidly Earthed:** When an electric device, appliance or electrical installation is connected to the earth electrode without a fuse, circuit breaker or resistance/Impedance, It is called “solidly earthed”.
- **Earth Electrode:** When a conductor (or conductive plate) buried in the earth for electrical earthing system. It is known to be Earth Electrode. Earth electrodes are in different shapes like, conductive plate, conductive rod, metal water pipe or any other conductor with low resistance.
- **Earthing Lead:** The conductor wire or conductive strip connected between Earth electrode and Electrical installation system and devices in called Earthing lead.
- **Earth Continuity Conductor:** The conductor wire, which is connected among different electrical devices and appliances like, distribution board, different plugs and appliances etc. in other words, the wire between earthing lead and electrical device or appliance is called earth continuity conductor. It may be in the shape of metal pipe (fully or partial), or cable metallic sheath or flexible wire.

- **Sub Main Earthing Conductor:** A wire connected between switch board and distribution board i.e. that conductor is related to sub main circuits.
- **Earth Resistance:** This is the total resistance between earth electrode and earth in Ω (Ohms). Earth resistance is the algebraic sum of the resistances of earth continuity conductor, earthing lead, earth electrode and earth.

points to be earthed

Earthing is not done anyhow. According to IE rules and IEE (Institute of Electrical Engineers) regulations,

- Earth pin of 3-pin lighting plug sockets and 4-pin power plug should be efficiently and permanently earthed.
- All metal casing or metallic coverings containing or protecting any electric supply line or apparatus such as GI pipes and conduits enclosing VIR or PVC cables, iron clad switches, iron clad distribution fuse boards etc should be earthed (connected to earth).
- The frame of every generator, stationary motors and metallic parts of all transformers used for controlling energy should be earthed by two separate and yet distinct connections with the earth.
- In a dc 3-wire system, the middle conductors should be earthed at the generating station.
- Stay wires that are for overhead lines should be connected to earth by connecting at least one strand to the earth wires.

UNIT IV ELECTRONIC CIRCUITS

PART A

1. What is meant by Q-factor?

Q-factor is known as the quality factor. It is used to measure the quality factor of the coils such as inductors, Capacitors etc.

2. Define transistor action.

A transistor consists of 2 coupled PN junctions. The base is a common region to both junctions and makes a coupling between them. Since the base regions are smaller, a significant interaction between junctions will be available. This is called transistor actions.

3. Define hybrid parameters.

Any linear circuit having input and output terminals can be analysed by four parameters(one measured on ohm, one in mho and two dimensionless) called hybrid or hparameters.

4. Which is the most commonly used transistor configuration? Why?

The CE Configuration is most commonly used. The reasons are

- High Current gain
- High voltage gain
- High power gain
- Moderate input to output ratio.

5. What are the types of transistors?

- Unipolar junction transistor
- Bipolar junction transistor.

6. Define delay time

It is defined as the time required for the current to rise from 0 to 10% of its maximum value.

7. Define rise time

It is the time required for the current to rise from 0 to 90 percentage of the maximum value.

8. Define Biasing

“**Biasing**” is providing minimum external voltage and current to activate the device to study its characteristics.

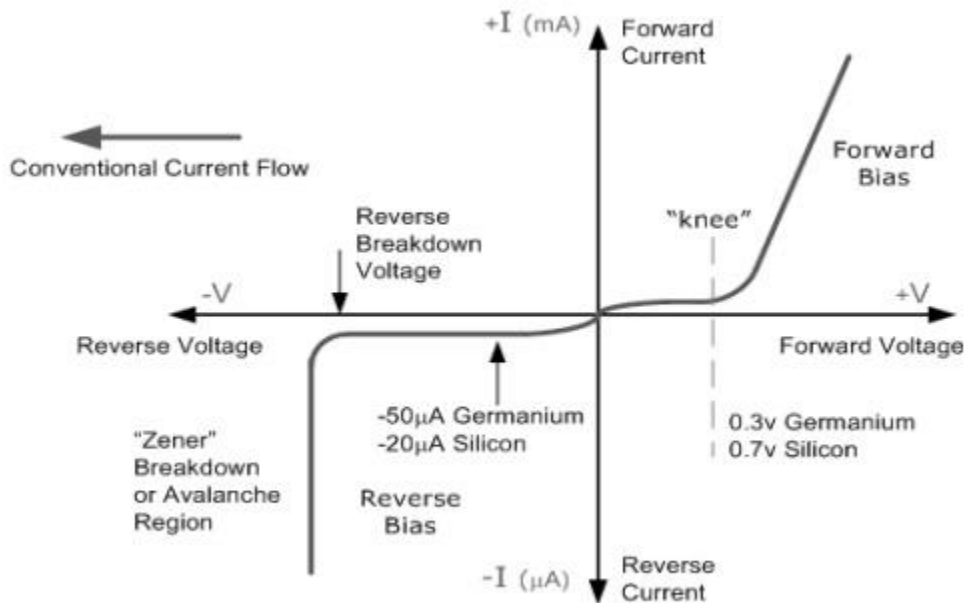
There are two operating regions and two "biasing" conditions for the standard Junction Diode and they are:

1. Zero Bias
2. Forward Bias
3. Reverse Bias

9. What is meant by zener effect?

When the doping is heavy, even the reverse voltage is low, the electric field at barrier will be so strong thus the electrons in the covalent bonds can break away from the bonds. This effect is known as **Zener effect**.

10. Draw thw V/I characteristics of PN diode?



PART B

1.Explain the V/I characteristics of ZENER diode?

Zener Diode

A diode which exhibits the zener effect is called a Zener Diode. Hence it is defined as a reverse biased heavily doped PN junction diode which operates in breakdown region.

zener effect

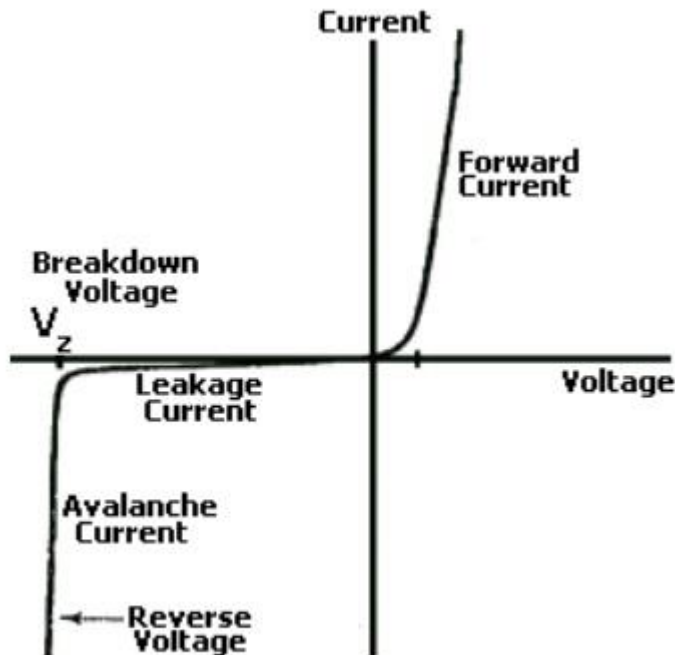
In a general purpose PN diode the doping is light; as a result of this the breakdown voltage is high. If a P and N region are heavily doped then the breakdown voltage can be reduced. When the doping is heavy, even the reverse voltage is low, the electric field at barrier will be so strong thus the electrons in the covalent bonds can break away from the bonds. This effect is known as **Zener effect**.

ZENER DIODE



A diode which exhibits the zener effect is called a **Zener Diode**. Hence it is defined as a reverse biased heavily doped PN junction diode which operates in breakdown region. The zener diodes have been designed to operate at voltages ranging from a few volts to several hundred volts. **Zener Breakdown** occurs in junctions which is heavily doped and have narrow depletion layers. The breakdown voltage sets up a very strong electric field. This field is so strong enough to break or rupture the covalent bonds thereby generating electron hole pairs. Even a small reverse voltage is capable of producing large number of current carrier. When a zener diode is operated in the breakdown region care must be taken to see that the power dissipation across the junction is within the power rating of the diode otherwise heavy current flowing through the diode may destroy it.

V-I characteristics of Zener diode



The illustration above shows this phenomenon in a current vs voltage graph with a zener diode connected in the forward direction .It behaves exactly as a standard diode. In the reverse direction however there is a very small leakage current between 0v and the zener voltage –i.e. just a tiny amount of current is able to flow. Then, when the voltage reaches the breakdown voltage (V_z), suddenly current can flow freely through it.

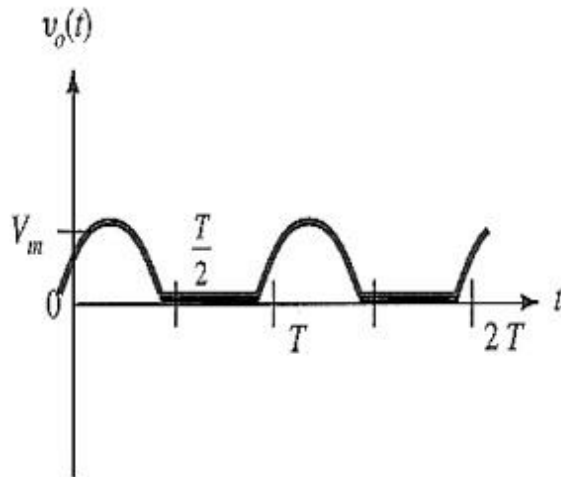
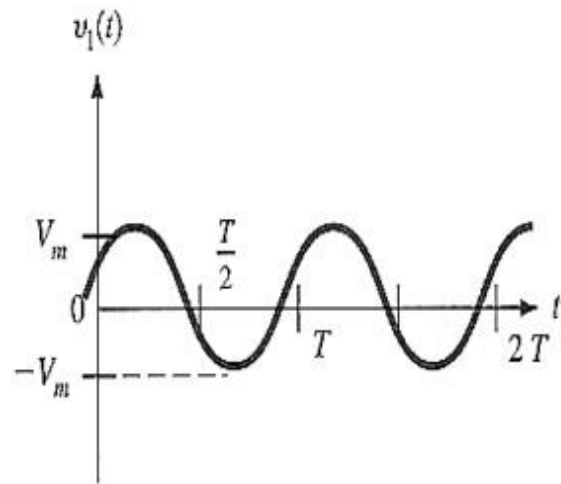
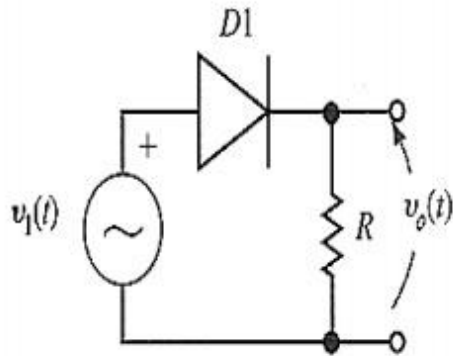
Application of Zener diode

- a) as voltage regulator
- b) as peak clippers
- c) for reshaping waveforms

2. Briefly describe the operation of Half wave Rectifier?

Principle

It is a circuit that converts alternating voltage or current into pulsating voltage or current for half the period of input cycle hence it is named as “**half wave rectifier**”



Operation

- During the positive half cycle of input, the diode D is forward biased, it offers very small resistance and it acts as closed switch and hence conducts the current through the load resistor.
- During the negative half cycle of the input diode D is heavily reverse biased, it offers very high resistance and it acts as open switch hence it does not conduct any current. The rectified output voltage will be in phase with AC input voltage for completely resistive load.

Construction

- It consists of step-down transformer, semiconductor diode and the load resistance.
- The step-down transformer – reduce the available ac voltage into required level of smaller ac voltage.
- The diode can be used to convert the ac into pulsating dc.

3. Briefly describe the operation of Full wave Rectifier?

Principle

A circuit that converts the ac voltage or current into pulsating voltage or current during both half cycle of input is known as “full wave rectifier”.

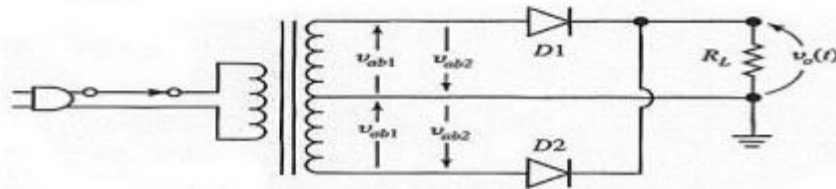
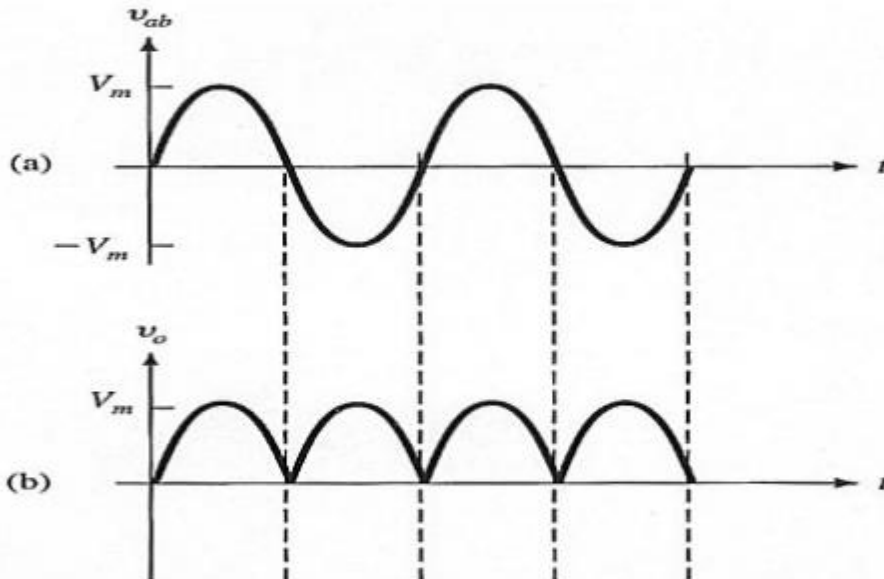
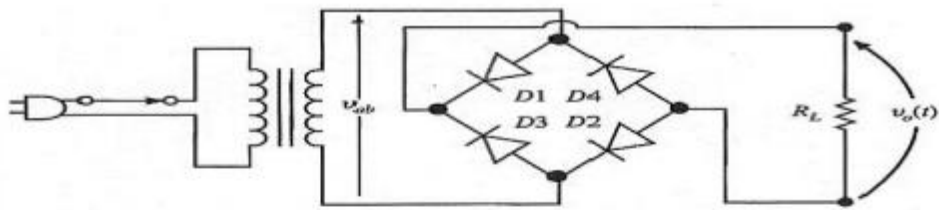
Operation

- During positive half cycle of ac input, diode D_1 becomes forward biased, provides very small resistance and acts as closed switch, resulting in the flow of current.
- During negative half cycle, diode D_1 reverse biased, offers high resistance and it acts as open circuit.

The “rectifier” is a circuit that converts AC voltages and currents into pulsating DC voltages and currents. It consists of DC components and the unwanted ac ripple or harmonic components which can be removed by using filter circuit. Thus the output obtained will be steady DC voltage and magnitude of DC voltage can be varied by varying the magnitude of AC voltage.

Filters: A circuit that removes ripples (unwanted ac components) present in the pulsating dc voltage.

Regulator: A circuit that maintains the terminal voltage as constant even if the input voltage or load current varying.



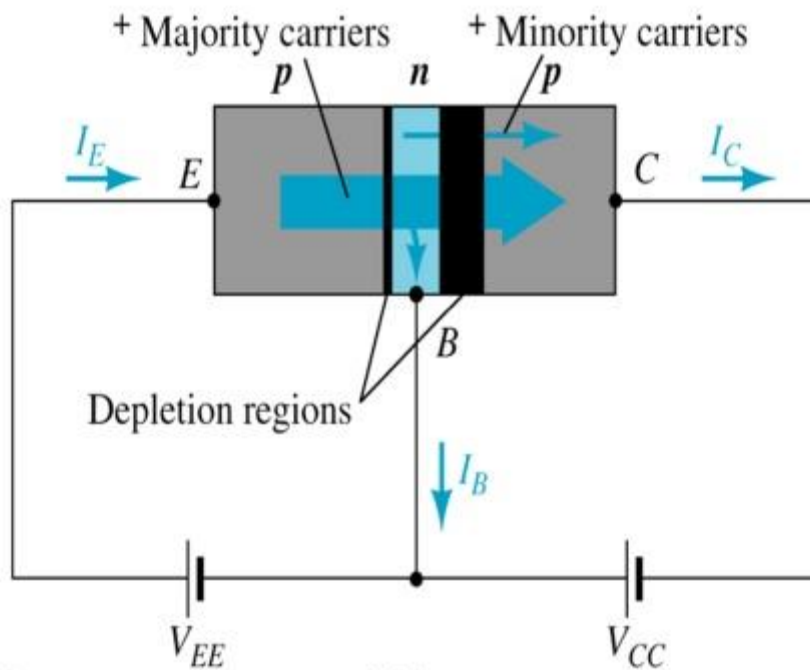
4. Explain the operation of BIPOLAR JUNCTION TRANSISTOR?

A bipolar junction transistor is a three terminal semiconductor device in which the operation depends on the interaction of majority and minority carriers. Transistor refers to Transfer Resistor i.e., signals are transferred from low resistance circuit into high resistance circuit. BJT consists of silicon crystal in which a layer of 'N' type silicon is sandwiched between two layers of 'P' type silicon. The semiconductor sandwiched is extremely smaller in size. In other words, it consists of two back to back PN junction joined together to form single piece of semiconductor crystal. These two junctions give three regions called Emitter, Base and Collector.

There are two types of transistors such as PNP and NPN. The arrow on the emitter specifies whether the transistor is PNP or NPN type and also determines the direction of flow of current, when the emitter base junction is forward biased.

Operation of Transistor

The basic operation will be described using the pnp transistor. The operation of the pnp transistor is exactly the same if the roles played by the electron and hole are interchanged. One p-n junction of a transistor is reverse-biased, whereas the other is forward-biased. Both biasing potentials have been applied to a pnp transistor and resulting majority and minority carrier flows indicated.



Majority carriers (+) will diffuse across the forward-biased p-n junction into the n-type material. A very small number of carriers (+) will through n-type material to the base terminal. Resulting I_B is typically in order of microamperes. The large number of majority carriers will diffuse across the reverse-biased junction into the p-type material connected to the collector terminal. Majority carriers can cross the reverse-biased junction because the injected majority carriers will appear as minority carriers in the n-type material.

Applying KCL to the transistor :

$$I_E = I_C + I_B$$

The comprises of two components – the majority and minority carriers

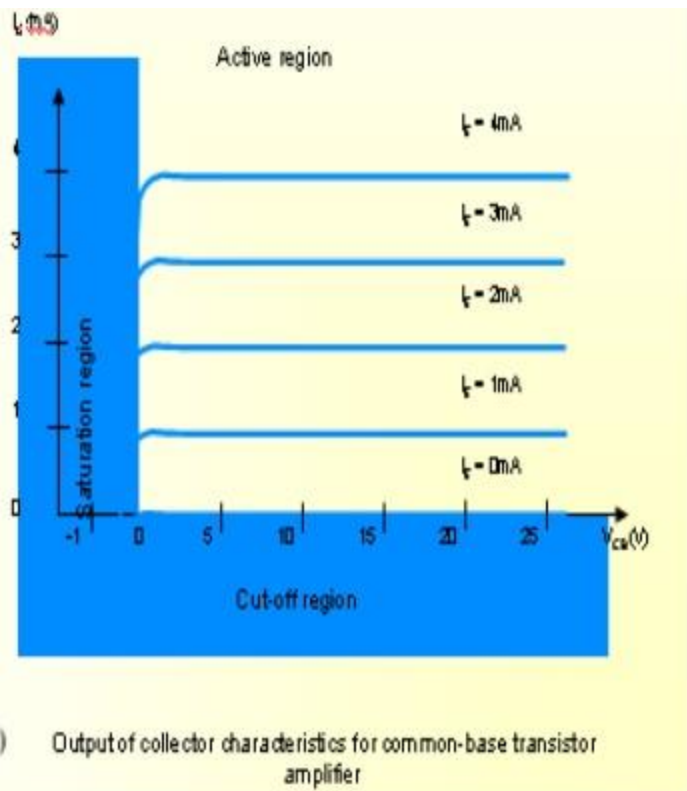
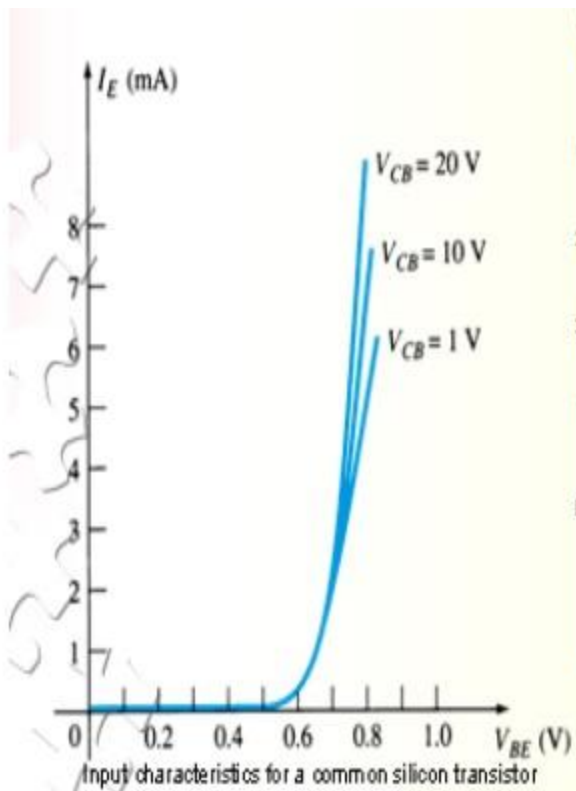
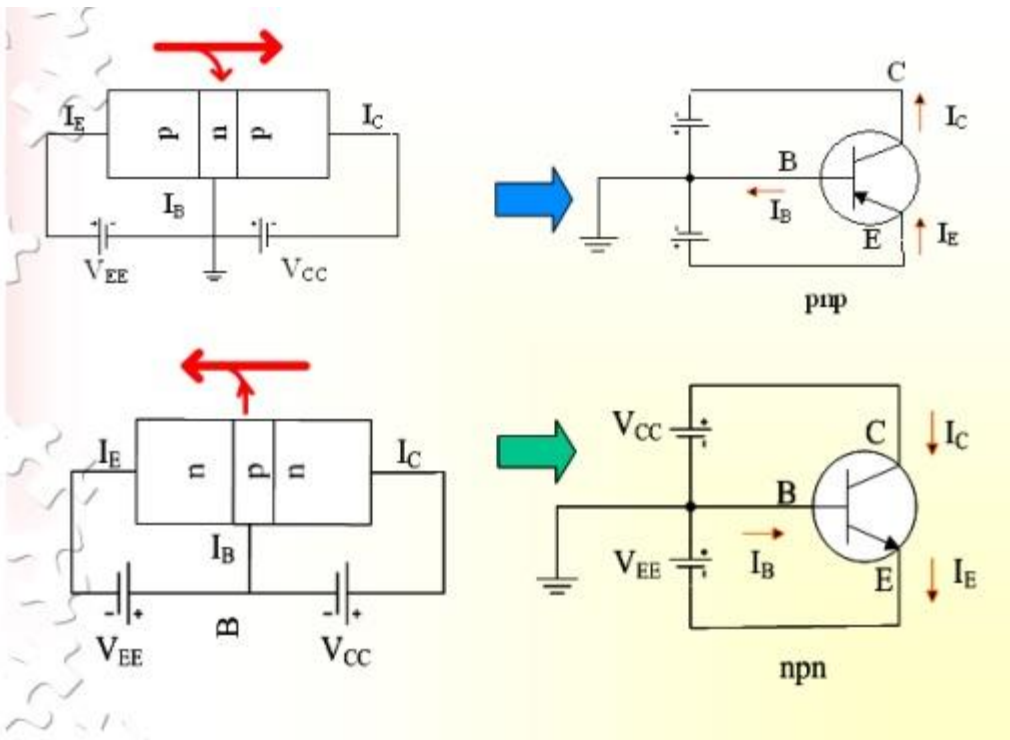
$$I_C = I_{C_{majority}} + I_{C_{minority}}$$

$I_{CO} - I_C$ current with emitter terminal open and is called leakage current.

5. Describe the Common Base configuration of transistor?

- Common-base terminology is derived from the fact that the :
 - base is common to both input and output of the configuration.
 - base is usually the terminal closest to or at ground potential.
- All current directions will refer to conventional (hole) flow and the arrows in all electronic symbols have a direction defined by this convention.
- Note that the applied biasing (voltage sources) are such as to establish current in the direction indicated for each branch.
- To describe the behavior of common-base amplifiers requires two set of characteristics:
 - 1 . Input or driving point characteristics.
 2. Output or collector characteristics
- The output characteristics has 3 basic regions:
 1. Active region –defined by the biasing arrangements
 - 2 . Cutoff region – region where the collector current is 0A
 - 3 . Saturation region- region of the characteristics to the left of

$$V_{CB} = 0V$$



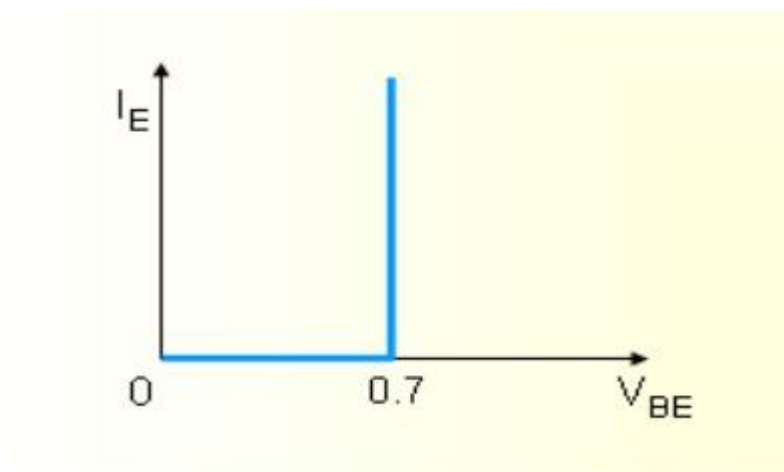
| Active region | Saturation region | Cut-off region |
|--|--|--|
| <ul style="list-style-type: none"> • I_E increased, I_C increased • BE junction forward bias and CB junction reverse bias • Refer to the graf, $I_C \approx I_E$ • I_C not depends on V_{CB} • Suitable region for the transistor working as amplifier | <ul style="list-style-type: none"> • BE and CB junction is forward bias • Small changes in V_{CB} will cause big different to I_C • The allocation for this region is to the left of $V_{CB} = 0$ V. | <ul style="list-style-type: none"> • Region below the line of $I_E = 0$ A • BE and CB is reverse bias • no current flow at collector, only leakage current |

The curves (output characteristics) clearly indicate that a first approximation to the relationship between I_E and I_C in the active region is given by

$$I_C \approx I_E$$

> Once a transistor is in the 'on' state, the base-emitter voltage will be assumed to be

$$V_{BE} = 0.7V$$



- In the dc mode the level of I_C and I_E due to the majority carriers are related by a quantity called alpha

$$\alpha = I_C / I_E$$

$$I_C = \alpha I_E + I_{CBO}$$

- It can then be summarize to $I_C = \alpha I_E$ (ignore I_{CBO} due to small value)
- For ac situations where the point of operation moves on the characteristics curve, an ac alpha defined by
- Alpha a common base current gain factor that shows the efficiency by calculating the current percent from current flow from emitter to collector. The value of is typical from

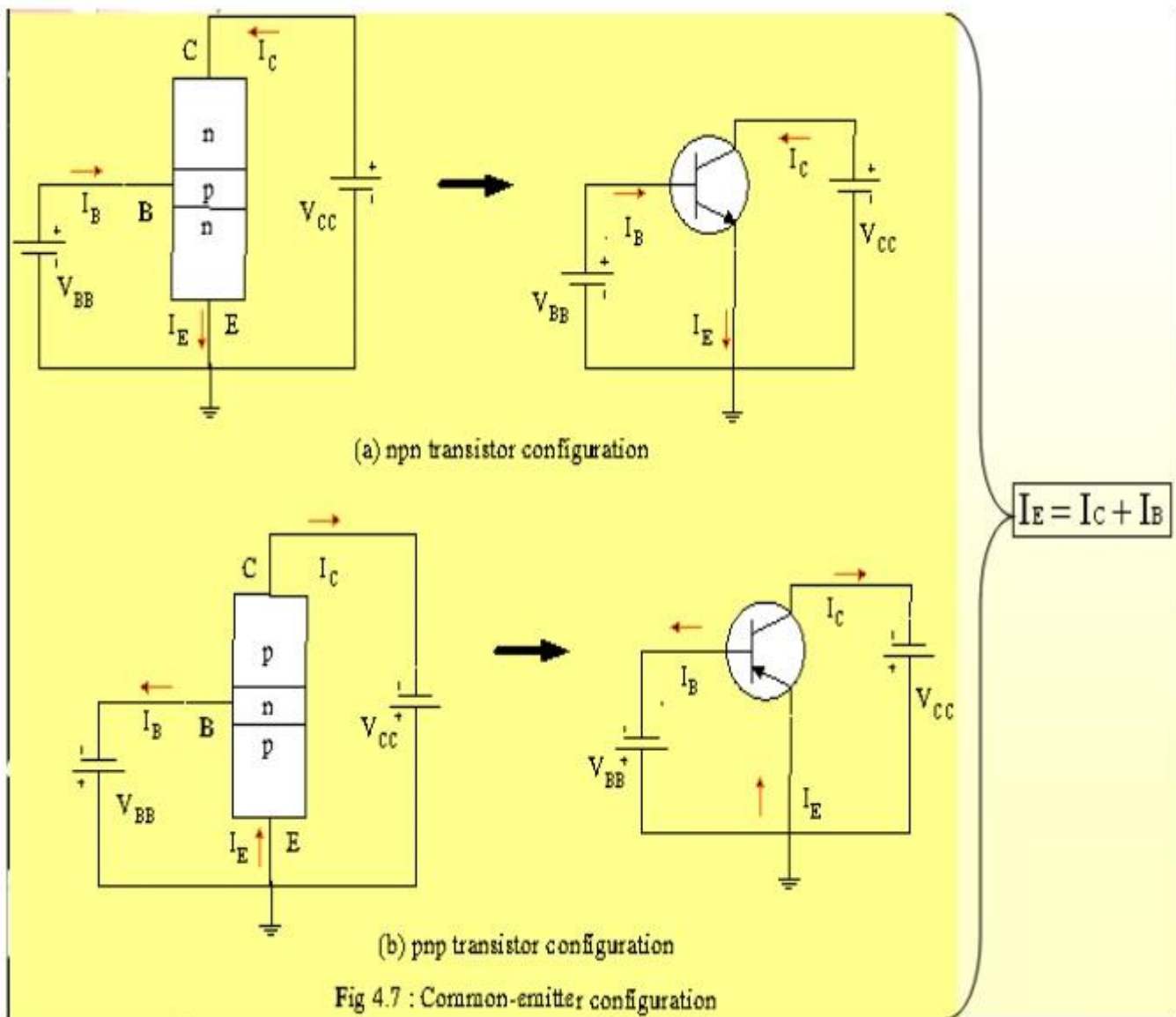
$$0.9 \sim 0.998.$$

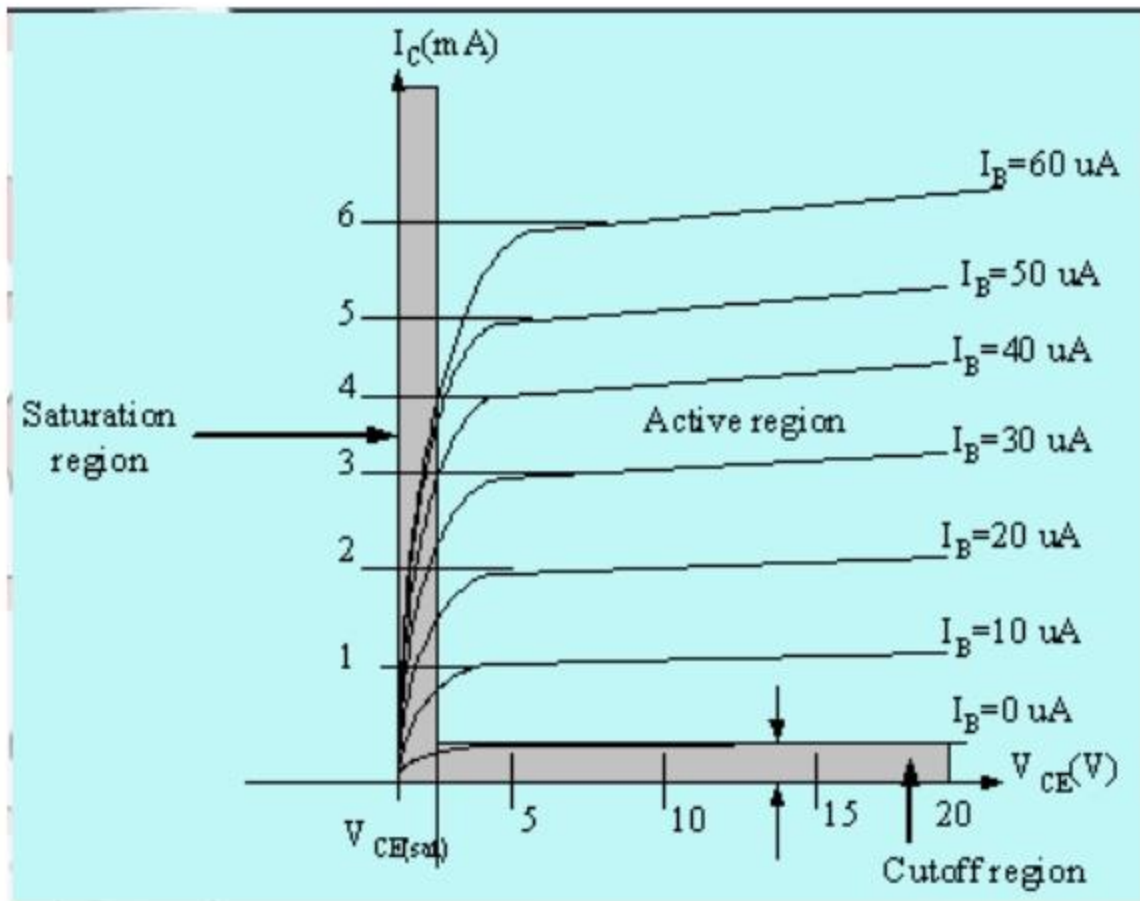
6. Describe the Common Emitter configuration of transistor?

- It is called common-emitter configuration since :
 - emitter is common or reference to both input and output terminals.
 - emitter is usually the terminal closest to or at ground potential.
- Almost amplifier design is using connection of CE due to the high gain for current and voltage.
- Two set of characteristics are necessary to describe the behavior for CE; input (base terminal) and output (collector terminal) parameters.

Input characteristics for CE configuration

- I_B in microamperes compared to milliamperes of I_C .
- I_B will flow when $V_{BE} > 0.7V$ for silicon and $0.3V$ for germanium
- Before this value I_B is very small and no I_B .
- Base-emitter junction is forward bias
- Increasing V_{CE} will reduce I_B for different values.





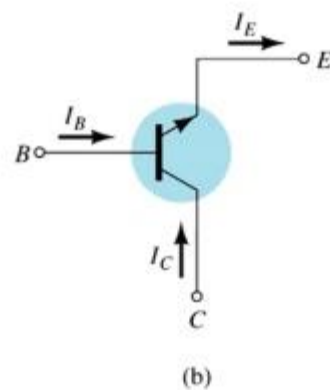
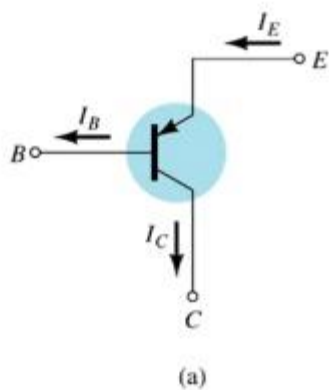
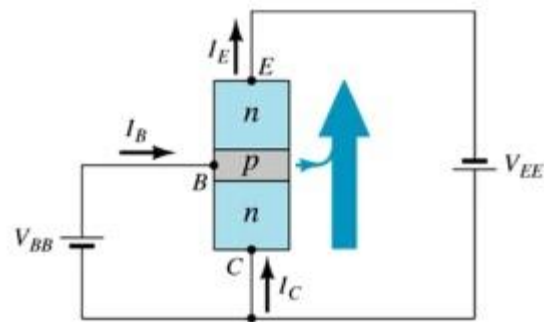
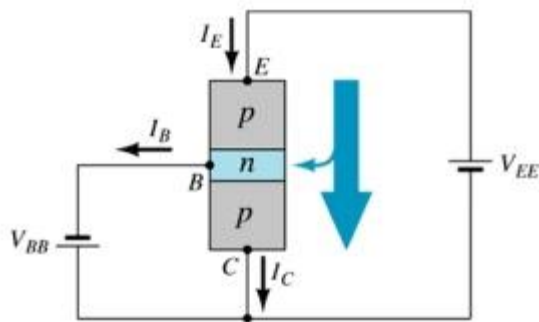
Output characteristics for CE configuration

Output characteristics for CE configuration

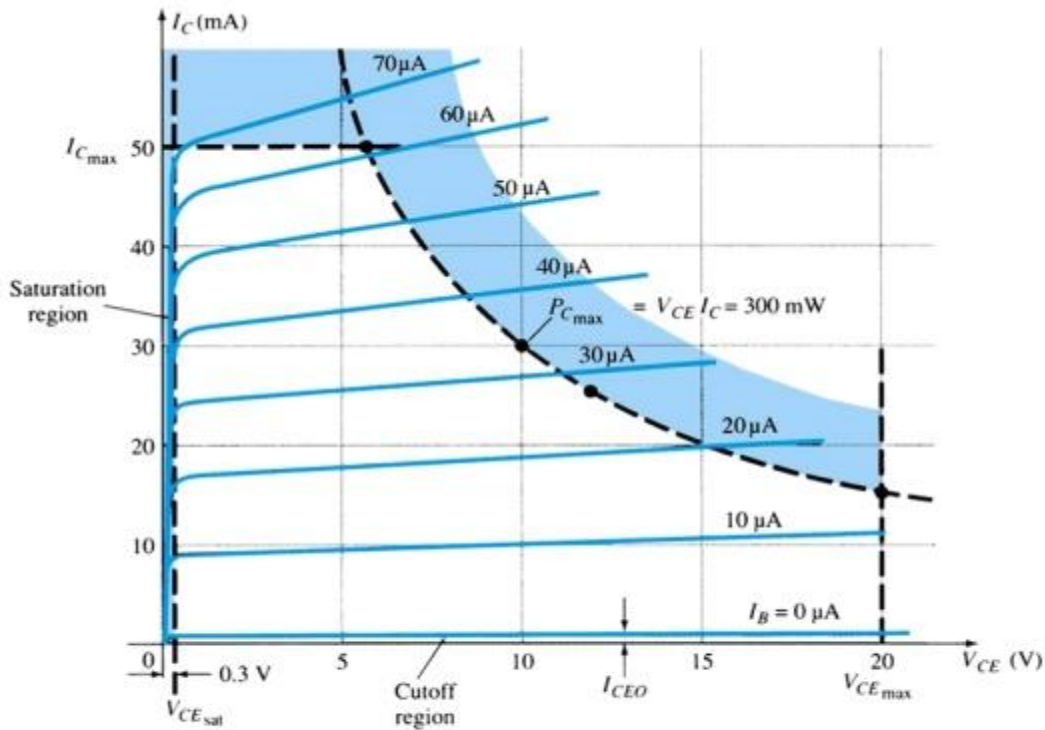
- For small V_{CE} ($V_{CE} < V_{CE(sat)}$, I_C increase linearly with increasing of V_{CE})
- $V_{CE} > V_{CE(sat)}$ I_C not totally depends on V_{CE} --> constant I_C
- I_B (μA) is very small compare to I_C (mA). Small increase in I_B cause big increase in I_C
- $I_B = 0 \text{ A}$ --> I_{CEO} occur.
- Noticing the value when $I_C = 0 \text{ A}$. There is still some value of current flows.

7. Describe the Common Collector configuration of transistor?

- Also called emitter-follower (EF).
- It is called common-emitter configuration since both the signal source and the load share the collector terminal as a common connection point.
- The output voltage is obtained at emitter terminal.
- The input characteristic of common-collector configuration is similar with common-emitter. configuration.
- Common-collector circuit configuration is provided with the load resistor connected from emitter to ground.
- It is used primarily for impedance-matching purpose since it has high input impedance and low output impedance.



- For the common-collector configuration, the output characteristics are a plot of I_E vs V_{CE} for a range of values of I_B .



Input characteristics of CC configuration

8. Explain how the Small Signal Amplifier has functioned in clear manner?

When the input signal is so weak as to produce small fluctuations in the collector current compared to its quiescent value, the amplifier is known as Small Signal Amplifier.

In other words, as the name indicates, the input applied to the circuit is $V_{in} \ll V_{th}$. It has only one amplifying device.

$$A = I_C / I_E$$

$$I_C = \alpha I_E + I_{CBO}$$

Voltage and current equation for hybrid parameters:

$$V_1 = h_{11}i_1 + h_{12}V_2$$

$$I_2 = h_{21}i_1 + h_{22}V_2$$

The values of h-parameters:

$$h_{11} = V_1 / i_1$$

$$h_{12} = V_1 / V_2$$

$$h_{21} = i_2 / i_1$$

$$h_{22} = i_2 / V_2$$

9. Explain the method of ADC in detail manner?

Successive Approximation Technique

When unknown voltage (V_a) is applied, the circuit starts up from 0000, as shown above. The output of SAR advances with each MSB. The output of SAR does not increase step-by-step in BCD bus pattern, but individual bit becomes high—starting from MSB. Then by comparison, the bit is fixed or removed. Thus, it sets first MSB (1000), then the second MSB (0100) and so on. Every time, the output of SAR is converted to equivalent analog voltage by binary ladder.

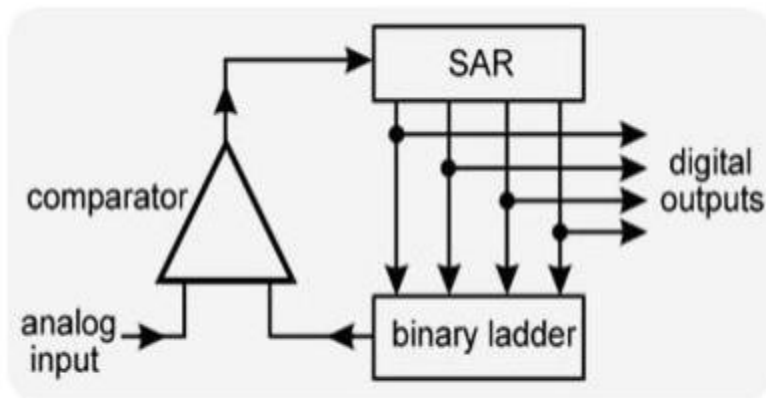


Figure: Successive Approximation Technique

Flash ADC

Also called the parallel A/D converter, this circuit is the simplest to understand. It is formed of a series of comparators, each one comparing the input signal to a unique reference voltage. The comparator outputs connect to the inputs of a priority encoder circuit, which then produces a binary output.

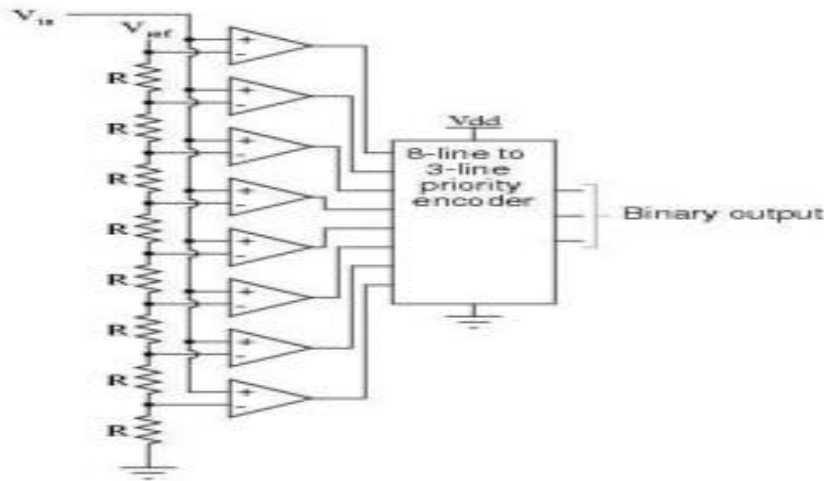


Figure: Flash ADC

10. Explain the method of DAC in detail manner?

Weighted resistors D/A converter

There are four resistors R , $2R$, $4R$ and $8R$ at the input terminals of the OPAMP with R as feedback resistor. The network of resistors at the input terminal of OPAMP is called as variable resistor network. The four inputs of the circuit are D, C, B & A. Input D is at MSB and A is at LSB. Here we shall connect 8V DC voltage as logic-1 level. So we shall assume that $0 = 0V$ and $1 = 8V$.

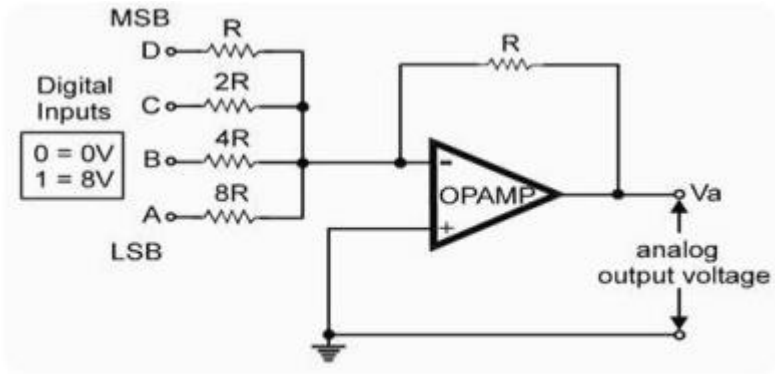


Figure: Weighted resistors D/A converter

Figure: Weighted resistors D/A converter

Now the working of the circuit is as follows. Since the circuit is summing amplifier, its output is given by the following equation

$$v_0 = R \left(\frac{D}{R} + \frac{C}{2R} + \frac{B}{4R} + \frac{A}{8R} \right)$$

Working of the circuit

When input DCBA = 0000, then putting these value in above equation (1) we get

$$v_0 = R \left(\frac{0}{R} + \frac{0}{2R} + \frac{0}{4R} + \frac{0}{8R} \right) = 0V$$

When digital input of the circuit DCBA = 0001, then putting these value in above equation (1) we get

$$v_0 = R \left(\frac{0}{R} + \frac{0}{2R} + \frac{0}{4R} + \frac{0}{8R} \right) = 0V$$

When digital input of the circuit DCBA = 0010, then putting these value in above equation (1) we get

$$v_0 = R \left(\frac{0}{R} + \frac{0}{2R} + \frac{8V}{4R} + \frac{0}{8R} \right) = -R \frac{8V}{4R} = -2V$$

R-2R Ladder D/A Converter

It is modern type of resistor network. It has only two values of resistors the R and 2R. These values repeat throughout in the circuit. The OPAMP is used at output for scaling the output voltage. The working of the circuit can be understood as follows. For simplicity, we ignore the OPAMP in the above circuit (this is because its gain is unity). Now consider the circuit, without OPAMP. Suppose the digital input is DCBA = 1000. Then the circuit is reduced to a small circuit.

$$\text{output} = \left(\frac{2R}{2R + 2R} \right) \times (+V) = \frac{V}{2}$$

Now suppose digital input of the same circuit is changed to DCBA = 0100. Then the output voltage will be V/4, when DCBA = 0010, output voltage will be V/8, for DCBA = 0001, output voltage will be V/16 and so on. The general formula for the above circuit of R–2R ladder, including the OPAMP also, will be –

$$v_0 = -R \left(\frac{D}{2R} + \frac{C}{4R} + \frac{B}{8R} + \frac{A}{16R} \right)$$

You can take (R) common from the above formula and simplify it. With the help of this formula, we can calculate any combination of digital input into its equivalent analog voltage at the output terminals.

UNIT V ELECTRICAL MEASUREMENTS

PART A

1. What are the basic elements of a generalized measurement system?

- Primary sensing element which is generally a transducer. Data conditioning element which further consists of variable conversion element and variable manipulation element.
- Data transmission and presentation elements which include data transmission system and data display system.

2. List any four Static characteristics of a measuring system.

Accuracy,
Precision,
Error,
Resolution,
Stability,
Linearity etc.

3. Define the term Accuracy.

The accuracy is defined as the degree of closeness with which the instrument reading approaches the true value of the quantity to be measured. It indicates the ability of an instrument to indicate the true value of the quantity.

4. What is an Error?

The algebraic difference between the indicated value and the true value of the quantity to be measured is called an error.

5. What is calibration?

Calibration is the process of making an adjustment or making a scale so that the readings of an instrument agree with the accepted value and the certified standard.

6. Define the term Precision.

It is the Measure of consistency or repeatability of measurements. It denotes the amount by which the individual readings are departed about the average of number of readings.

7. Give the applications of measurement systems.

- The instruments and measurement systems are used for
- Monitoring of processes and operations.
- Control of processes and operations.
- Experimental engineering analysis.

8. What are the various Dynamic characteristics?

- Various Dynamic characteristics are
- Fidelity
- Speed of Response
- Time Delay
- Lag

➤ Dynamic error

9. Explain the function of measurement system.

The measurement system consists of a transducing element which converts the quantity to be measured in an analogous form. The analogous signal is then processed by some intermediate means and is then fed to the end device which presents the results of the measurement.

10. Define static characteristics.

Static characteristics are the set of rules or criteria that is defined for those instruments that varies very slowly with time or remains a constant.

PART B

1.Explain the Static characteristics of measurement?

The set of criteria defined for the instruments, which are used to measure the quantities which are slowly varying with time or mostly constant, i.e., do not vary with time, is called 'static characteristics'.

Accuracy:

It is the degree of closeness with which the reading approaches the true value of the quantity to be measured. The accuracy can be expressed in following ways:

Conformity:

Consider a resistor having true value as 2385692 , which is being measured by an ohmmeter. But the reader can read consistently, a value as 2.4 M due to the non availability of proper scale. The error created due to the limitation of the scale reading is a precision error.

Number of significant figures:

The precision of the measurement is obtained from the number of significant figures, in which the reading is expressed. The significant figures convey the

actual information about the magnitude & the measurement precision of the quantity. The precision can be mathematically expressed as:

$$P=1-\frac{\overline{X_n-X_n}}{\overline{X_n}}$$

Where, P = precision

X_n = Value of nth measurement

$\overline{X_n}$ = Average value the set of measurement values

Sensitivity:

The sensitivity denotes the smallest change in the measured variable to which the instrument responds. It is defined as the ratio of the changes in the output of an instrument to a change in the value of the quantity to be measured.

Reproducibility:

It is the degree of closeness with which a given value may be repeatedly measured. It is specified in terms of scale readings over a given period of time.

Repeatability:

It is defined as the variation of scale reading & random in nature Drift: Drift may be classified into three categories:

a) zero drift:

If the whole calibration gradually shifts due to slippage, permanent set, or due to undue warming up of electronic tube circuits, zero drift sets in.

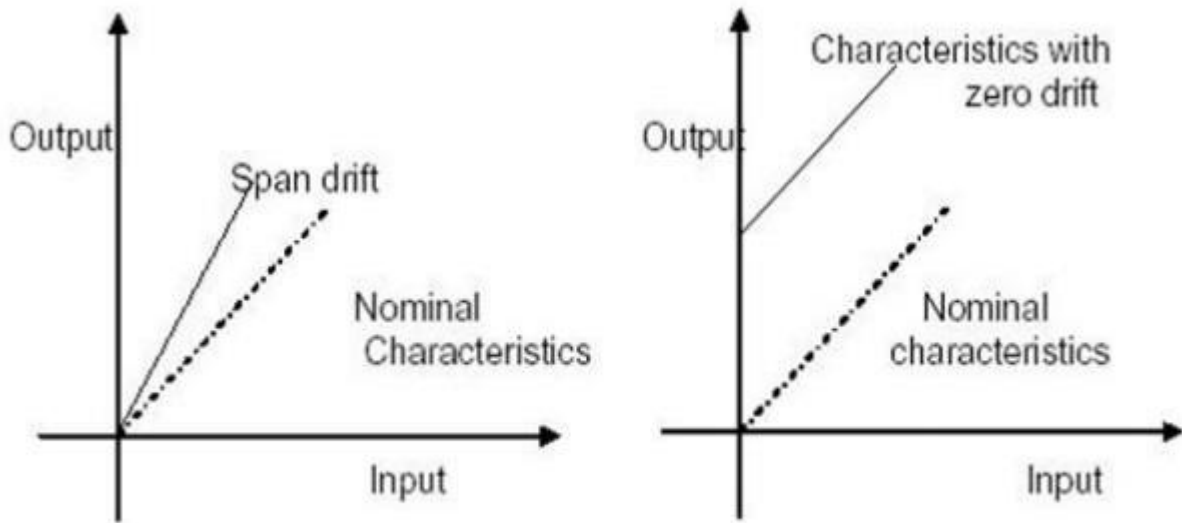


Fig. 4.1 Span Drift and Zero Drift

b) Span drift or sensitivity drift

If there is proportional change in the indication all along the upward scale, the drift is called span drift or sensitivity drift.

c) Zonal drift:

In case the drift occurs only a portion of span of an instrument, it is called zonal drift.

Resolution:

If the input is slowly increased from some arbitrary input value, it will again be found that output does not change at all until a certain increment is exceeded. This increment is called resolution.

Threshold:

If the instrument input is increased very gradually from zero there will be some minimum value below which no output change can be detected. This minimum value defines the threshold of the instrument.

Stability:

It is the ability of an instrument to retain its performance throughout its specified operating life.

Tolerance:

The maximum allowable error in the measurement is specified in terms of some value which is called tolerance.

Range or span:

The minimum & maximum values of a quantity for which an instrument is designed to measure is called its range or span.

2.Explain the Dynamic characteristics of measurement?

The set of criteria defined for the instruments, which are changes rapidly with time, is called 'dynamic characteristics'.

The various static characteristics are:

Speed of response

Measuring lag

Fidelity

Dynamic error

Speed of response:

It is defined as the rapidity with which a measurement system responds to changes in the measured quantity.

Measuring lag:

It is the retardation or delay in the response of a measurement system to changes in the measured quantity. The measuring lags are of two types:

a) Retardation type:

In this case the response of the measurement system begins immediately after the change in measured quantity has occurred.

b) Time delay lag:

In this case the response of the measurement system begins after a dead time after the application of the input. Fidelity: It is defined as the degree to which a

measurement system indicates changes in the measurand quantity without dynamic error.

3.Explain the types of measurement errors in detailed manner?

Errors in Measurement

1.Theoretical errors:

The explicit or implicit model on which we base our interpretation of our measurements may be inapplicable or inaccurate. Range of Validity: A model is applicable only within a limited range of Conditions. Beyond that, it will give inaccurate predictions. Approximation: Models have finite precision even within their range of validity

2.Static errors:

i. Parallax: Analog meters use a needle as a pointer to indicate the measured value. Reading this at an oblique angle causes a misreading, known as a parallax reading error.

ii. Interpolation: The needle often rests between two calibrated marks. Guessing its position by interpolation is subject to an error that depends on the size of the scale, and on the visual acuity and experience of the person reading the meter.

iii. Last-digit bobble: Digital readouts are often observed to oscillate between two neighboring values, for example a digital voltmeter (DVM) may alternately show 3.455 and 3.456 volts. This occurs when the actual value is about midway between the two displayed values. Small variations in the system under test, or in the meter itself, are sufficient to change the reading when it is delicately poised between the two values.

3.Environmental errors:

Measurements can be affected by change in ambient factors

I. Temperature

II. Pressure

III. Electromagnetic field: Static electric or magnetic fields, dynamic (changing) fields, and propagating fields (radiation) can interfere with measurements.

4. Characteristic errors:

Zero Offset: a meter (for example) may read zero when the actual value is nonzero. This is a common form of calibration error.

Gain error: amplifiers are widely used in instruments such as CRO probes, and we may trust that “times 10” means precisely what it says only when the amplification has been carefully calibrated.

Processing error: modern instruments contain complex processing devices such as analog computers which can introduce errors into the process leading to the displayed value of a measurement.

Repeatability error: instruments change over time, which is why they must be regularly calibrated, just as a car must be serviced. Instruments change, however slightly, even between consecutive measurements.

Nonlinearity: Drive an amplifier to too high a gain and it will operate in its nonlinear regions, producing a severely distorted output signal.

Resolution: devices can only resolve (that is, distinguish) values that are sufficiently separated. For example, optical instruments cannot easily resolve objects less than one wavelength apart.

5. Drift:

Drift is a complex phenomenon for which the observed effects are that the sensitivity and offset values vary. It also can alter the accuracy of the instrument differently at the various amplitudes of the signal present.

4. Describe the working principle of strain gauges in detailed manner?

Principle of Working of Strain Gauges

When force is applied to any metallic wire its length increases due to the strain. The more is the applied force, more is the strain and more is the increase in length of the wire. If L_1 is the initial length of the wire and L_2 is the final length after application of the force, the strain is given as: $\epsilon = (L_2 - L_1) / L_1$ Further, as the length of the stretched wire increases, its diameter decreases

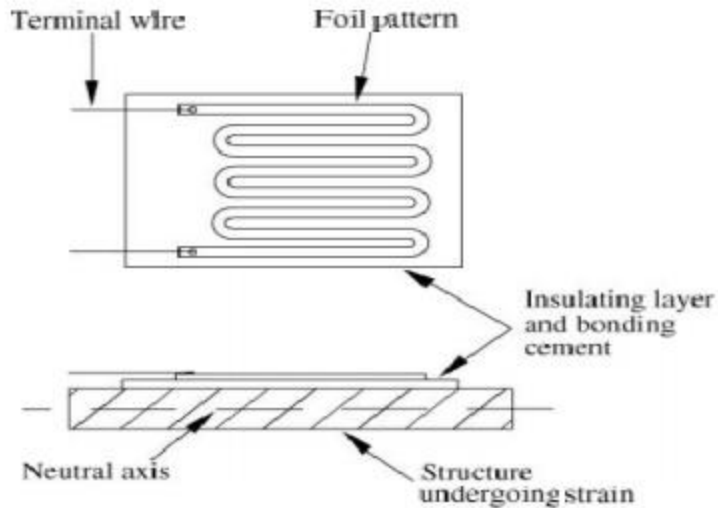


Fig. 4.5 Strain Gauge

Now, we know that resistance of the conductor is the inverse function of the length. As the length of the conductor increases its resistance decreases. This change in resistance of the conductor can be measured easily and calibrated against the applied force. Thus strain gauges can be used to measure force and related parameters like displacement and stress. The input and output relationship of the strain gauges can be expressed by the term gauge factor or gauge gradient, which is defined as the change in resistance R for the given value of applied strain ϵ .

Applications of the Strain Gauges

- 1) Measurement of strain
- 2) Measurement of other quantities

5. Describe the working principle of LVDT in detailed manner?

LVDT (Linear Variable Differential Transformer)

The LVDT operation does not require electrical contact between the moving part (probe or core rod assembly) and the transformer, but rather relies on electromagnetic coupling; this and the fact that they operate without any built-in electronic circuitry are the primary reasons why LVDTs have been widely used in

applications where long life and high reliability under severe environments are a required, such Military/Aerospace applications.

The LVDT consists of a primary coil (of magnet wire) wound over the whole length of a non-ferromagnetic bore liner (or spool tube) or a cylindrical coil form. Two secondary coils are wound on top of the primary coil for “long stroke” LVDTs (i.e. for actuator main RAM) or each side of the primary coil for “Short stroke” LVDTs (i.e. for electro-hydraulic servo-valve or EHSV). The two secondary windings are typically connected in “opposite series” (or wound in opposite rotational directions). A ferromagnetic core, which length is a fraction of the bore liner length, magnetically couples the primary to the secondary winding turns that are located above the length of the core.

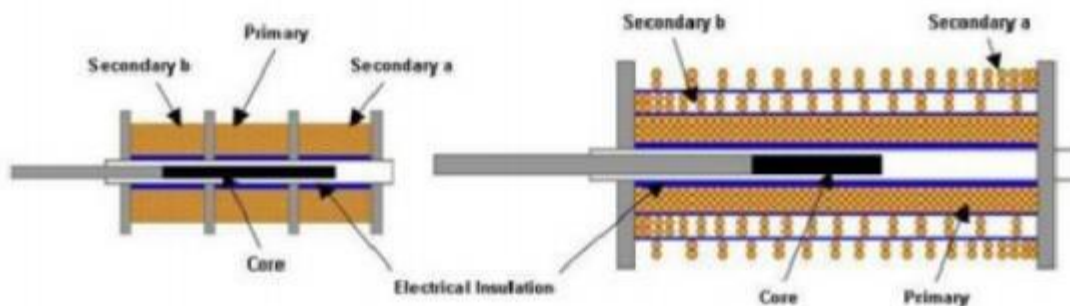
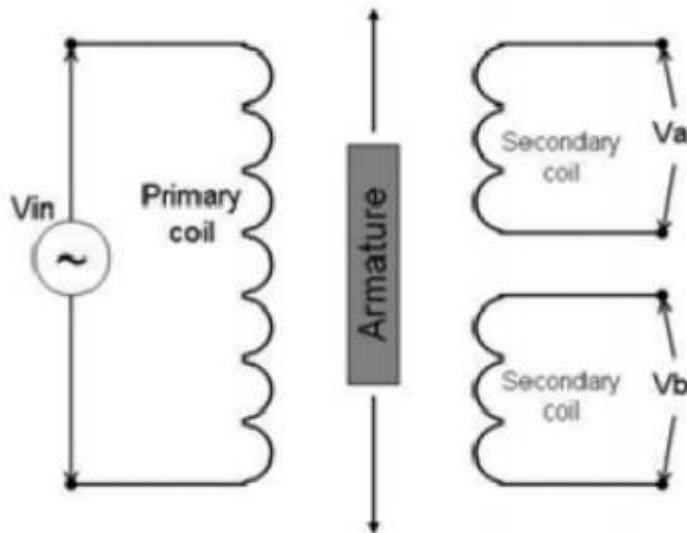


Fig. 4.9 LVDT

When the primary coil is excited with a sine wave voltage (V_{in}), it generate a variable magnetic field which, concentrated by the core, induces the secondary

voltages (also sine waves). While the secondary windings are designed so that the differential output voltage ($V_a - V_b$) is proportional to the core position from null, the phase angle (close to 0 degree or close to 180 degrees depending of direction) determines the direction away from the mechanical zero. The zero is defined as the core position where the phase angle of the ($V_a - V_b$) differential output is 90 degrees.

The differential output between the two secondary outputs ($V_a - V_b$) when the core is at the mechanical zero (or “Null Position”) is called the Null Voltage; as the phase angle at null position is 90 degrees, the Null Voltage is a “quadrature” voltage. This residual voltage is due to the complex nature of the LVDT electrical model, which includes the parasitic capacitances of the windings

6. Describe the working principle of RVDT in detailed manner?

Rotary Variable Differential Transformer (RVDT)

Rotary Variable Differential Transformer (RVDT) A Rotary Variable Differential Transformer (RVDT) is an electromechanical transducer that provides a variable alternating current (AC) output voltage that is linearly proportional to the angular displacement of its input shaft.

When energized with a fixed AC source, the output signal is linear within a specified range over the angular displacement. RVDT's utilize brushless, non-contacting technology to ensure long-life and reliable, repeatable position sensing with infinite resolution. Such reliable and repeatable performance assures accurate position sensing under the most extreme operating conditions. Moog offers seven frequency optimized RVDT's in a basic size 8 configured housing. Each is designed to operate at a specific frequency. Frequency optimization provides the benefit of an increased operating range of angular displacement with a reduction in sensor size and weight.

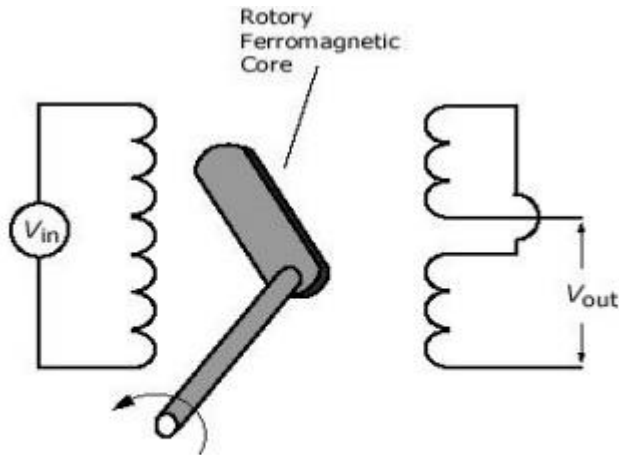


Fig. 4.10 RVDT

The Rotational Variable Differential Transformer (RVDT) is used to measure rotational angles and operates under the same principles as the LVDT sensor. Whereas the LVDT uses a cylindrical iron core, the RVDT uses a rotary ferromagnetic core.

7. Describe the working principle of Piezoelectric Transducer

in detailed manner?

The main principle of a piezoelectric transducer is that a force, when applied on the quartz crystal, produces electric charges on the crystal surface. The charge thus produced can be called as piezoelectricity. Piezo electricity can be defined as the electrical polarization produced by mechanical strain on certain class of crystals. The rate of charge produced will be proportional to the rate of change of force applied as input. As the charge produced is very small, a charge amplifier is needed so as to produce an output voltage big enough to be measured. The device is also known to be mechanically stiff. For example, if a force of 15 KN is given to the transducer, it may only deflect to a maximum of 0.002mm. But the output response may be as high as 100 KHz. This proves that the device is best applicable for dynamic measurement.

Piezoelectric Transducer can measure *pressure* in the same way a *force* or an *acceleration* can be measured. For low pressure measurement, possible vibration of the amount should be compensated for. The pressure measuring quartz disc stack faces the pressure through a *diaphragm* and on the other side of this stack, the compensating mass followed by compensating quartz.

Advantages

- Very high frequency response.
- Self generating, so no need of external source.
- Simple to use as they have small dimensions and large measuring range.

8.Describe the working principle of Energy Meter in detailed manner?

Introduction

The energy meter is an electrical measuring device, which is used to record Electrical Energy Consumed over a specified period of time in terms of units. Electric meters are typically calibrated in billing units, the most common one being the kilowatt hour.

A periodic reading of electric meters establishes billing cycles and energy used during a cycle.

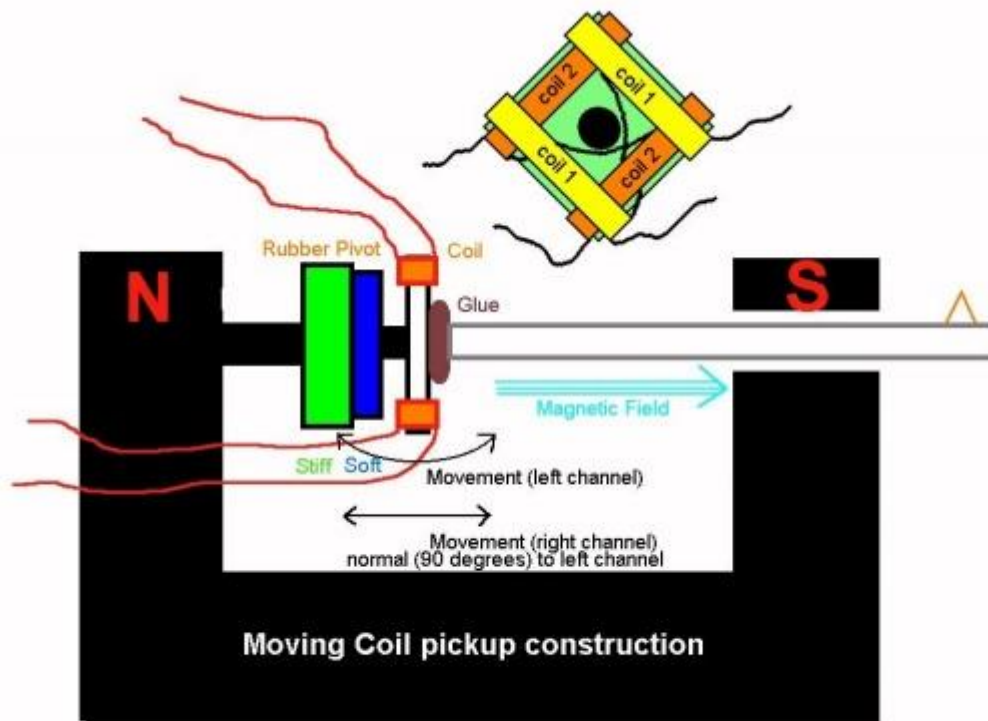
Features:

- * Display of current time (24 hours type), week, load power and cost tariff.
- * Display of total on time, total used energy and accrued energy cost.
- * Display of total record time, total on time and percentage.
- * Dual programmable power tariffs.
- * Connection, operation settings.

9. Describe the working principle of moving coil in detailed manner?

Moving Coil

Moving Coil Instruments are used for measuring DC quantities. They can be used on AC systems when fed through bridge rectifiers. Center magnet system is incorporated in our moving coil instruments which completely shields the movement from the effect of external magnetic fields. The movement is pivoted between synthetic sapphire jewel bearings for frictionless operation.



10. Describe the working principle of moving iron in detailed manner?

Moving Iron

Moving Iron Instruments are generally used for measuring AC Voltage and Currents. A feature of the moving element is that it is fitted with synthetic sapphire jewels. The movement is light, quick acting, but extremely robust. An efficient system of fluid damping is employed. The movement is efficiently shielded against the effect of external magnetic fields.

