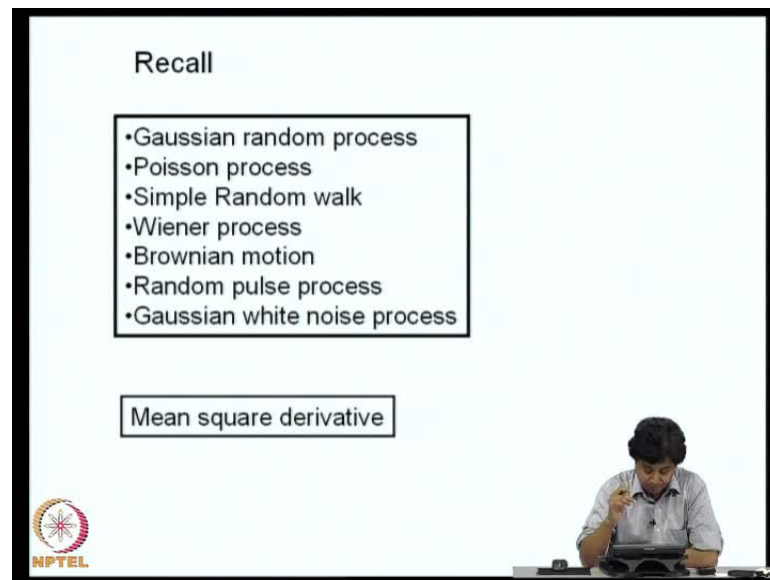


Stochastic Structural Dynamics
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Lecture No. # 09
Random Processes-4
Random Vibrations of sdof Systems-1

(Refer Slide Time: 00:27)



We have been reviewing topics in probability and random process. This would be the last lecture **at the part of this review, we will quickly recall** what we did in the previous lecture. We discuss properties of Gaussian random process, Poisson process, simple random walk, wiener process, Brownian motion, random pulse process and Gaussian white noise process.

So, will have opportunity to return to some of these processes, later when we use them to model, loads and responses, we also introduced the notion of mean square derivative; that was based on mean square convergence of random variables, in the sense of mean square, I briefly touched upon the topic of 2 random processes. In the previous lecture will continue with that.

(Refer Slide Time: 01:10)

Two random processes: $U(t)$ and $V(t)$


Description of $U(t)$

- $P[U(t) \leq u_1]$
- $P[U(t_1) \leq u_1 \cap U(t_2) \leq u_2]$
- ⋮
- $P\left[\bigcap_{i=1}^n U(t_i) \leq u_i\right]$

$p_U(\tilde{u}; \tilde{t})$

$$m_U(t) = \int_{-\infty}^{\infty} u p_U(u, t) du$$

$$C_{UU}(t_1, t_2) = \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} [u_1 - m_U(t_1)][u_2 - m_U(t_2)] p_{UU}(u_1, u_2; t_1, t_2) du_1 du_2$$



So, let us consider two random processes: U of t and V of t the description of U of t to recapitulate is in terms of first order probability distribution function or second order probability distribution function or in general the n th order probability distribution function. Associated with this we also have the n th order probability density function.

(Refer Slide Time: 02:07)


Description of $V(t)$

- $P[V(t) \leq v_1]$
- $P[V(t_1) \leq v_1 \cap V(t_2) \leq v_2]$
- ⋮
- $P\left[\bigcap_{i=1}^n V(t_i) \leq v_i\right]$

$p_V(\tilde{v}; \tilde{t})$

$$m_V(t) = \int_{-\infty}^{\infty} v p_V(v, t) dv$$

$$C_{VV}(t_1, t_2) = \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} [v_1 - m_V(t_1)][v_2 - m_V(t_2)] p_{VV}(v_1, v_2; t_1, t_2) dv_1 dv_2$$



We could define expectations this is the mean, this is covariance and so on and so for. So, this completes the description of a random process U of t and we say that U of t is completely specified, if I am able to specify the n th order joint density function for any

choice of these $t_1, t_2, t_3, \dots, t_n$ and for any choice of n . In a similar way **in** we can completely describe random process V of n th order density function or expectation, so on and so for.

(Refer Slide Time: 02:17)

Joint description of $U(t)$ and $V(t)$

- $P[U(t_1) \leq u_1 \cap V(t_2) \leq v_2]$
- $P\left[\bigcap_{i=1}^n U(t_i) \leq u_i \cap \bigcap_{j=1}^m V(s_j) \leq v_j\right]$

$p_{UV}(\vec{u}, \vec{v}; \vec{t}, \vec{s})$

$$C_{UV}(t_1, t_2) = \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} [u_1 - m_U(t_1)][v_2 - m_V(t_2)] p_{UV}(u_1, v_2; t_1, t_2) du_1 dv_2$$

$\langle U(t) \rangle = 0 \quad \langle V(t) \rangle = 0$

$C_{UV}(t_1, t_2) = \langle U(t_1)V(t_2) \rangle$

Cross covariance function

Now, when it comes to question of joint description of U of t and V of t , we can considered 2 time instance t equal to t_1 and t equal to t_2 and I select 1 random variable U of t_1 from ensemble of U of t and another random variable V of t_2 from ensemble of V of t . So, this can be viewed as the joint probability distribution function of the random process U of t and V of t at 2 time instance. So, this can be generalize for n time instance from ensemble of U of t and some other m time instance from ensemble of V of t and if I am able to specify, this kind of joint probability distribution function or the associated probability density function, we can say that the 2 random processes U of t and V of t are completely specified.

We can also define the moments with respect to these joint density functions. For instance, I can define, what is known as cross covariance function, which is actually the expectation of C of $U V$ of t_1, t_2 is expectation of suppose mean is 0, suppose if I assume U of t is 0 and V of t is 0, this is nothing but U of t_1 into V of t_2 , if mean is not 0 this is what I have shown here. This is expectation of u of t_1 minus m_U of t_1 v of t_2 minus m_V of t_2 into this density function. And this double integral gives the cross covariance function between U and V ; the word cross originates from the fact that the random

variable U of t_1 emanates from ensemble of U of t and the random variable V of t_2 emanates from ensemble of V of t .

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Joint stationarity

$U(t)$ and $V(t)$ are said to be jointly stationary (wide sense)

$$m_U(t) = m_U$$

$$m_V(t) = m_V$$

$$C_{UU}(t, t + \tau) = C_{UU}(\tau)$$

$$C_{VV}(t, t + \tau) = C_{VV}(\tau)$$

$$C_{UV}(t, t + \tau) = C_{UV}(\tau)$$

Definition for strong sense stationarity could also be provided in terms of joint pdf-s.

So, this is in contrast to the definition of auto covariance function, which we introduced previously, where I talked about 2 random variables originating from the same ensemble. So, the world although both functions are covariance, we add this prefix auto and cross to indicate the choice of the underline random variables. The notion of stationarity can be generalize to stationarity of two processes. Suppose, we ask the question, when do we say that U of t and V of t are said to be jointly stationary. First of all, U of t must be stationary, V of t must stationary, in the chosen sense of stationarity. Suppose, we are talking about wide sense stationarity, U of t should be wide sense stationary in it is own **right** V of t should be stationary in it is own **right**. That would mean of U would be a constant mean of V would be a constant and auto covariance of U of t that is C_{UU} of $t, t + \tau$ will be a function of only the time difference. Therefore C_{UU} of $t, t + \tau$ would be C_{UU} of τ .

(Refer Slide Time: 06:36)

Covariance matrix

$$C(t_1, t_2) = \begin{bmatrix} C_{UU}(t_1, t_2) & C_{UV}(t_1, t_2) \\ C_{VU}(t_1, t_2) & C_{VV}(t_1, t_2) \end{bmatrix}$$

$$C(\tau) = \begin{bmatrix} C_{UU}(\tau) & C_{UV}(\tau) \\ C_{VU}(\tau) & C_{VV}(\tau) \end{bmatrix}$$

$$C_{UV}(\tau) = \langle U(t)V(t+\tau) \rangle = \langle V(t+\tau)U(t) \rangle$$

$$\Rightarrow C_{UV}(\tau) = C_{VU}(-\tau)$$

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Handwritten diagram illustrating the relationship between covariance functions and time shifts for two processes, $U(t)$ and $V(t)$.

Left side (Process U):

- Waveform labeled $U(t)$.
- Time points $t_1 = t$ and $t_1 + \tau = t_2$ are marked with asterisks.
- Expression: $C_{UV} = \langle U(t_1) V(t_2) \rangle = C_{UV}(t_1, t_2)$

Right side (Process V):

- Waveform labeled $V(t)$.
- Time points $t_1 = t$ and $t_1 + \tau = t_2$ are marked with asterisks.
- Expressions:
 - $\equiv \langle V(t_2) U(t_1) \rangle$
 - $= \langle V(t_1 + \tau) U(t_1) \rangle$
 - $= C_{VU}(t_2, t_1)$

Bottom:

- Expression: $C_{VU}(t_1, t_2)$
- A red arrow points from the $C_{VU}(t_1, t_2)$ label up to the $C_{UV}(t_1, t_2)$ label, indicating the relationship $C_{UV}(\tau) = C_{VU}(-\tau)$.

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Similarly, C_{VV} of $t, t + \tau$ is C_{VV} of τ . The additional requirement for wide sense stationarity of two random processes is that the cross covariance function should also satisfy this equal that is C_{UV} of $t, t + \tau$ must be function of the time difference start in the same sense. We can also discuss about strong sense stationarity in terms of probability density functions. So, I take joint density function between a random variable originating from U of t and V of t at some time difference and based on the notion of stationarity that we consider for U of t we could extend the definitions and we can talk about strong sense stationarity of two random processes.

Now, let us look at what exactly is meant by, we talked about C_{UV} of is expectation of U of t_1 V of t_2 . So, I will call it as C_{UV} of t_1, t_2 ; suppose you have an ensemble, suppose this is ensemble of U of t and this is an ensemble of V of t , what we are saying is there is the time instant t , the another time instant t plus τ , there is this same time instant t and time instant t plus τ . So, there are basically four random variables here two from U of t and two from v of t . When I says C_{UV} of t_1, t_2 it is from U of t_1 that is, this is t_1 say this is t_1 is equal to t and this is a t_2 say t_1 equal to t this is equal to t ; U of t_1 that is this and V of t_2 that is this. This is clearly different from that covariance between these 2 random variables.

So, this is also equal to that is this. This also equal to expectation of V of t_2 into U of t_1 . So, this is V of t plus τ and U of t so this can be written as C_{VU} of t_2, t_1 . The point that I am making is this is not a symmetric function, because if you simply consider C_{VU} of t_1, t_2 then you will be looking at correlation between two other random variables. For example, if you now consider C_{VU} of t_1, t_2 . This is correlation between this random variable and this random variable. So, there is no reason why the co relation covariance between these two random variables should be same as covariance between this random variable. So, is not symmetric, because there are four random variables we are talking about.



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Covariance matrix

$$C(t_1, t_2) = \begin{bmatrix} C_{UU}(t_1, t_2) & C_{UV}(t_1, t_2) \\ C_{VU}(t_1, t_2) & C_{VV}(t_1, t_2) \end{bmatrix}$$

$$C(\tau) = \begin{bmatrix} C_{UU}(\tau) & C_{UV}(\tau) \\ C_{VU}(\tau) & C_{VV}(\tau) \end{bmatrix}$$

$$C_{UV}(\tau) = \langle U(t)V(t+\tau) \rangle = \langle V(t+\tau)U(t) \rangle$$

$$\Rightarrow C_{UV}(\tau) = C_{VU}(-\tau)$$



So, we can write this information covariance in the metric form C of t_1, t_2 C_{UU} of t_1, t_2 is the auto covariance of U of t_1, t_2 C_{UV} of t_1, t_2 is cross covariance between U of t_1 and V of t_2 C_{VU} of t_1, t_2 is cross covariance between V of t_1 and U of t_2 and this is C_{VV} of t_1, t_2 is auto covariance of V of t_1 and V of t_2 . This is not symmetric, because we are not talking about covariance between the same two random variables. There are four random variables. So, although this is a covariance matrix it should be understood that it is not in the sense of covariance between two random variables there are actually four random variables here that is why this asymmetric arises C of τ is this, if this process is stationary this will be suppose t is t_1 is t and t_2 is t plus τ this is the covariance matrix. So, this is what I am trying to say here C_{UV} of τ is U of t into V of t plus τ . So, this is commutative, I can write this first and this next this is V of t plus τ U of t and this is nothing but C_{VU} of t minus t minus τ . Therefore, C_{UV} of τ is C_{VU} of minus τ . So, in some sense there is the kind of an anti-symmetric, we can see here. So, this would be understood, what is the corresponding definition of power spectral density functions?

(Refer Slide Time: 11:10)


PSD matrix

$$S(\omega) = \begin{bmatrix} S_{UU}(\omega) & S_{UV}(\omega) \\ S_{VU}(\omega) & S_{VV}(\omega) \end{bmatrix}$$

$$S_{UU}(\omega) = \lim_{T \rightarrow \infty} \frac{1}{T} \langle \overset{U_T}{X_T}(\omega) \overset{U_T^*}{X_T^*}(\omega) \rangle \Rightarrow S_{UU}(\omega) = S_{UU}(-\omega)$$

$$S_{UV}(\omega) = \lim_{T \rightarrow \infty} \frac{1}{T} \langle U_T(\omega) V_T^*(\omega) \rangle \Rightarrow$$

$$S_{UV}(-\omega) = \lim_{T \rightarrow \infty} \frac{1}{T} \langle U_T(-\omega) V_T^*(-\omega) \rangle = \lim_{T \rightarrow \infty} \frac{1}{T} \langle U_T^*(\omega) V_T(\omega) \rangle = S_{VU}(\omega)$$

$$S_{VU}^*(\omega) = \lim_{T \rightarrow \infty} \frac{1}{T} \langle U_T^*(\omega) V_T(\omega) \rangle = \lim_{T \rightarrow \infty} \frac{1}{T} \langle U_T(-\omega) V_T^*(-\omega) \rangle = S_{UV}(-\omega)$$


So, we define power spectral density matrix of power spectral density functions S_{UU} of ω is the auto power spectral density function of process U of T U of T is taken to be stationary. And do you will also assume that it has 0 mean, S_{UV} of ω is known as cross power spectral density function between U and V , this is S_{VU} of U this S_{VV} of ω . So, you recall how did we define S_{UU} of ω , it was limit of 1 by T

expectation of X^T of ω into it is conjugate. Now, if you replace ω by minus ω what happens this will be X^T of minus ω is X^T star of minus ω X^T of minus ω is nothing but X^T star of ω and X^T star of minus ω is nothing but X^T of ω . Therefore, S_{UU} of ω is S_{UU} of minus ω .

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$$S_{UV}(\omega) = \int_{-\infty}^{\infty} R_{UV}(\tau) \exp(i\omega\tau) d\tau$$

$$R_{UV}(\tau) = \frac{1}{2\pi} \int_{-\infty}^{\infty} S_{UV}(\omega) \exp(-i\omega\tau) d\omega$$

$$S_{UV}(\omega) = |S_{UV}(\omega)| \exp[-i\phi(\omega)]$$

$$|S_{UV}(\omega)| = \text{amplitude of cross PSD function}$$

$$\phi(\omega) = \text{phase spectrum}$$


$$p_{UV}(\omega) = \text{Re}[S_{UV}(\omega)] = \text{co-spectrum}$$

$$q_{UV}(\omega) = \text{Im}[S_{UV}(\omega)] = \text{quadrature spectrum}$$

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Now, if you consider S_{UV} of ω , this is expectation of U^T of ω V^T star of ω , I think here there is a mistake here this S_{UU} of ω this X should be U^T U of ω is defined like this now if I replace ω by minus ω so this is U^T of ω is U^T of minus ω V^T star of ω is V^T star of minus ω this is U^T star of ω and this is V^T of ω . Therefore, this is S_{VU} of ω now you take conjugate of S_{UV} of ω this will be defined as U^T star of ω V^T of ω , because you are taking conjugation of this this is nothing but S_{UV} of minus ω the point that you have to notice is S_{UV} of ω , which is the cross power spectral density function it is a complex valued function it has a nonzero real part and a nonzero imaginary part, you can show that this indeed is the Fourier transform of the cross covariance function, we were assuming mean to be 0. Therefore, there is no distinction between cross covariance and cross correlation function.

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Complex coherency function

$$\text{coh}_{UV}(\omega) = \frac{S_{UV}(\omega)}{\sqrt{S_{UU}(\omega)S_{VV}(\omega)}}$$
$$\text{coh}_{UV}(\omega) = |\text{coh}_{UV}(\omega)| \exp(-i\theta(\omega))$$

Coherency

$$|\text{coh}_{UV}(\omega)| = \frac{|S_{UV}(\omega)|}{\sqrt{S_{UU}(\omega)S_{VV}(\omega)}}$$
$$0 \leq |\text{coh}_{UV}(\omega)| \leq 1$$

$|\text{coh}_{UV}(\omega)| = 0$
=> lack of linear dependency between two processes

Two processes are linearly related

$$|\text{coh}_{UV}(\omega)| = 1$$

11

So, this is cross correlation because mean is 0. This is a Fourier transform pair constituted by S_{UV} of ω and R_{UV} of τ sense. This is not a symmetric function S_{UV} of ω will have a non-zero real part and imaginary part and we can write it in the polar form as an amplitude and a phase. So, this amplitude we called it as amplitude of cross PSD function and this is the phase spectrum. The real part of S_{UV} of ω is called co spectrum and imaginary part is known as quadrature spectrum. This is a nomenclature later on when we model earthquake loads using multiple random processes will be will have a occasion to use some of this nomenclature, we can introduce what is known as Complex coherency function, which is S_{UV} of ω divided by square root of S_{UU} of ω S_{VV} of ω this again is a complex function and the modulus of this complex function is known as coherency.

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Example

Let $U(t)$ and $V(t)$ be defined as

$$U(t) = S(t)$$

$$V(t) = S(t + \alpha)$$

$S(t)$ = stationary Gaussian random process with zero mean.


$$C_{UV}(\tau) = \langle U(t)V(t+\tau) \rangle = \langle S(t)S(t+\alpha+\tau) \rangle = C_{SS}(\alpha+\tau)$$

$$\Rightarrow S_{UV}(\omega) = \int_{-\infty}^{\infty} C_{UV}(\tau) \exp(-i\omega\tau) d\tau = \int_{-\infty}^{\infty} C_{SS}(\alpha+\tau) \exp(-i\omega\tau) d\tau$$

Substitute $\beta = \alpha + \tau \Rightarrow$

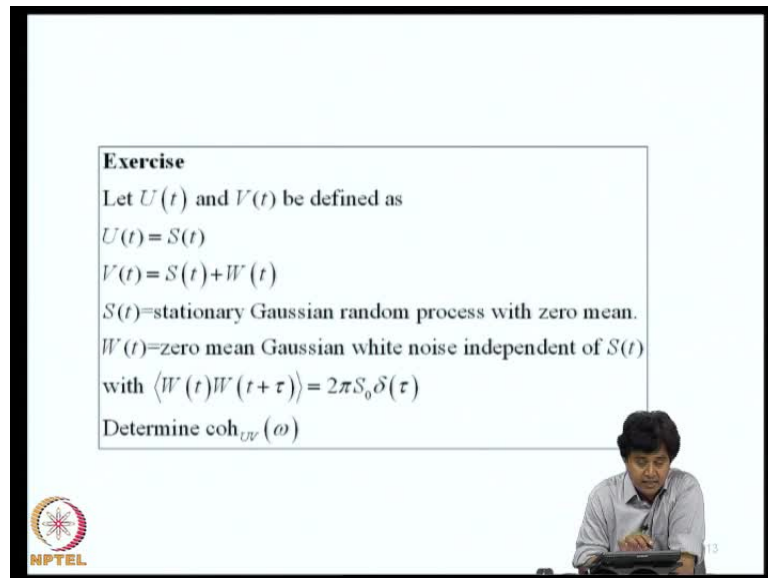
$$S_{UV}(\omega) = \int_{-\infty}^{\infty} C_{SS}(\beta) \exp(-i\omega\tau + i\omega\alpha) d\beta = \exp(i\omega\alpha) S_{SS}(\omega)$$

$$\Rightarrow \rho_{UV}(\omega) = \exp(i\omega\alpha)$$





This coherency you can show that takes value between 0 and 1 and if U and V are independent, you can show that coherency will be 0 and if they share a linear relationship between them, the coherency will be 1. So, this again is somewhat analogous to the cross correlation function that we were discussing earlier, we will consider a simple example, where I defined two random processes U of t and V of t in terms of a parent process S of t I call U of t as S of t and V of t is S of t plus α . Let S of t be a stationary Gaussian random process with 0 mean. So, what is the coherency or a cross power spectral density function between U and V that is the question C_{UV} of τ can be expressed as expectation of U of t into V of t plus τ mean of U and mean of V are 0 because mean of S and S is stationary and mean is 0. Therefore, expected value of U and V would be 0.

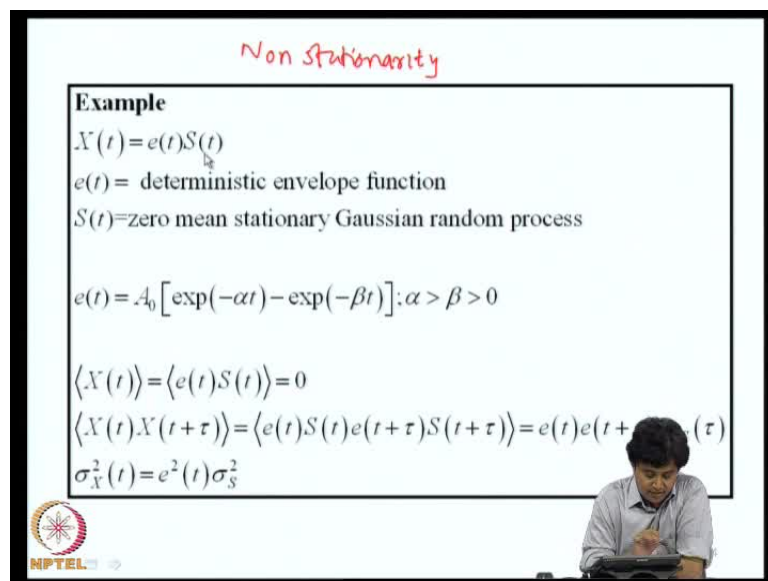
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Exercise
 Let $U(t)$ and $V(t)$ be defined as
 $U(t) = S(t)$
 $V(t) = S(t) + W(t)$
 $S(t)$ = stationary Gaussian random process with zero mean.
 $W(t)$ = zero mean Gaussian white noise independent of $S(t)$
 with $\langle W(t)W(t+\tau) \rangle = 2\pi S_0 \delta(\tau)$
 Determine $\text{coh}_{UV}(\omega)$



 

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Non stationarity

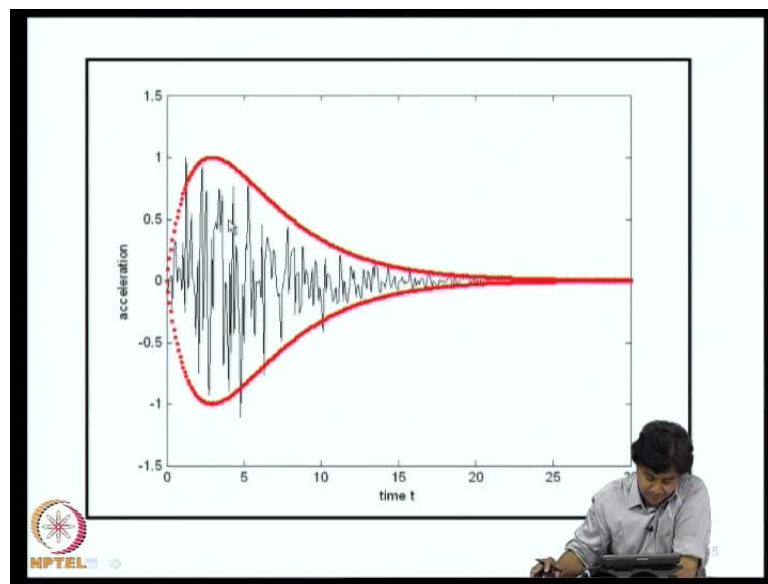
Example
 $X(t) = e(t)S(t)$
 $e(t)$ = deterministic envelope function
 $S(t)$ = zero mean stationary Gaussian random process
 $e(t) = A_0 [\exp(-\alpha t) - \exp(-\beta t)]$; $\alpha > \beta > 0$
 $\langle X(t) \rangle = \langle e(t)S(t) \rangle = 0$
 $\langle X(t)X(t+\tau) \rangle = \langle e(t)S(t)e(t+\tau)S(t+\tau) \rangle = e(t)e(t+\tau) \langle S(t)S(t+\tau) \rangle$
 $\sigma_X^2(t) = e^2(t)\sigma_S^2$

So, $U(t) = S(t)$ and $V(t) = S(t) + W(t)$. Therefore, this is $C_{UV}(\tau) = \langle U(t)V(t+\tau) \rangle = \langle S(t)[S(t+\tau) + W(t+\tau)] \rangle = \langle S(t)S(t+\tau) \rangle + \langle S(t)W(t+\tau) \rangle$. So, if you consider now the fourier transform of this, we get the cross power spectral density function between U and V , this is given by this $C_{UV}(\tau) = \int_{-\infty}^{\infty} e(t)e(t+\tau) \langle S(t)S(t+\tau) \rangle \delta(\tau - \tau') d\tau$. So, using this relation for C_{UV} here we get this. We can now make a substitution $\beta = \alpha + \tau$ and limits of integration would remain unchanged, but there will be $C_{UV}(\tau)$ will become $C_{SS}(\beta)$ and this τ will be $\beta - \alpha$ and if you use that we get $S_{UV}(\omega) = \int_{-\infty}^{\infty} C_{UV}(\tau) e^{-j\omega\tau} d\tau = \int_{-\infty}^{\infty} e(t)e(t+\tau) \langle S(t)S(t+\tau) \rangle \delta(\tau - \tau') e^{-j\omega\tau} d\tau$. This would mean

coherency U V of ω is exponential $I \omega \alpha$, we can consider as a similar problem again there is a parent process S of t and at another process W of t U of t is S of t and V of t is S of t plus W of t S of t is taken to be a stationary Gaussian random process with 0 mean and W of t is 0 mean Gaussian white noise process, which is independent of S of t . The question is what is coherency between U and V ? So, this is an exercise that you can try.

(Refer Slide Time: 18:43)



We now briefly touch upon, question of how to deal with non-stationary? One of the simple technique to deal with non stationarity is to express a non-stationary process x of t as a product of stationary process S of t and deterministic envelope function or modulating function e of t is a deterministic envelope function, S of t for purpose of illustration will take it as 0 mean stationary Gaussian random process, X of t would be Gaussian because this a linear transformation it is a constant in to a Gaussian quantity. So, the X of t would also be Gaussian. In modeling earth quake ground acceleration, this is one of the popular models that will be using we can see first how it looks like the black line that you are seeing here is X of t and this red line is actually the envelope e of t .

So, this envelope multiplies the sample of a stationary random process and what you get is this black line. So, I am talking about modeling of X of t where non stationarity as this kind of transient behavior, this is typically how an earthquake record might look like it will last for about 30 seconds and there is a strong motion phase and gentle d k. so, we

are interested in characterizing the behavior of the structure during the strong motion phase that is what typically we are interested in.

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Non stationarity

Example

$$X(t) = e(t)S(t)$$


$e(t)$ = deterministic envelope function
 $S(t)$ = zero mean stationary Gaussian random process

$$e(t) = A_0 [\exp(-\alpha t) - \exp(-\beta t)]; \alpha > \beta > 0$$

$$\langle X(t) \rangle = \langle e(t)S(t) \rangle = 0$$

$$\langle X(t)X(t+\tau) \rangle = \langle e(t)S(t)e(t+\tau)S(t+\tau) \rangle = e(t)e(t+\tau)\langle S(t)S(t+\tau) \rangle$$

$$\sigma_X^2(t) = e^2(t)\sigma_S^2$$



So, a brief description of this random process can be obtained by considering its mean and covariance function. Since X of t is Gaussian the mean and covariance function completely specify X of t . So, what is expected value of X of t ? Expected value of X of t is expected value of e of t into S of t , which is zero because expected value of S of t is 0. Then expected value of X of t into X of t plus τ is nothing but e of t into s of t e of t plus τ into S of t plus τ then this is e of t into e of t plus τ and $R S S$ of τ . So, the covariance of X of t now depends on t t plus τ this is a stationary part but X of t by virtual presents of these 2 terms continued to be non-stationary.

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Markov Property

Let $X(t)$ be a random process with continuous state and continuous parameter (time t).


Let $t_1 < t_2 < \dots < t_n$ be n time instants.

This defines n random variables $X(t_1), X(t_2), \dots, X(t_n)$.

$X(t)$ is said to possess Markov property if

$$P[X(t_n) \leq x_n \mid X(t_{n-1}) \leq x_{n-1}, X(t_{n-2}) \leq x_{n-2}, \dots, X(t_1) \leq x_1]$$
$$= P[X(t_n) \leq x_n \mid X(t_{n-1}) \leq x_{n-1}]$$

for any n and any choice of $t_1 < t_2 < \dots < t_n$.

36

If you consider for example, τ equal to 0 this expectation is nothing but a variance of X of t and you can see here that variance is e square of t into σ^2 that would mean the variance are the standard deviation is the function of time. So, it is not a white noise stationary process and since its process is Gaussian even the first order density function will be a changing function of time. Therefore, it is not stationary. There is one property which will be very useful in modeling that is known as Markov property. We use the term Markov random processes it should be noted that Markov property does not describe a distribution function for instants a Gaussian random process refers to the fact that at a given time t the random variable X of t is Gaussian distributed. So, you take two random variables there jointly Gaussian and so on and so forth.

When I talk about Markov processes, there is no underline probability distribution that I am talking about for example, a Gaussian random process could have a Markov property. So, what is this Markov property? Let us consider a random process X of t , let us assume that it has continuous state and the time parameter t is also continuous, will consider n time instants t_1, t_2 up to t_n ordered in this particular manner. This defines n random variables X of t_1, X of t_2 and so on and so forth. X of t_n to characterize a Markov property, we need the conditional probability distribution of X of t_1 less than or equal to x_n . So, I consider the probability that X of t_1 is less than or equal to x_n condition on the fact that X of t_{n-1} is less than x_{n-1}, X of t_{n-2} is less than or equal to t_{n-2} and so on and so forth.

If this is equal to X of t_n less than or equal to x_n condition only on the value of X of t_{n-1} , the previous time instant we say that the process is Markov. This should be true for any choice of n and any choice of t_1 and t_n that means, if you imagine that n th time instant defines a future and $n-1$ time is the present. What Markov property says is that future depends upon the present and not on the past. How the present was arrived at has no bearing, on what happens tomorrow. It has so called one step memory, you recall when I talked about white noise are a sequence of independent random variables, they have no memory, there is no memory there.

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⇒

$$P_X(x_n, t_n | x_{n-1}, t_{n-1}; x_{n-2}, t_{n-2}; \dots; x_1, t_1) = P_X(x_n, t_n | x_{n-1}, t_{n-1})$$

$$P_X(x_n, t_n | x_{n-1}, t_{n-1}; x_{n-2}, t_{n-2}; \dots; x_1, t_1) = P_X(x_n, t_n | x_{n-1}, t_{n-1})$$

Description of a Markov process

- $p(x_1, t_1)$
- $p(x_2, t_2; x_1, t_1) = p(x_2, t_2 | x_1, t_1) p(x_1, t_1)$
- $p(x_3, t_3; x_2, t_2; x_1, t_1) = p(x_3, t_3 | x_2, t_2; x_1, t_1) p(x_2, t_2 | x_1, t_1) p(x_1, t_1)$
 $= p(x_3, t_3 | x_2, t_2) p(x_2, t_2 | x_1, t_1) p(x_1, t_1)$
- \vdots
- $p(x_n, t_n; x_{n-1}, t_{n-1}; \dots; x_1, t_1) = \prod_{v=2}^n p(x_v, t_v | x_{v-1}, t_{v-1}) p(x_1, t_1)$

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Here, if there is one step memory, in terms of distribution function the Markov property described as the probability distribution function of at t_n that is x_n, t_n conditioned on $x_{n-1}, t_{n-1}, x_{n-2}, t_{n-2}$ so on and so forth is equal to probability distribution of x_n, t_n conditioned on x_{n-1}, t_{n-1} . The fact that all these has no bearing on this function, associated with this, we can define the n th order probability density function this also displace the same trait that the knowledge of what was happening from t_1 to t_{n-2} has no bearing on this function. So, again you can see there is a one-step memory.

So, how do we describe a Markov random process? We can start by talking about first order probability distribution function. There is 1 random variable x of t_1 that is completely describe in terms of a density function, two random variables, I will write it

in terms of joint density function or a conditional density function and the first order density function. How about third order probability density function? This can be define in terms of a third order joint density function or two conditional density functions and 1 first order density function, if x of t is Markov this density function conditional density function will be equal to this because what happens at t_3 depends only on what is happening a t_2 and not on what has happened t_1 . so this density function becomes independent of x_1 and t_1 .

So, here we need two conditional probability density functions and one probability density first order density function. There is you can see here, this describes the transition from t_2 to t_3 . This describe transition from t_1 from t_2 and this conditional probability density function is known as transitional probability density function of the Markov process x of t . So, if you consider now a n th order joint density function by using this argument, we can show that this depends on the first order density function P of $X_1 t_1$ and products of n minus 1 conditional density functions and these are known as I said conditional probability density functions.

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Transition probability density function

tpdf : $p(x_v, t_v | x_{v-1}, t_{v-1})$

• $p(x_1, t_1)$ and $p(x_v, t_v | x_{v-1}, t_{v-1})$
completely specify a Markov process

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Chapman - Kolmogorov - Smoluchowski Equation

$t = t_1$

$t = \tau$

$t = t_2$

$$\begin{aligned}
 p(x_2, t_2; x_1, t_1) &= \underline{p(x_2, t_2 | x_1, t_1)} p(x_1, t_1) \\
 &= \int p(x_2, t_2; x, \tau; x_1, t_1) dx \\
 &= \int \underline{p(x_2, t_2 | x, \tau; x_1, t_1)} p(x, \tau | x_1, t_1) p(x_1, t_1) dx \\
 \Rightarrow \\
 p(x_2, t_2 | x_1, t_1) &= \int p(x_2, t_2 | x, \tau; x_1, t_1) p(x, \tau | x_1, t_1) dx \\
 &= \int p(x_2, t_2 | x, \tau) p(x, \tau | x_1, t_1) dx
 \end{aligned}$$



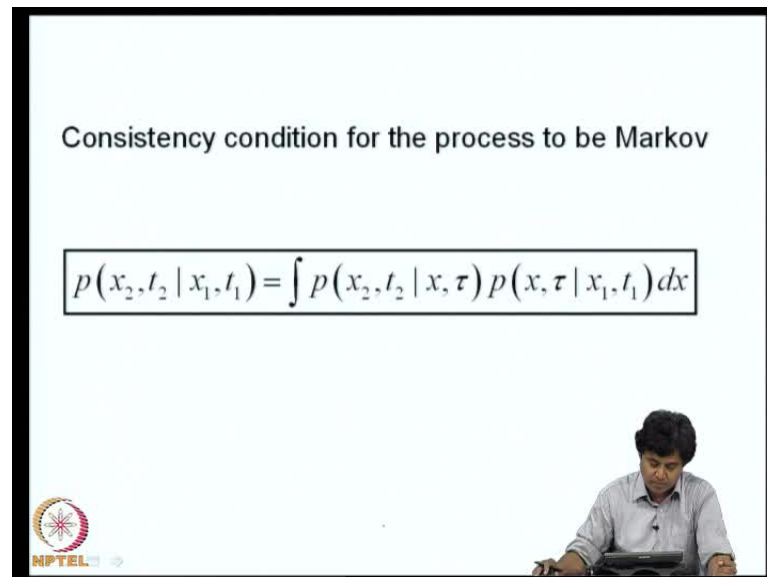
So, this is denote in this transition probability density function is called as t p d f. And this is given as p of x new, t new conditioned on x new minus one t new minus one. So, a complete description of a Markov process is through a first order density function and this transition probability density function. There is the further requirement on this transition density function which is known as a compatibility requirement or consistency requirement that can be explain as follows, you consider two time instants t equal to t 1 and t equal to t 2 and a intermediate time instant t equal to tau.

Now, let us consider the joint density function between x of t 1 and x of t 2 and that is given by x 2, t 2 x 1, t 1 this can be written as probability density function x 2, t 2 conditioned on x 1 t 1 into probability of x 1 t 1. This conditional density function **itself...** This joint density function itself can also be interrupted as a marginal density function of a third order probability density function that means, I consider now three random variables: one here, one here and one here and integrate with respect to states of this random variable. So, the second order density function is a marginal density function of a third order density function.

Now, let us use a conditional density is now. This I will to since t 1 tau and t 2 are order in the manner that we are talking probability of x 2, t 2 conditioned on x of tau and x 1 of t 1 into probability of x of tau conditioned on x 1 t 1 into probability of x 1 t 1 d x. Now, the process is Markov. Therefore, this density function can be approximated this density

function, this is same as this approximated by x_2, t_2 of τ . Now, see what is this equality this equality is between this term and this term there is p of x_1, t_1 on the left hand side and p of x_1, t_1 on the right hand side, they can be cancel. So, I am left with p of x_2, t_2 conditioned on x_1, t_1 is this integral.

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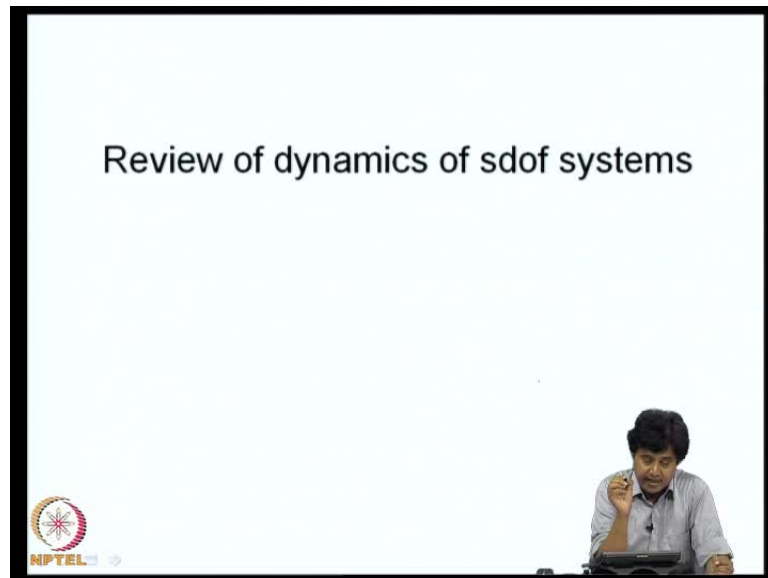
Consistency condition for the process to be Markov

$$p(x_2, t_2 | x_1, t_1) = \int p(x_2, t_2 | x, \tau) p(x, \tau | x_1, t_1) dx$$

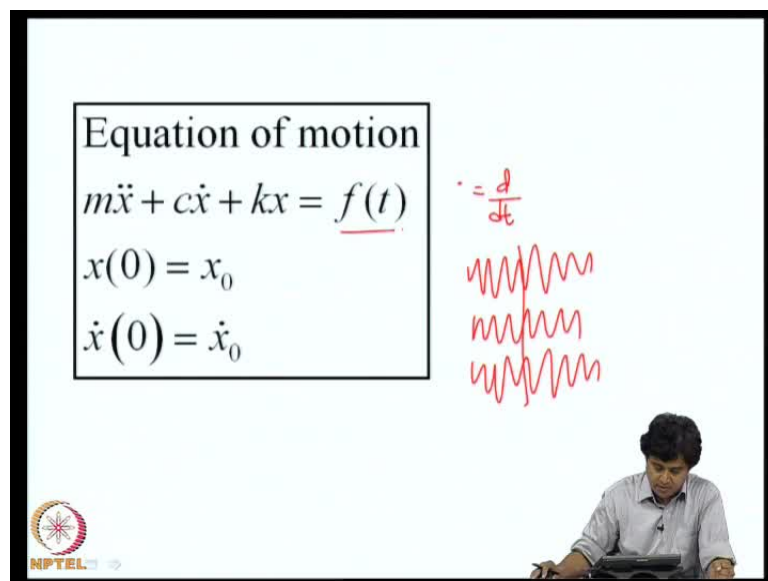
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Now, I have invoked the Gaussian property and this is the expression I get on the right hand side, what this says is basically in moving from t_2 to t_1 . The conditional density function should satisfy this internal requirement. So, this is the requirement for the transition probability density function of a Markov random process. So, this the consistency condition for the process to be Markov. So, this is so called Chapman Kolmogorov Smoluchowski relation, we will written to this related in this course, where we will use Markov property of responses to characterize the probability distribution of structural response to random loads that will come slightly later in the course.

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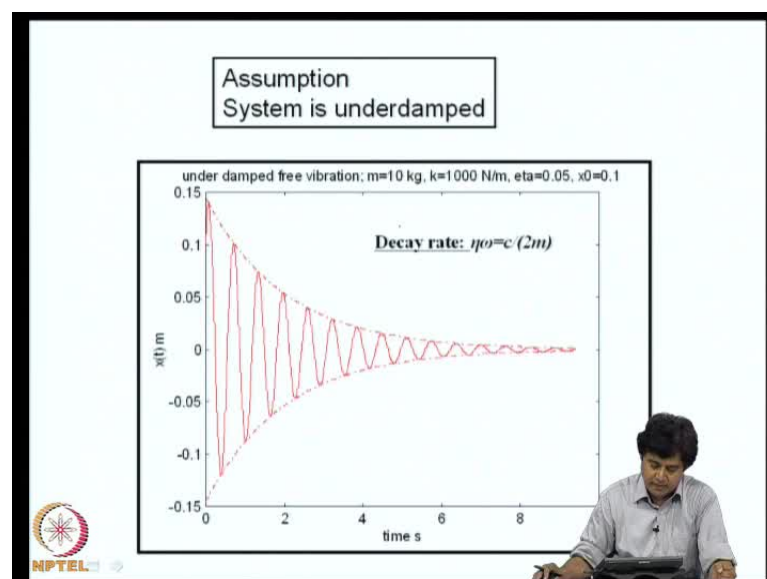


Now, with these we have completed the review of topics and probability and random processes. And we are now ready to start discussion discussing about stochastic structural dynamics or random vibration problems. As a prelude to that we will begin by quickly reviewing, the basic principles of dynamics of single degree freedom systems as this is a familiar equation of motion for a single degree freedom system m is the mass, c is a damping, k is a stiffness and f of t is the excitation, x naught and \dot{x} naught are the initial conditions. This dot represents derivative with respect to time.

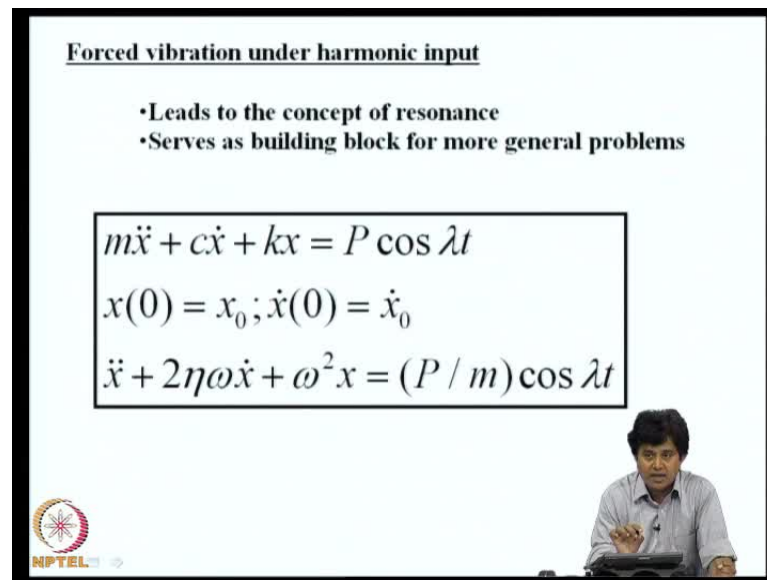
So, this is a linear second order differential equation, It is ordinary differential equation t is the independent variable x is a dependent variable. This is an initial value problem and we need to solve this problem to analyze, the response of a structure, which can be modeled as a single degree freedom system. I will touch upon certain descriptors of the single degree freedom system, in terms of it is initial response, impulse response, frequency response, function and the concept of resonance, which will serve as starting point to understand for example, the notion of a steady state when f of t is modeled as a random process.

If f of t is a random process, what it means is there is an ensemble of f of t . Apart from the fact that at any time t it is a random variable and this family of random variables etcetera. The implication of this on this equation of equilibrium is that this is not a single differential equation, because f of t has is itself of an ensemble of time histories this equation of motion represents an ensemble of equation of motion we can take sample by sample and solve this problem, but that is not what we are going to do. The question we would ask is if I know the complete description of f of t as a random process - that would mean the definition of n th order probability density function for some any n . The question is what is the corresponding description of x of t , can I get the complete description of the random process x of t ? That would be the first question or in other words, how does the uncertainty in specifying f of t propagates through the system and manifest as uncertainty in x of t .

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Forced vibration under harmonic input

- Leads to the concept of resonance
- Serves as building block for more general problems

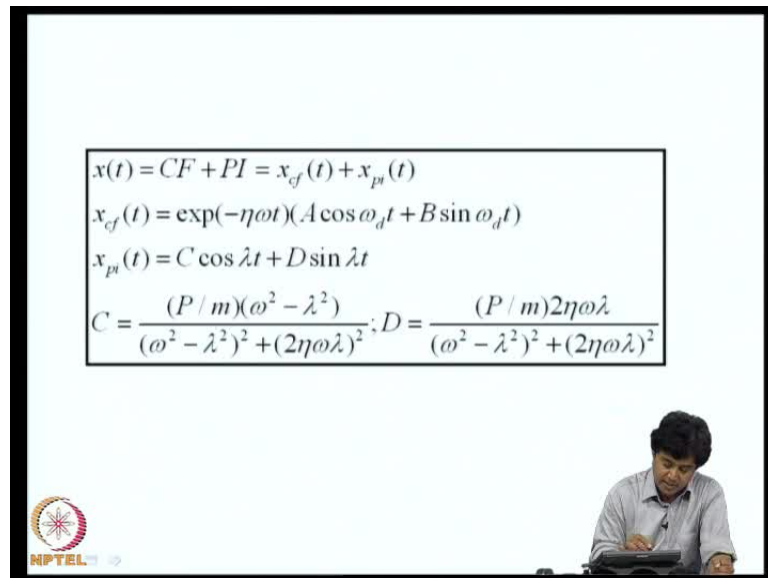
$$m\ddot{x} + c\dot{x} + kx = P \cos \lambda t$$
$$x(0) = x_0; \dot{x}(0) = \dot{x}_0$$
$$\ddot{x} + 2\eta\omega\dot{x} + \omega^2x = (P/m) \cos \lambda t$$

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If f of t is a random process, x of t is also a random process. So, the basic problem in random vibration analysis of a single degree freedom system lies in answering this question. What is the given the complete description of f of t , what is a complete description of x of t ? This question will take up shortly but before we do that we will restrict our attention of f of t being deterministic. Will begin by considering f of t to be harmonic say, $P \cos \lambda t$. A point to be noted in our discussion is the value of this damping c as there are three possible ranges of this damping value, the so called under damp systems and critically damp system and over damp system. We focus our attention in this course on under damp system unless I state it. Otherwise, the default model for damping is that it is systems are under damp by that what it means is in free vibration, the structures oscillate and the vibrations decay exponentially as shown here. That means there is oscillatory decay to the state of rest in free vibration that is when the f of t is 0. So, this is an underline assumption of all our analysis the part of this course in am again restricting attention to only linear systems, if the system in non-linear or it has a different damping model I will have to explicitly make that exception, when I describe the problem.

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$$\begin{aligned}x(t) &= CF + PI = x_{cf}(t) + x_{pi}(t) \\x_{cf}(t) &= \exp(-\eta\omega t)(A \cos \omega_d t + B \sin \omega_d t) \\x_{pi}(t) &= C \cos \lambda t + D \sin \lambda t \\C &= \frac{(P/m)(\omega^2 - \lambda^2)}{(\omega^2 - \lambda^2)^2 + (2\eta\omega\lambda)^2}; D = \frac{(P/m)2\eta\omega\lambda}{(\omega^2 - \lambda^2)^2 + (2\eta\omega\lambda)^2}\end{aligned}$$

So, in standard notations, we divide by m and use the standard notation, where η is the ratio of a damping to the critical damping $\ddot{x} + 2\eta\omega\dot{x} + \omega^2 x = P/m \cos \lambda t$. These are the initial conditions so, x of t as a complementary function and a particular integral. This so the complementary function for this case can be shown to be sums of exponentially decaying harmonics that is exponential minus $\eta\omega t$ a $\cos \omega_d t$ plus $b \sin \omega_d t$ and the particular integral. In this case, we assume it to be harmonic super position of cosine and sine functions these constant C and D have to be obtain by performing by a harmonic balance of the equation of motion with a right hand side, as $p \cos \lambda t$; where as these constant a and b are the arbitrary constants of integration, which have to be evaluated based on the initial conditions of the system.

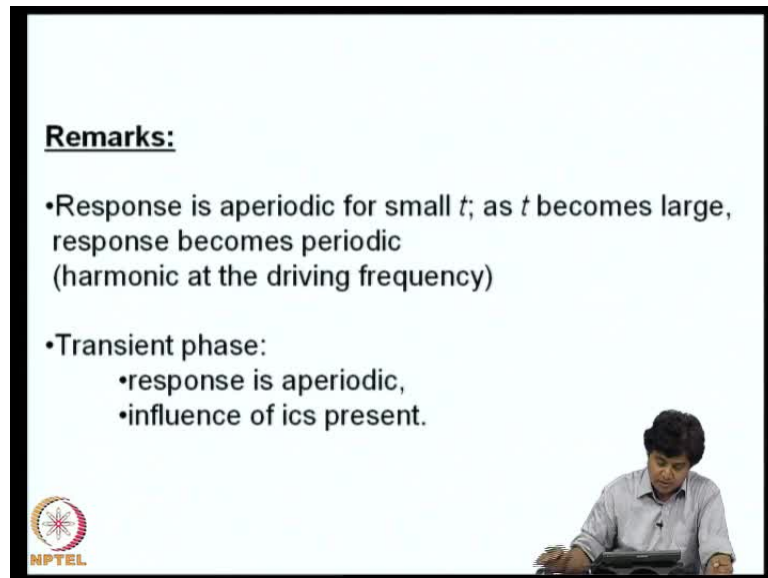
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$$x(t) = \exp(-\eta\omega t) (A \cos \omega_d t + B \sin \omega_d t) + \frac{(P/m) \cos(\lambda t - \theta)}{[(\omega^2 - \lambda^2)^2 + (2\eta\omega\lambda)^2]^{\frac{1}{2}}}$$
$$x_0 = A + \frac{(P/m) \cos(\theta)}{[(\omega^2 - \lambda^2)^2 + (2\eta\omega\lambda)^2]^{\frac{1}{2}}}$$
$$\dot{x}_0 = -\eta\omega A + B\omega_d + \frac{(P/m)\lambda \sin(\theta)}{[(\omega^2 - \lambda^2)^2 + (2\eta\omega\lambda)^2]^{\frac{1}{2}}}$$

So, if you do this harmonic balancing I get for C this quantity for D this quantity, if I now substitute those values of C and D I get the displacement to be given in this form; where A and B are still the constants, which have to be yet determine based on initial conditions and if you actually do that you can show that A and B depend on x naught and x naught dot through this pair of equations. So, I have determined A and B in terms of system parameters and given initial conditions and this is the particular integral, so this is a complete solution. If you now look at the nature of this solution, we can immediately see that x of t is aperiodic, because there is an exponentially $\eta\omega t$, which decays in time, but this component itself the second component itself is periodic with period 2π by λ , which is actually the period of excitation.

Now, as time becomes large this exponential $\eta\omega t$ decays to 0 in which case this part, the first part goes to 0. And consequently what happens? The effect of initial condition x naught and x naught dot are encapsulated in this constant A and B. This part is independent of A and B. So, as t tense to infinity this entire part goes to 0 and effect of initial conditions also consequently will vanish from the response.

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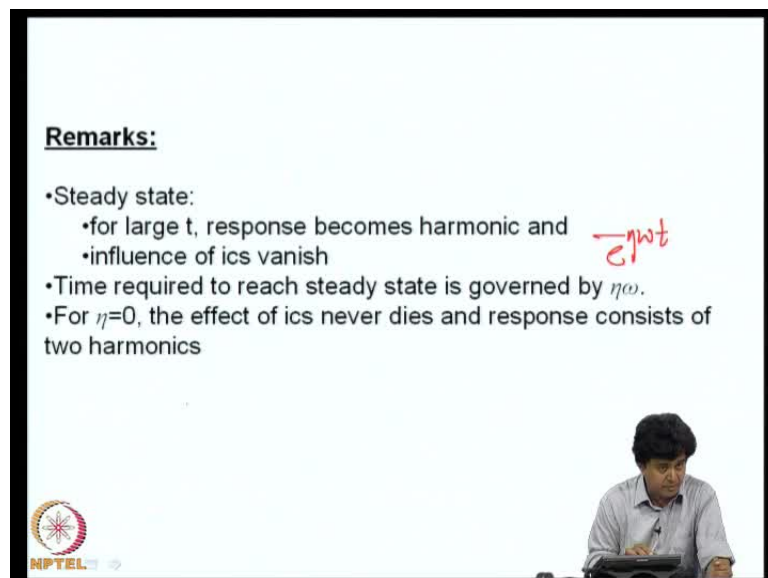
Remarks:

- Response is aperiodic for small t ; as t becomes large, response becomes periodic (harmonic at the driving frequency)
- Transient phase:
 - response is aperiodic,
 - influence of ics present.

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So, if that happens, we say that x of t has reached harmonic steady state. In that situation, the response is periodic at the driving frequency and is independent of initial conditions. The question that will be asking shortly is there a counted part of harmonic steady state when f of t is a random process that is a question that we have to answer. Before that will run through some of the details. So, few remarks response is aperiodic for small t as t becomes large response, becomes periodic in fact it is harmonic at the driving frequency.

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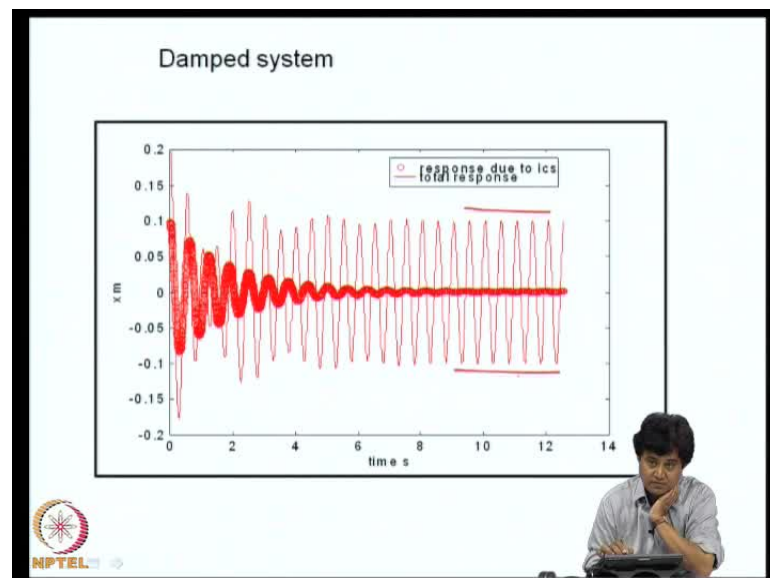
Remarks:

- Steady state:
 - for large t , response becomes harmonic and *time*
 - influence of ics vanish
- Time required to reach steady state is governed by $\eta\omega$.
- For $\eta=0$, the effect of ics never dies and response consists of two harmonics

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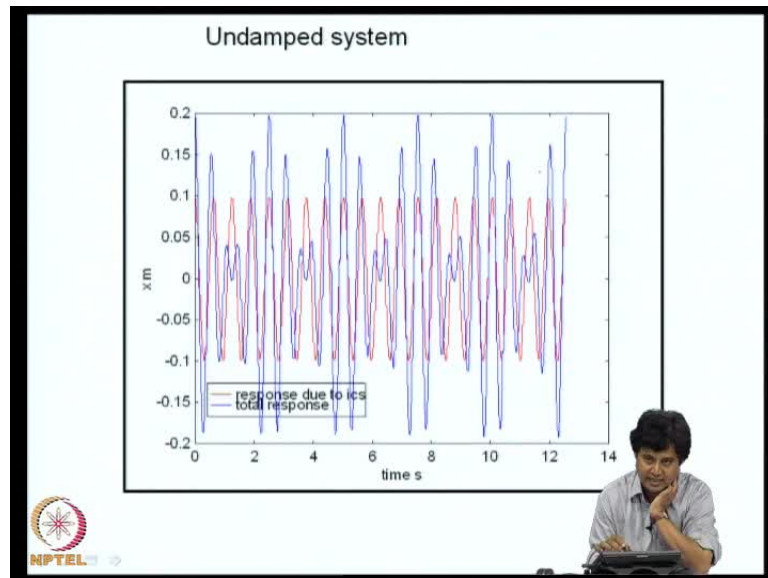
So, there is a transient phase and a steady phase a transient phase is characterized by aperiodic response, in which influence of initial conditions are still present. The steady phase is when the response become periodic and effect of initial conditions vanish, does such things happen, if f of t is stationary is a question that we have to bear in mind. So, steady state for large t response becomes harmonic and influence of initial conditions vanish, the system indeed take some time to reach steady state and that is gone by the factor η of ω . How does is originate? There is a term $\eta \omega t$ so steady state is reached when this term becomes 0. So, the rate at which this goes to 0 depends on value of $\eta \omega$ **large** larger the value of η , ω sooner will be the steady state. If η is 0, system will never reach steady state that is for η equal to 0, the effect of initial condition never dies and response consists of two harmonics and sums of two harmonics could be periodic but in general they are not periodic.

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So, this is the graphical display of this solution. This red line, thin red line that we are seeing is the total response - that means that transient phase plus the steady phase. So, this is the actually the x of t , but on this if I now plot only the terms which are functions of initial conditions, you can see that they oscillate and go to 0. So, initially if you follow now this full lines, you will see that response is aperiodic and after some vanish time they amplitudes have become constant and if you see the phases, they are also constant that take some effort and we see we say that the system has reached steady state, the notion of harmonic steady state. We will see later is related to notion of stationarity.

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So, what does steady state mean? Steady state does not mean that x of t becomes constant. Some property of x of t , namely it is amplitude and it is phase without forcing function that becomes constant. when the function itself is continuous to be function of time it varies with time, if damping is absent as I was telling the initial condition, there is no reason why initial conditions should die and the blue line is the total response and the red line is the component due to initial conditions and total response has two harmonics and this goes on oscillating forever and initial condition effect is felt throughout the time all the values of time.

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Nature of steady state response

$$\lim_{t \rightarrow \infty} x(t) \rightarrow \frac{(P/m) \cos(\lambda t - \theta)}{[(\omega^2 - \lambda^2)^2 + (2\eta\omega\lambda)^2]^{\frac{1}{2}}}$$

$$= \frac{(P/k) \cos(\lambda t - \theta)}{[(1-r^2)^2 + (2\eta r)^2]^{\frac{1}{2}}} = X \cos(\lambda t - \theta); r = \frac{\lambda}{\omega}$$

$$\frac{X}{(P/k)} = DMF = \frac{1}{[(1-r^2)^2 + (2\eta r)^2]^{\frac{1}{2}}}$$

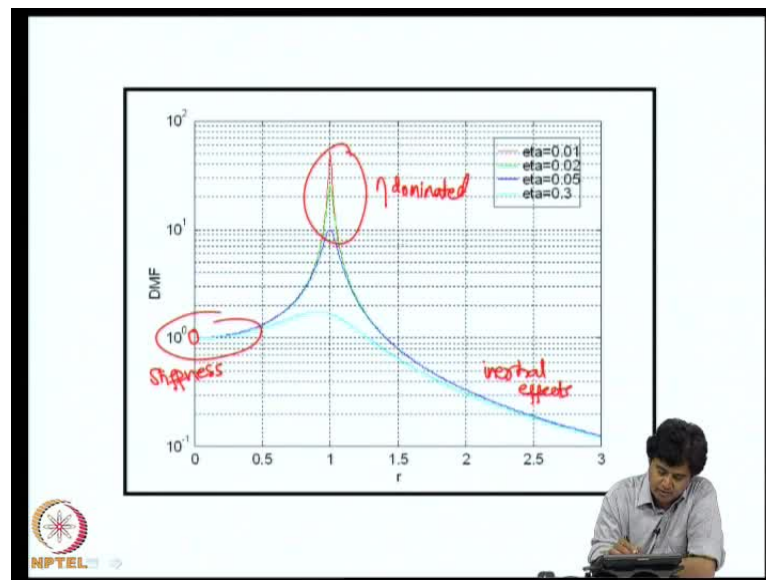
$$\theta = \tan^{-1} \left(\frac{2\eta r}{1-r^2} \right)$$

- DMF=dynamic magnification factor (modification?)
- (P/k)=static response under force P

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Now, we can see the few details about nature of this steady state. Suppose, in the expression for x of t , we allow t to infinity that terms that are actually multiplied by e raise to $\eta \omega t$ go to 0 and what remains is this function. We can slightly reorganize these terms the details can be digested, but if you actually do that we can define a non-dimensional quantity, which is x which is amplitude of response of x of t divide by P by k P by k is a static response. If I where to apply a static force p , the displacement of the system would be P by k , this ratio is known as dynamic magnification factor and it is given by this quantity, this r is the ratio of the driving frequency and the natural frequency of the system this θ , which is the phase is given by $\tan^{-1} \frac{2 \eta r}{1 - r^2}$.

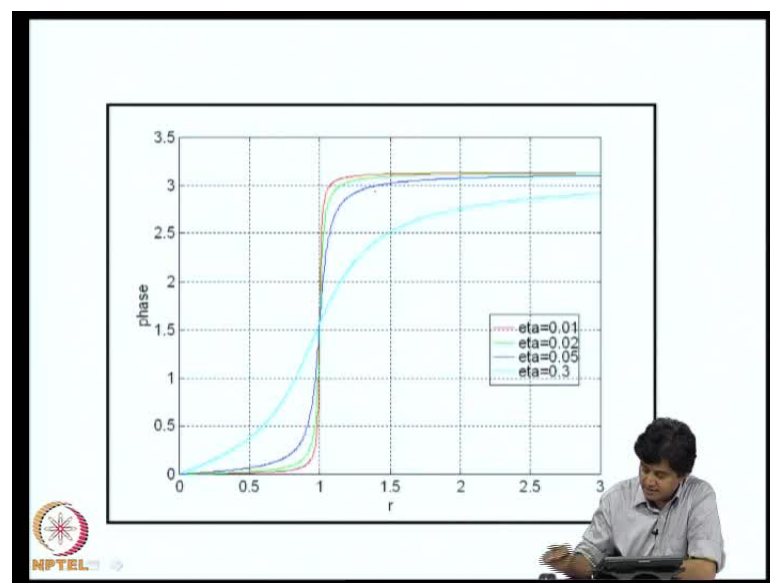
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So, this dynamic magnification factor tells what is a magnification on the static response, because a system is vibrating under dynamic loads, if I where to plot. Now, this dynamic magnification factor as a function of r and η , we will see this familiar curve here, η is changing along these curves, thus see the line blue, green and red lines and η is changing there; an on x-axis of frequency ratio is varying. So, if you as you see r goes to 0 that means driving a $\cos \lambda t$ and λ is 0, the amplitude will be one and there is no dynamic magnification factor, I get as one that means this is the static behavior, this is the static point of static behavior. As frequency increases, we see that in the neighborhood of λ being equal to ω the dynamic magnification factor claims up it is as I has 20, 30, 50 depending upon value of damping, subsequently as further

increase in r the dynamic magnification factor goes on reducing and beyond a point. In fact, it becomes less than the static response and we can divide this graph into three regions: one is the region here which is characterizes by low value of r , where the stiffness of the system dominates - that means response is nearly static and it is dominated by the stiffness characteristics here in the region r close to one. This is eta dominated as r becomes large this behavior is dominated by inertial effects that means the load keeps changing it is sign. So, rapidly that system fails to recognize that that means inertial effects dominate the structural behavior.

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So, this is mass control region, this is damping control region, and this is stiffness control region. Similarly, if you look at the phase plot at in the region of r equal to 1, the phase rapidly changes and goes through a transition, through a value of π and a rapid change in phase is what characterizes is so called resonant condition in the neighborhood of r equal to 1, where there is significant dynamic magnification, we say that structure is under resonance the precise value of r at which this reaches the highest value is not equal to r , equal to one is slightly different and that is actually the point of resonance in the phase diagram. The manifestation of resonance is through a rapid change in phrase angle.

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Resonance in undamped system



$$\ddot{x} + \omega^2 x = (P/m) \cos \lambda t$$

$$x(t) = A \cos \omega t + B \sin \omega t + \frac{(P/m) \cos \lambda t}{\omega^2 - \lambda^2}$$

$$x(0) = 0; \dot{x}(0) = 0$$

$$x(t) = (P/m) \frac{\cos \lambda t - \cos \omega t}{\omega^2 - \lambda^2}$$

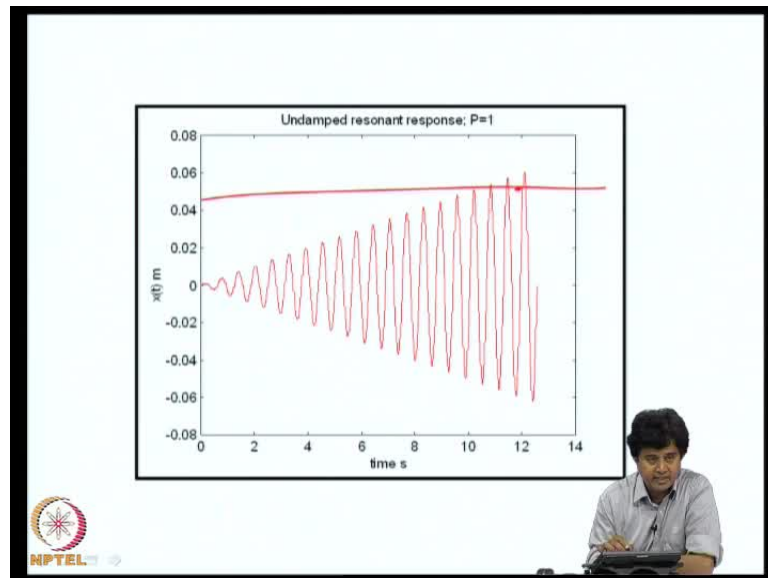
$$\lim_{\lambda \rightarrow \omega} x(t) \rightarrow \frac{Pt}{2m\omega} \sin \omega t$$

$$\lim_{\substack{\lambda \rightarrow \omega \\ t \rightarrow \infty}} x(t) \rightarrow \infty$$



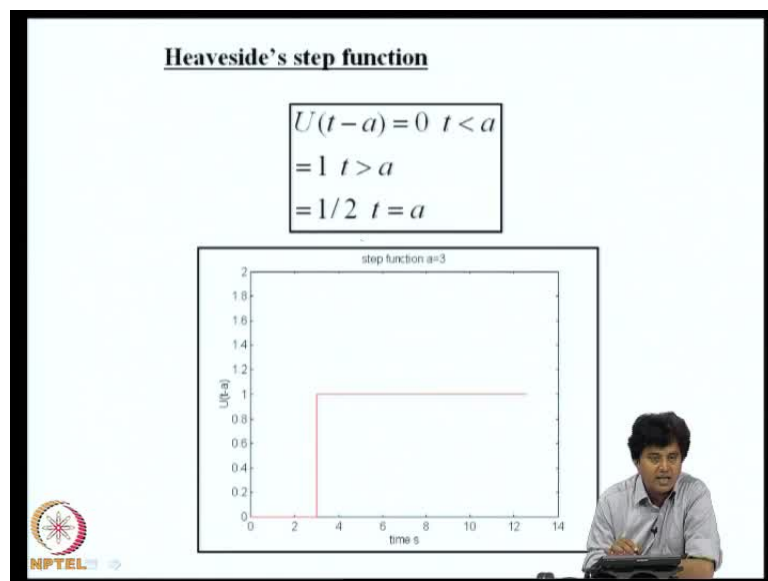
If we now look at this dynamic magnification factor, you see that as damping is reducing the amplification factor is claiming up. So, the question would arise if damping indeed goes to 0. What will be the magnification? It appears from this graph that the dynamic magnification factor becomes unbounded, so what really happens if you apply a harmonic load on a un damped system near resonance, what really happens, if you want to analyze that we can consider the equilibrium equation with no damping here this is a complementary function and this is the particular integral and let us assume the systems starts from rest and we can show that if you impose these initial conditions we can evaluate a and b and if you were to do that I get these expression and in this if I now simulate the resonant condition namely omega going to lambda or limit of lambda going to omega what happens is cos lambda t minus cos lambda t is goes to 0 omega square minus lambda square that also goes to 0 it is indeterminate form.

So, we have to use this one hospital rule and get the limit has lambda goes to omega. So, we have to actually differentiate with respect to lambda. So, if you do that I get this function P t in to sin omega t divided by 2 m omega so as time becomes large this is not a periodic function although the system is driven harmonically and it is system is linear, but address on is not a harmonic it is not even a periodic function because as t becomes large this this amplitude becomes large and in the limit of t tend in to infinity the respond becomes unbounded.

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So, if you plot the time history, you see that there is a linear growth of amplitude as time passes and if there is a critical limit on the say force in the spring, which is k in to x that may be at this point. The structure would simply back, but it would still survive this 12 second of shaking before it actually breaks - that means an un-damped system resonance does not cause instantaneous failure. The amplitudes not growing and at some time when the amplitude crosses a threshold it fails. The description of dynamical is in the description of dynamical system, we introduce what is known as indicial response and impulse response. Indicial response is a response of the system to the unit step function.

A step function if you recall we define as U of t minus a is 0 for t less than a and it is 1 for t greater than a that is as shown here. This is an a is 3 seconds, so till reach the value of 3 it is 0 and yet three jumps. So, if you actually want value of this function at 3 seconds you can take the left hand average of the left hand limit and right hand limit, which is 0 and 1, you can call it as half it is a matter of convention.

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Indicial response analysis

$$m\ddot{x} + c\dot{x} + kx = U(t)$$

$$x(0) = 0; \dot{x}(0) = 0$$

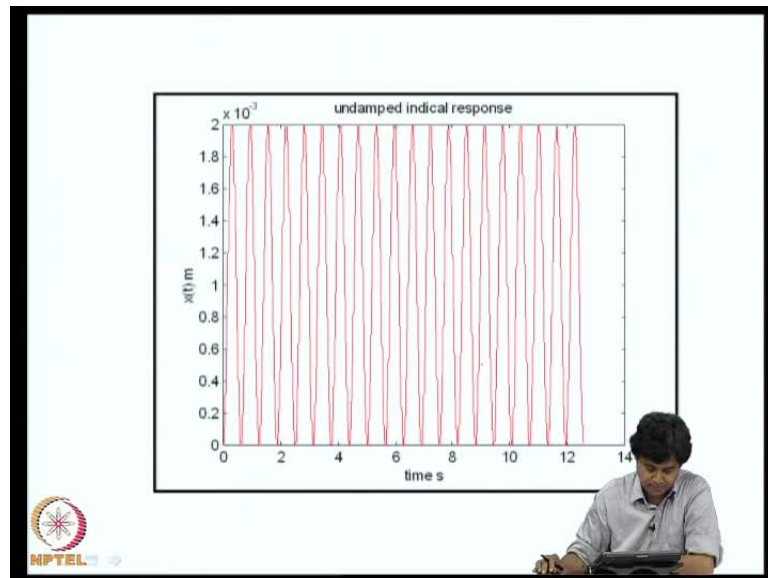
$$x(t) = \exp(-\eta\omega_d t)(A \cos \omega_d t + B \sin \omega_d t) + (1/k)$$

$$x(t) = (1/k) \left[1 - \exp(-\eta\omega_d t) \left(\cos \omega_d t + \frac{\eta}{\sqrt{1-\eta^2}} \sin \omega_d t \right) \right] = G(t)$$

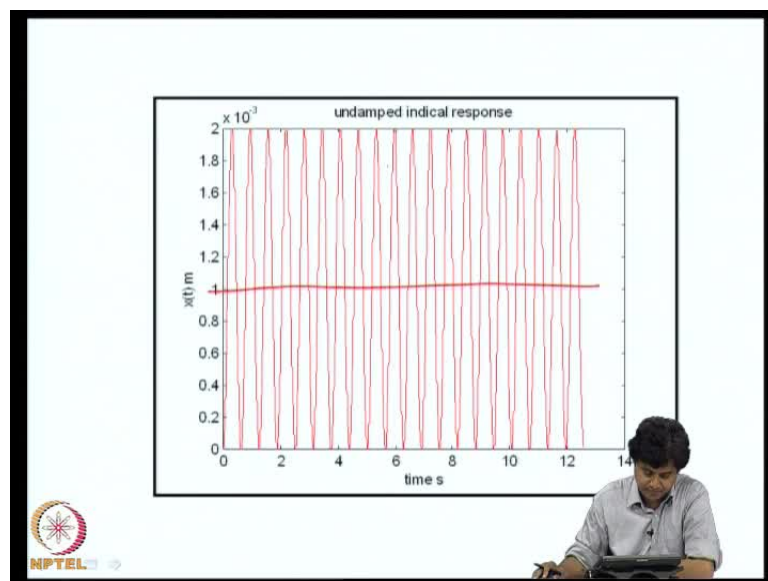



Now, if we consider the response of the system to the suddenly applied load. We can model the load as a u side step function. So, U of t means there is a suddenly applied load at t equal to 0 suppose a system is addressed and you apply a suddenly applied load, constant load, we can analyze this problem there is again the solution will have a complementary function and a particular integral and if you now select the initial conditions, impose the initial condition x of 0 is 0 x dot of 0 is 0 we can evaluate A and B and I get this response. And we call this we denote this by G of t and this is known as indicial response of the system. This is one of the important descriptors of a linear time in variant system.

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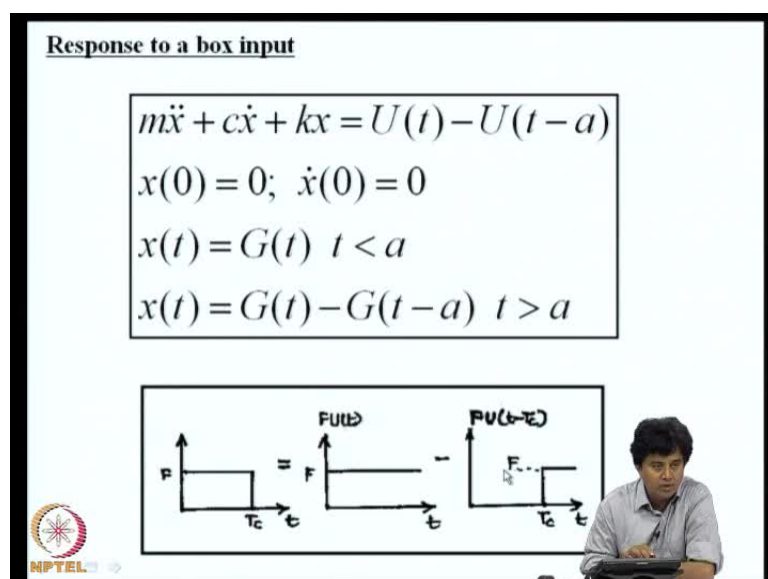


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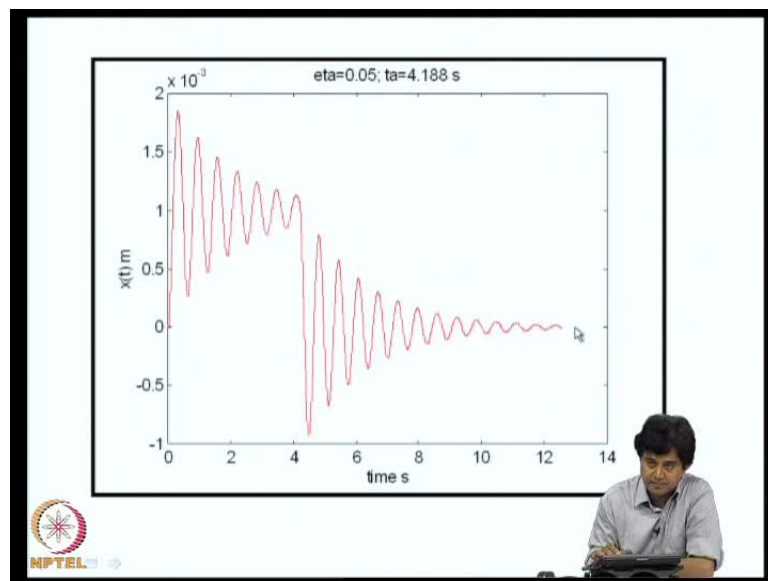
How does indicial response look for a un damped system? You can see here if dumping is 0. This will be 1 and this part is 0 I get 1 minus cos omega t. So that is depicted here the function is oscillating about value of 1 not about 0 about 1, because it is 1 minus cos omega t and this is the un damped indicial response a damped indicial response will oscillate about one, this is one. So, it will respond and it will decay to a constant value s t becomes large. This is nothing but actually the static response of the system under unit low s t becomes large, the response of the system is nothing but the static response because a damped system.

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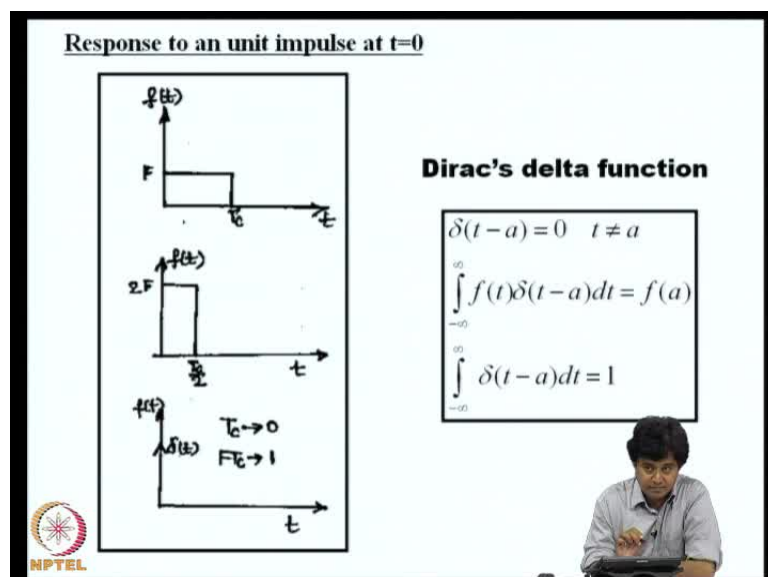


Now, we can consider the response of the system to a box kind of load, as shown here. So, this can be modeled in terms of two give sides step functions, this is f suppose f is 1 here f I have taken as one in this equation this is U of t minus U of t minus a . So, this is a box of width a . So, I can analyze this response, the response will be G of t for t less than or equal to a and for t greater than equal to a it will super position of response due to U of t and U of t minus a which is G of t minus G of t minus a is might have I will **I will** super imposing the force in terms of 2 functions and I am super imposing the response in terms of corresponding responses to these 2 forces.

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So, for a specific case, we can see here this is the response. This is the time at which the box ends box excitation ends so as long as there is the load, it will be oscillating and the load is suddenly removed; it will oscillate and come to rest. This is the response of the system to a box input. Now, what we will do is we shrink this box, now suppose there is the step function it amplitude f and time t_c . Now, it will allow this t_c to go to 0, this and at the same time I will increase f . So that $f t_c$ goes to unit. So these we have to discussed earlier when I discussed about discrete random variables and the relation to continuous random variables, we get what is known as in the limit that I am mentioning we get what is known as direct delta function.

(Refer Slide Time: 52:37)

Response to an unit impulse at $t=0$

$$f(t) = F[U(t) - U(t - T_c)]$$



$$x(t) = FG(t) \quad t < T_c$$

$$x(t) = F[G(t) - G(t - T_c)] \quad t > T_c$$

$$f(t) = FT_c \frac{[U(t) - U(t - T_c)]}{T_c}$$

$$\lim_{\substack{T_c \rightarrow 0 \\ FT_c \rightarrow 1}} f(t) = \frac{dU}{dt} = \delta(t)$$

$$x(t) = FT_c \frac{[G(t) - G(t - T_c)]}{T_c}$$



$$\lim_{\substack{T_c \rightarrow 0 \\ FT_c \rightarrow 1}} x(t) = \frac{dG}{dt} = h(t)$$



So, direct delta function is defined as direct delta of t minus a equal to 0, when t is equal to a , area under this function is one and if it is actually f of t . This function will have value f if the area under the function is 1 at t equal to a itself, we do not define direct delta function, the direct delta function is defined through an integral. This direct delta function can be used to model impulsive loads on the structure. So, how do we get the so called response of the system to an impulse. So, I will model f of t as a box function and the response will be accordingly F in to G of t for t less than or equal to T_c and this is this for t greater than T_c .

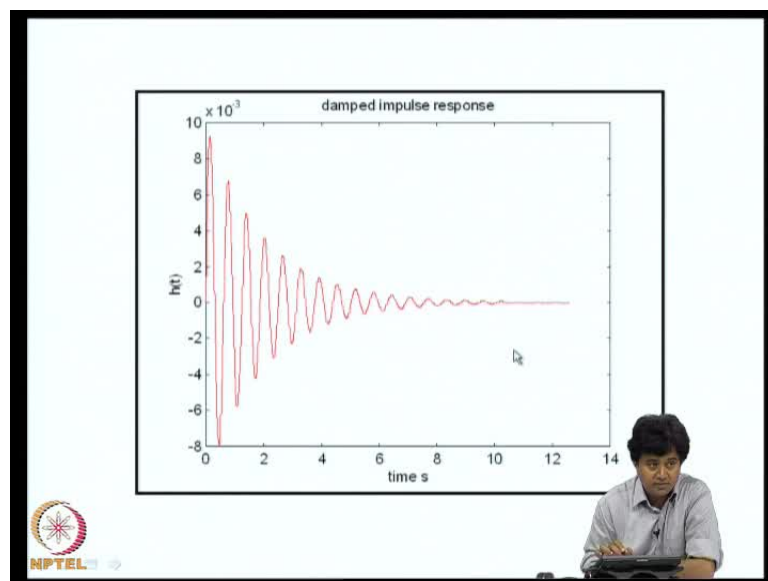
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Impulse response function $h(t)$

$$G(t) = (1/k) \left[1 - \exp(-\eta at) \left(\cos \omega_d t + \frac{\eta}{\sqrt{1-\eta^2}} \sin \omega_d t \right) \right]$$

$$h(t) = \frac{dG}{dt} = \frac{1}{m\omega_d} \exp(-\eta at) \sin \omega_d t$$



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Now, I will rewrite $f(t)$ as $F T c$ divided by $T c$. Now, the limit that I am talking about this $T c$ going to 0 $F T c$ going to 1 and the forcing function goes to $d u$ by $d t$ which is direct delta function, under the same limiting condition what happens to $x(t)$? It will be F of $T c$ in to G of t minus G minus $T c$ by $T c$ so under this limiting condition you see that this quantity is nothing but $d G$ by $d t$. This $h(t)$ is therefore derivative of indicial response. This we known as call it as unit impulse response function. And the expression for that is displayed here this is G of t and this is the differential of G of t . And for an undamped system when η is 0 ω_d will be ω , this is 1 by $m \omega$ sine ωt

and that is the sinusoidal function that is shown here. For a damped system it is a exponentially decaying harmonic sinusoidal function. So, in the next class we will see how the impulse response function can be used to construct the response of the structure to an arbitrary load f of t .

So, the lecture ends with this.