## INSTITUTE OF AERONAUTICAL ENGINEERING

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# PPT ON PROBABILITY THEORY \&STOCHASTIC PROCESS 

II B.Tech I semester (JNTUH-R15)

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# probability introduced through sets and relative frequency 

- Experiment:- a random experiment is an action or process that leads to one of several possible outcomes

Experiment

## Outcomes

Flip a coin

Exam Marks

Assembly Time
Course Grades

Heads, Tails
Numbers: $0,1,2, \ldots$, 100
$\mathrm{t}>0$ seconds
F, D, C, B, A, A+

## Sample Space

- List: "Called the Sample Space"
- Outcomes: "Called the Simple Events"

This list must be exhaustive, i.e. ALL possible outcomes included.

- Die roll $\{1,2,3,4,5\} \quad$ Die roll $\{1,2,3,4,5,6\}$

The list must be mutually exclusive, i.e. no two outcomes can occur at the same time:

Die roll \{odd number or even number\}
Die roll\{ number less than 4 or even number\}

## Sample Space

- A list of exhaustive [don't leave anything out] and mutually exclusive outcomes [impossible for 2 different events to occur in the same experiment] is called a sample space and is denoted by S .
- The outcomes are denoted by $\mathrm{O}_{1}, \mathrm{O}_{2}, \ldots, \mathrm{O}_{\mathrm{k}}$
- Using notation from set theory, we can represent the sample space and its outcomes as:

$$
\cdot S=\left\{O_{1}, O_{2}, \ldots, O_{k}\right\}
$$

- Given a sample space $\mathrm{S}=\left\{\mathrm{O}_{1}, \mathrm{O}_{2}, \ldots, \mathrm{O}_{\mathrm{k}}\right\}$, the probabilities assigned to the outcome must satisfy these requirements:
(1) The probability of any outcome is between 0 and 1
i.e. $0 \leq \mathrm{P}\left(\mathrm{O}_{\mathrm{i}}\right) \leq 1$ for each $i$, and
(2) The sum of the probabilities of all the outcomes equals 1

$$
\text { i.e. } \mathrm{P}\left(\mathrm{O}_{1}\right)+\mathrm{P}\left(\mathrm{O}_{2}\right)+\ldots+\mathrm{P}\left(\mathrm{O}_{\mathrm{k}}\right)=1
$$

$$
\sum_{i=1}^{k} P\left(O_{i}\right)=1
$$

## Relative Frequency

Random experiment with sample space $S$. we shall assign non-negative number called probability to each event in the sample space.
Let $A$ be a particular event in $S$. then "the probability of event $A^{\prime \prime}$ is denoted by $P(A)$.
Suppose that the random experiment is repeated $n$ times, if the event $A$ occurs $n_{A}$ times, then the probability of event A is defined as "Relative frequency "

- Relative Frequency Definition: The probability of an
- event $A$ is defined as

$$
P(A)=\lim _{n \rightarrow \infty} \frac{n_{A}}{n}
$$

## Axioms of Probability

For any event $A$, we assign a number $P(A)$, called the probability of the event $A$. This number satisfies the following three conditions that act the axioms of probability.
(i) $\quad P(A) \geq 0 \quad$ (Probabili ty is a nonnegativ e number)
(ii) $P(\Omega)=1 \quad$ (Probabili ty of the whole set is unity)
(iii) If $A \cap B=\phi$, then $P(A \cup B)=P(A)+P(B)$.
(Note that (iii) states that if $A$ and $B$ are mutually exclusive (M.E.) events, the probability of their union is the sum of their probabilities.)

## Events

- The probability of an event is the sum of the probabilities of the simple events that constitute the event.
- E.g. (assuming a fair die) $S=\{1,2,3,4,5,6\}$ and
- $P(1)=P(2)=P(3)=P(4)=P(5)=P(6)=1 / 6$
- Then:
- $\mathrm{P}($ EVEN $)=\mathrm{P}(2)+\mathrm{P}(4)+\mathrm{P}(6)=1 / 6+1 / 6+1 / 6=$ $3 / 6=1 / 2$


## Conditional Probability

- Conditional probability is used to determine how two events are related; that is, we can determine the probability of one event given the occurrence of another related event.
- Experiment: random select one student in class.
- $P($ randomly selected student is male) $=$
- $P$ (randomly selected student is male/student is on $3^{\text {rd }}$ row) $=$
- Conditional probabilities are written as $\mathbf{P}(\mathbf{A} \mid \mathbf{B})$ and read as "the probability of A given B " and is calculated as

$$
P(A \mid B)=\frac{P(A \text { and } B)}{P(B)}
$$

- $P(A$ and $B)=P(A) * P(B / A)=P(B) * P(A / B)$ both are true
- Keep this in mind!

$$
\begin{aligned}
& P(A \mid B)=\frac{P(A \text { and } B)}{P(B)} \\
& P(B \mid A)=\frac{P(A \text { and } B)}{P(A)}
\end{aligned}
$$

## Bayes' Law

- Bayes' Law is named for Thomas Bayes, an eighteenth century mathematician.
- In its most basic form, if we know $P(B \mid A)$,
- we can apply Bayes' Law to determine $P(A \mid B)$

- The probabilities $P(A)$ and $P\left(A^{C}\right)$ are called prior probabilities because they are determined prior to the decision about taking the preparatory course.
- The conditional probability $\mathrm{P}(\mathrm{A} \mid \mathrm{B})$ is called a posterior probability (or revised probability), because the prior probability is revised after the decision about taking the preparatory course.


## Total probability theorem

- Take events $\mathrm{A}_{\mathrm{i}}$ for $\mathrm{I}=1$ to k to be:
- Mutually exclusive: $A_{i} \cap A_{j}=0 \quad$ for all i,j
- Exhaustive:

$$
A_{1} \cup \cdots \cup A_{k}=S
$$

For any event B on S

$$
\begin{aligned}
& p(B)=p\left(B \mid A_{1}\right) p\left(A_{1}\right)+\cdots+p\left(B \mid A_{k}\right) p\left(A_{k}\right) \\
& p(B)=\sum_{i=1}^{k} p\left(B \mid A_{i}\right) p\left(A_{i}\right)
\end{aligned}
$$

Bayes theorem follows

$$
p\left(A_{j} \mid B\right)=\frac{p\left(A_{j} \cap B\right)}{p(B)}=\frac{p\left(B \mid A_{j}\right) \cdot p(A)}{\sum_{i=1}^{k} p\left(B \mid A_{i}\right) p\left(A_{i}\right)}
$$

## Independence

- Do $A$ and $B$ depend on one another?
- Yes! B more likely to be true if A.
- A should be more likely if B.
- If Independent

$$
\begin{gathered}
p(A \cap B)=p(A) \cdot p(B) \\
p(A \mid B)=p(A) \quad p(B \mid A)=p(B)
\end{gathered}
$$

- If Dependent

$$
\begin{gathered}
p(A \cap B) \neq p(A) \cdot p(B) \\
p(A \cup B)=p(A)+p(B)-p(A \cap B) \\
p(A \cap B)=p(B \mid A) \cdot p(A)
\end{gathered}
$$

## Random variable

- Random variable
- A numerical value to each outcome of a particular experiment

- Example 1 : Machine Breakdowns
- Sample space : $S=$ \{electrical, mechanical,misuse\}
- Each of these failures may be associated with a repair cost
- State space : $\{50,200,350\}$
- Cost is a random variable :50, 200, and 350
- Probability Mass Function (p.m.f.)
- A set of probability value assigned to each of the values taken by the discrete random variable $x_{i}$
$-0 \leq p_{i} \leq 1$ and $\quad \sum_{i} p_{i}=1$
- Probability : $P\left(X=x_{i}\right)=p_{i}$


## Continuous and Discrete random variables

- Discrete random variables have a countable number of outcomes
- Examples: Dead/alive, treatment/placebo, dice, counts, etc.
- Continuous random variables have an infinite continuum of possible values.
- Examples: blood pressure, weight, the speed of a car, the real numbers from 1 to 6.
- Distribution function:

$$
F_{X}(x)=P(X \leq x),-\infty<x<\infty
$$

- If $F_{X}(x)$ is a continuous function of $x$, then $X$ is a continuous random variable.
$-F_{X}(x)$ : discrete in $x \rightarrow$ Discrete rv's
$-F_{X}(x)$ : piecewise continuous $\rightarrow$ Mixed rv's
- PROPERTIES:
- $0 \leq F_{X}(x) \leq 1,-\infty<x<\infty$
- $F_{X}(x)$ : monotonically increasing func. of $x$
. $x \rightarrow-\infty F_{X}(x)=0$ and $x \xrightarrow{\lim } \infty F_{X}(x)=1$


## Probability Density Function (pdf)

- $X$ : continuous rv, then, $f(x)=\frac{d F(x)}{d x}$ is the $p d f$ of $X$.

$$
\begin{aligned}
& C D F \leftarrow \rightarrow p d f \\
& P(X \leq x)=F(x)=\int^{x} f(u) d u,-\infty<x<\infty \\
& P(X \in(a, b])=P(a<X \leq b)=\int_{a}^{b} f_{X}(u) d u
\end{aligned}
$$

- pdf properties:

1. $\quad f(x) \geq 0$ for all $x$.
2. $\int_{-\infty}^{\infty} f(x) d x=1$. $F(t)=\int_{-\infty}^{t} f(x) d x$

$$
=\int_{0}^{t} f(x) d x
$$

## Binomial

- Suppose that the probability of success is $p$
- What is the probability of failure?

$$
q=1-p
$$

- Examples
- Toss of a coin ( $S=$ head): $p=0.5 \Rightarrow q=0.5$
- Roll of a die $(S=1): p=0.1667 \Rightarrow q=0.8333$
- Fertility of a chicken egg ( $S=$ fertile): $p=0.8 \Rightarrow q=0.2$


## binomial

- Imagine that a trial is repeated $n$ times
- Examples
- A coin is tossed 5 times
- A die is rolled 25 times
- 50 chicken eggs are examined
- Assume $p$ remains constant from trial to trial and that the trials are statistically independent of each other
- Example
- What is the probability of obtaining 2 heads from a coin that was tossed 5 times?
$P(H H T T T)=(1 / 2)^{5}=1 / 32$


## Poisson

- When there is a large number of trials, but a small probability of success, binomial calculation becomes impractical
- Example: Number of deaths from horse kicks in the Army in different years
- The mean number of successes from $n$ trials is $\mu=n p$
- Example: 64 deaths in 20 years from thousands of soldiers

If we substitute $\mu / n$ for $p$, and let $n$ tend to infinity, the binomial distribution becomes the Poisson distribution:

$$
P(x)=\frac{e^{-\mu} \mu^{x}}{x!}
$$

## poisson

- Poisson distribution is applied where random events in space or time are expected to occur
- Deviation from Poisson distribution may indicate some degree of non-randomness in the events under study
- Investigation of cause may be of interest


## Exponential Distribution

The random variable $X$ that equals the distance between successive counts of a Poisson process with mean $\lambda>0$ is an exponential random variable with parameter $\lambda$. The probability density function of $X$ is

$$
\begin{equation*}
f(x)=\lambda e^{-\lambda x} \text { for } 0 \leq x<\infty \tag{4-14}
\end{equation*}
$$

If the random variable $X$ has an exponential distribution with parameter $\lambda$,

$$
\begin{equation*}
\mu=E(X)=\frac{1}{\lambda} \quad \text { and } \quad \sigma^{2}=V(X)=\frac{1}{\lambda^{2}} \tag{4-15}
\end{equation*}
$$

In a large corporate computer network, user log-ons to the system can be modeled as a Poisson process with a mean of $25 \log$-ons per hour. What is the probability that there are no $\log$ ons in an interval of 6 minutes?

Let $X$ denote the time in hours from the start of the interval until the first log-on. Then, $X$ has an exponential distribution with $\lambda=25 \log$-ons per hour. We are interested in the probability that $X$ exceeds 6 minutes. Because $\lambda$ is given in log-ons per hour, we express all time units in hours. That is, 6 minutes $=0.1$ hour. The probability requested is shown as the shaded area under the probability density function in Fig. 4-23. Therefore,

$$
P(X>0.1)=\int_{0.1}^{\infty} 25 e^{-25 x} d x=e^{-25(0.1)}=0.082
$$

Also, the cumulative distribution function can be used to obtain the same result as follows:

$$
P(X>0.1)=1-F(0.1)=e^{-25(0.1)}
$$

An identical answer is obtained by expressing the mean number of log-ons as $0.417 \log$ ons per minute and computing the probability that the time until the next log-on exceeds 6 minutes. Try it.

What is the probability that the time until the next $\log$-on is between 2 and 3 minutes? Upon converting all units to hours,

$$
P(0.033<X<0.05)=\int_{0.033}^{0.05} 25 e^{-25 x} d x=-\left.e^{-25 x}\right|_{0.033} ^{0.05}=0.152
$$

An alternative solution is

$$
P(0.033<X<0.05)=F(0.05)-F(0.033)=0.152
$$

Determine the interval of time such that the probability that no log-on occurs in the i val is 0.90 . The question asks for the length of time $x$ such that $P(X>x)=0.90$. Now,

$$
P(X>x)=e^{-25 x}=0.90
$$

Take the (natural) $\log$ of both sides to obtain $-25 x=\ln (0.90)=-0.1054$. Therefore,

$$
x=0.00421 \text { hour }=0.25 \text { minute }
$$

Furthermore, the mean time until the next log-on is

$$
\mu=1 / 25=0.04 \text { hour }=2.4 \text { minutes }
$$

The standard deviation of the time until the next log-on is

$$
\sigma=1 / 25 \text { hours }=2.4 \text { minutes }
$$

## Uniform

All (pseudo) random generators generate random deviates of $\mathrm{U}(0,1)$ distribution; that is, if you generate a large number of random variables and plot their empirical distribution function, it will approach this distribution in the limit.
$U(a, b) \rightarrow$ pdf constant over the $(a, b)$ interval and CDF is the ramp function

$$
f(x)= \begin{cases}\frac{1}{b-a}, & a<x<b \\ 0, & \text { otherwise }\end{cases}
$$



## Uniform distribution

$\mathrm{F}(\mathrm{x})= \begin{cases}0, & \mathrm{x}<\mathrm{a}, \\ \frac{\mathrm{x}-a}{b-a}, & \mathrm{a}<\mathrm{x}<\mathrm{b} \\ 1, & \mathrm{x}>\mathrm{b} .\end{cases}$

## Gaussian (Normal) Distribution

- Bell shaped pdf - intuitively pleasing!
- Central Limit Theorem: mean of a large number of mutually independent rv's (having arbitrary distributions) starts following Normal distribution as $n \rightarrow$

$$
f(x)=\frac{1}{\sigma \sqrt{2 \pi}} \mathrm{e}^{-\frac{1}{2}\left(\frac{x-\mu}{\sigma}\right)^{2}},-\infty<x<\infty
$$

- $\mu$ : mean, $\sigma$ : std. deviation, $\sigma^{2}$ : variance ( $N(\mu$, $\left.\sigma^{2}\right)$ )
- $\mu$ and $\sigma$ completely describe the statistics. This is significant in statistical estimation/signal processing/communication theory etc.
- $N(0,1)$ is called normalized Guassian.
- $N(0,1)$ is symmetric i.e.
$-f(x)=f(-x)$
$-F(z)=1-F(z)$.
- Failure rate $h(t)$ follows IFR behavior.
- Hence, $N($ ) is suitable for modeling long-term wear or aging related failure phenomena


## Exponential Distribution

$$
\begin{aligned}
& f(t)=\sum_{i=1}^{k} \alpha_{i} \lambda_{i} \mathrm{e}^{-\lambda_{i} t}, t>0, \lambda_{i}>0, \alpha_{i}>0, \sum_{i=1}^{k} \alpha_{i} \\
& \\
& F(t)=\sum_{i} \alpha_{i}\left(1-\mathrm{e}^{-\lambda_{i} t}\right), t \geq 0 \\
& \\
& h(t)=\frac{\sum_{i} \alpha_{i} \lambda \mathrm{e}^{-\lambda_{i} t}}{\sum_{i} \alpha_{i} \lambda \mathrm{e}^{-\lambda_{i} t}}, t \geq 0
\end{aligned}
$$

## Conditional Distributions

- The conditional distribution of $Y$ given $X=1$ is:
- While marginal distributions are obtained from the bivariate by summing, conditional distributions are obtained by "making a cut" through the bivariate distribution


## The Expectation of a Random Variable

Expectation of a discrete random variable with p.m.f

$$
E(X)=\sum_{i}^{P\left(X=x_{i}\right)=p_{i}} p_{i} x_{i}
$$

Expectation of a continuous random variable with p.d.f $f(x)$

$$
E(X)=\int_{\text {state space }} x f(x) d x
$$

expectation of $X=$ mean of $X=$ average of $X$

$$
\begin{array}{ll}
E[X]=\bar{X}=\int_{-\infty}^{\infty} x f_{X}(x) d x & \text { continuous r.v. } \\
E[X]=\bar{X}=\sum_{i=1}^{N} x_{i} P\left(x_{i}\right) & \text { discrete r.v. }
\end{array}
$$

$$
f_{X}(x+a)=f_{X}(-x+a), \forall x \Rightarrow E[X]=a
$$

$X$ r.v. $\Rightarrow Y=g(X)$ r.v. $\quad$ Ex: $Y=g(X)=X^{2}$
$P(X=0)=P(X=-1)=P(X=1)=\frac{1}{3} \quad P(Y=0)=\frac{1}{3} \quad P(Y=1)=\frac{2}{3}$
Expectation
expectation of a function of a r.v. $X$

$$
\begin{array}{ll}
E[g(X)]=\int_{-\infty}^{\infty} g(x) f_{X}(x) d x & \text { continuous r.v. } \\
E[g(X)]=\sum_{i=1}^{\infty} g\left(x_{i}\right) P\left(x_{i}\right) & \text { discrete r.v. }
\end{array}
$$

conditional expectation of a r.v. $X$

$$
\begin{gathered}
E[X \mid B]=\int_{-\infty}^{\infty} x f_{X}(x \mid B) d x \\
E[X \mid B]=\sum_{i=1}^{\infty} x_{i} P\left(x_{i} \mid B\right)
\end{gathered}
$$

continuous r.v.
discrete r.v.

Ex: $\quad B=\{X \leq b\}$
$f_{X}(x \mid X \leq b)= \begin{cases}\frac{f_{X}(x)}{\int_{-\infty}^{b} f_{X}(x) d x}, & x<b \\ 0, & x \geq b\end{cases}$

$$
E[X \mid X \leq b]=\frac{\int_{-\infty}^{b} x f_{X}(x) d x}{\int_{-\infty}^{b} f_{X}(x) d x}
$$

## Moments

$n$-th moment of a r.v. $X$

$$
\begin{gathered}
m_{n}=E\left[X^{n}\right]=\int_{-\infty}^{\infty} x^{n} f_{X}(x) d x \\
m_{n}=E\left[X^{n}\right]=\sum_{i=1}^{N} x_{i}^{n} P\left(x_{i}\right) \\
m_{0}=1 \\
m_{1}=\bar{X}
\end{gathered}
$$

discrete r.v.
properties of expectation:

$$
\begin{aligned}
& \text { (1) } E[c]=c \quad c \text {-- constant } \\
& \text { (2) } E[\operatorname{ag}(X)+b h(X)]=a E[g(X)]+b E[h(X)] \\
& \text { PF: } E[c]=\int_{-\infty}^{\infty} c f_{X}(x) d x=c \int_{-\infty}^{\infty} f_{X}(x) d x=c \\
& E[a g(X)+b h(X)]=\int_{-\infty}^{\infty}\{a g(x)+b h(x)\} f_{X}(x) d x \\
& =a \int_{-\infty}^{\infty} g(x) f_{X}(x) d x+b \int_{-\infty}^{\infty} h(x) f_{X}(x) d x=a E[g(X)]+b E[h(X)]
\end{aligned}
$$

variance of a r.v. $X$

$$
\begin{aligned}
\sigma_{X}^{2}=\mu_{2} & =E\left[(X-\bar{X})^{2}\right]=E\left[X^{2}-2 \bar{X} X+\bar{X}^{2}\right] \\
& =E\left[X^{2}\right]-2 \bar{X} E[X]+\bar{X}^{2}=m_{2}-m_{1}^{2}
\end{aligned}
$$

standard deviation of a r.v. $X=\sigma_{X}(\geq 0)$
skewness of a r.v. $X=\frac{\mu_{3}}{\sigma_{X}^{3}}$
Ex 3.2-1 \& Ex3.2-2:

$$
f_{X}(x) \text { symmetric about } x=\bar{X} \Rightarrow \mu_{3}=0
$$

$$
f_{X}(x)= \begin{cases}\frac{1}{b} e^{-\frac{x-a}{b}}, & x>a \\ 0, & \mathrm{x}<a\end{cases}
$$

$$
\begin{gathered}
m_{1}=E[X]=\int_{a}^{\infty} x \frac{1}{b} e^{-\frac{x-a}{b}} d x=a+b \\
m_{2}=E\left[X^{2}\right]=\int_{a}^{\infty} x^{2} \frac{1}{b} e^{-\frac{x-a}{b}} d x=(a+b)^{2}+b^{2} \\
\sigma_{X}^{2}=\mu_{2}=m_{2}-m_{1}^{2}=b^{2} \\
m_{3}=E\left[X^{3}\right]=\int_{a}^{\infty} x^{3} \frac{1}{b} e^{-\frac{x-a}{b}} d x=a^{3}+3 a^{2} b+6 a b^{2}+6 b^{3} \\
u_{3}=E\left[(X-\bar{X})^{3}\right]=E\left[X^{3}-3 X^{2} \bar{X}+3 X \bar{X}^{2}-\bar{X}^{3}\right]=m_{3}-3 m_{1} m_{2}+3 m_{1}^{2} m_{1}-m_{1}^{3} \\
=a^{3}+3 a^{2} b+6 a b^{2}+6 b^{3}-3(a+b)\left\{(a+b)^{2}+b^{2}\right\}+2(a+b)^{3}=2 b^{3}
\end{gathered}
$$

skewness of a r.v. $X=\frac{\mu_{3}}{\sigma_{X}^{3}}=\frac{2 b^{3}}{b^{3}}=2$

Chebychev's inequality $P[|X-\bar{X}| \geq \varepsilon] \leq \frac{\sigma_{X}^{2}}{\varepsilon^{2}}$

$$
\begin{array}{r}
\sigma_{X}^{2}=\int_{-\infty}^{\infty}(x-\bar{X})^{2} f_{X}(x) d x \geq \int_{|x-\bar{X}| \geq \varepsilon}(x-\overline{\bar{X}})^{2} f_{X}(x) d x \\
\geq \varepsilon^{2} \int_{|x-\bar{X}| \geq \varepsilon} f_{X}(x) d x=\varepsilon^{2} P[|X-\bar{X}| \geq \varepsilon]
\end{array}
$$

Markov's inequality

$$
P[X<0]=0 \Rightarrow P[X \geq a] \leq \frac{E[X]}{a}
$$

Ex 3.2-3: $P\left[|X-\bar{X}| \geq 3 \sigma_{X}\right] \leq \frac{\sigma_{X}^{2}}{9 \sigma_{X}^{2}}=\frac{1}{9}$

Characteristic function of r.v. $X$

$$
\begin{aligned}
& \Phi_{X}(\omega)=E\left[e^{j \omega X}\right]=\int_{-\infty}^{\infty} f_{X}(x) e^{j \omega x} d x \\
& f_{X}(x)=\frac{1}{2 \pi} \int_{-\infty}^{\infty} \Phi_{X}(\omega) e^{-j \omega x} d \omega \quad \text { Fourier transform } \\
& \left|\Phi_{X}(\omega)\right| \leq \int_{-\infty}^{\infty}\left|f_{X}(x)\right|\left|e^{j \omega x}\right| d x \leq \int_{-\infty}^{\infty} f_{X}(x) d x=1=\Phi_{X}(0) \\
& \left.\frac{d^{n} \Phi_{X}(\omega)}{d \omega^{n}}\right|_{\omega=0}=\left.\int_{-\infty}^{\infty} f_{X}(x) j^{n} x^{n} e^{j \omega x} d x\right|_{\omega=0}=j^{n} \int_{-\infty}^{\infty} f_{X}(x) x^{n} d x=j^{n} E\left[X^{n}\right] \\
& m_{n}=\left.(-j)^{n} \frac{d^{n} \Phi_{X}(\omega)}{d \omega^{n}}\right|_{\omega=0}
\end{aligned}
$$

## Functions That Give Moments

Moment generating function of r.v. $X$

$$
M_{X}(v)=E\left[e^{v X}\right]=\int_{-\infty}^{\infty} f_{X}(x) e^{v x} d x
$$

$$
\left.\frac{d^{n} M_{X}(v)}{d v^{n}}\right|_{v=0}=\left.\int_{-\infty}^{\infty} f_{X}(x) x^{n} e^{v x} d x\right|_{v=0}=\int_{-\infty}^{\infty} f_{X}(x) x^{n} d x=m_{n}
$$

Ex 3.3-1 \& Ex 3.3-2:

$$
f_{X}(x)= \begin{cases}\frac{1}{b} e^{-\frac{x-a}{b}}, & x>a \\ 0, & \mathrm{x}<a\end{cases}
$$

$$
\begin{gathered}
\Phi_{X}(\omega)=E\left[e^{j \omega X}\right]=\frac{1}{b} e^{\frac{a}{b}} \int_{a}^{\infty} e^{-\left(\frac{1}{b}-j \omega\right) x} d x=\left.\frac{1}{b} e^{\frac{a}{b}} \frac{e^{-\left(\frac{1}{b}-j \omega\right) x}}{-\left(\frac{1}{b}-j \omega\right)}\right|_{x=a} ^{\infty} \\
=\frac{1}{b} e^{\frac{a}{b}} \frac{e^{-\left(\frac{1}{b}-j \omega\right) a}}{\left(\frac{1}{b}-j \omega\right)}=\frac{e^{j \omega a}}{1-j \omega b} \quad \frac{d \Phi_{X}(\omega)}{d \omega}=\frac{j a e^{j \omega a}(1-j \omega b)+e^{j \omega a} j b}{(1-j \omega b)^{2}} \\
M_{X}(v)=E\left[e^{v X}\right]=\frac{e^{v a}}{1-v b} \quad \frac{d M_{X}(v)}{d v}=\frac{a e^{v a}(1-v b)+e^{v a} b}{(1-v b)^{2}} \\
m_{1}=\left.(-j) \frac{d \Phi_{X}(\omega)}{d \omega}\right|_{\omega=0}=a+b \quad m_{1}=\left.\frac{d M_{X}(v)}{d v}\right|_{v=0}=a+b
\end{gathered}
$$

Chernoff's inequality
Ex 3.3-3:

$$
v>0
$$

$$
\begin{aligned}
P[X \geq a]=\int_{a}^{\infty} & f_{X}(x) d x=\int_{-\infty}^{\infty} f_{X}(x) u(x-a) d x \\
& \leq \int_{-\infty}^{\infty} f_{X}(x) e^{v(x-a)} d x=e^{-v a} M_{X}(v)
\end{aligned}
$$

## Transformations of a Random Variable

$$
Y=T(X) \quad f_{X}(
$$

monotone increasing $\Leftrightarrow$

$$
T\left(x_{1}\right)<T\left(x_{2}\right) \text { for any } x_{1}<x_{2}
$$

monotone decreasing $\Leftrightarrow$
$T\left(x_{1}\right)>T\left(x_{2}\right)$ for any $x_{1}<x_{2}$


Assume monotone increasing $T(\bullet) \quad Y=T(X)$

$$
F_{Y}\left(y_{0}\right)=P\left[Y \leq y_{0}\right]=P\left[X \leq x_{0}\right]=F_{X}\left(x_{0}\right)
$$

$$
\int_{-\infty}^{y_{0}} f_{Y}(y) d y=\int_{-\infty}^{T^{-1}\left(y_{0}\right)} f_{X}(x) d x
$$

$$
f_{Y}\left(y_{0}\right)=f_{X}\left[T^{-1}\left(y_{0}\right)\right] \frac{d T^{-1}\left(y_{0}\right)}{d y_{0}}
$$

$$
f_{Y}(y)=f_{X}\left[T^{-1}(y)\right] \frac{d T^{-1}(y)}{d y}=f_{X}(x) \frac{d x}{d y}
$$

Assume monotone decreasing $T(\bullet) \quad Y=T(X)$

$$
F_{Y}\left(y_{0}\right)=P\left[Y \leq y_{0}\right]=P\left[X \geq x_{0}\right]=1-F_{X}\left(x_{0}\right)
$$

$$
f_{Y}(y)=-f_{X}(x) \frac{d x}{d y}
$$

monotone $T(\bullet) \Rightarrow \quad f_{Y}(y)=f_{X}(x)\left|\frac{d x}{d y}\right|=f_{X}(x) \frac{1}{\left|\frac{d y}{d x}\right|}$
nonmonotone $T(\bullet)$

$$
\begin{aligned}
Y & =T(X) \\
f_{Y}(y) & =\sum_{n} \frac{f_{X}\left(x_{n}\right)}{\left.\left|\frac{d T(x)}{d x}\right|_{x=x_{n}} \right\rvert\,}
\end{aligned}
$$



Ex 3.4-2:

$$
\begin{aligned}
Y & =T(X)=c X^{2} \quad \text { nonmonotons } \\
f_{Y}(y) & =f_{X}(\sqrt{y / c})\left|\frac{d \sqrt{y / c}}{d y}\right| \\
& +f_{X}(-\sqrt{y / c})\left|\frac{-d \sqrt{y / c}}{d y}\right| \\
= & \frac{f_{X}(\sqrt{y / c})+f_{X}(-\sqrt{y / c})}{2 \sqrt{c y}}, \quad y \geq 0
\end{aligned}
$$

## MULTIPLE RANDOM VARIABLES and OPERATIONS: MULTIPLE RANDOM VARIABLES : <br> Vector Random Variables

A vector random variable $X$ is a function that assigns a vector of real numbers to each outcome $\zeta$ in $S$, the sample space of the random experiment

## Events and Probabilities

## EXAMPLE 4.4

Consider the tow-dimensional random variable $\mathbf{X}=(X, Y)$. Find the region of the plane corresponding to the events

$$
\begin{aligned}
& A=\{X+Y \leq 10\}, \\
& B=\{\min (X, Y) \leq 5\}, \text { and } \\
& C=\left\{X^{2}+Y^{2} \leq 100\right\} .
\end{aligned}
$$

The regions corresponding to events $A$ and $C$ are straightforward to find and are shown in Fig. 4.1.

## FIGURE 4.1

Examples of two-dimensional events




## Independence

If the one-dimensional random variable $X$ and $Y$ are "independent," if $A_{1}$ is any event that involves $X$ only and $A_{2}$ is any event that involves $Y$ only, then

$$
P\left[X \text { in } A_{1}, Y \text { in } A_{2}\right]=P\left[X \text { in } A_{1}\right] P\left[Y \text { in } A_{2}\right] .
$$

In the general case of n random variables, we say that the random variables $X_{1}, X_{2}, \ldots, X_{\mathrm{n}}$ are independent if

$$
\begin{equation*}
P\left[X_{1} \text { in } A_{1}, \ldots, X_{n} \text { in } A_{n}\right]=P\left[X_{1} \text { in } A_{1}\right] \cdots P\left[X_{n} \text { in } A_{n}\right] \tag{4.3}
\end{equation*}
$$

where the $A_{k}$ is an event that involves $X_{k}$ only.

FIGURE 4.3
Some two-dimensional non-product-form events.



## Pairs of Discrete Random Variable

Let the vector random variable $\mathbf{X}=(X, Y)$ assume values from some countable se§ $=\left\{\left(x_{j}, y_{k}\right), j=1,2, \ldots, k=1,2, \ldots\right\}$.The joint probability mass function of $\mathbf{X}$ specifies the probabilities of the product-form event

$$
\begin{align*}
& \quad\left\{X=x_{j}\right\} \cap\left\{Y=y_{k}\right\}: \\
& p_{X, Y}\left(x_{j}, y_{k}\right)=P\left[\left\{X=x_{j}\right\} \cap\left\{Y=y_{k}\right\}\right] \\
& \equiv P\left[X=x_{j}, Y=y_{k}\right] \quad j=1,2, \ldots \quad k=1,2, \ldots \tag{4.4}
\end{align*}
$$

The probability of any event $A$ is the sum of the pmf over the outcomes in $A$

$$
\begin{equation*}
P[X \text { in } A]=\sum_{\left(x_{j}, y_{k}\right)} \sum_{\text {in } A} p_{X, Y}\left(x_{j}, y_{k}\right) . \tag{4.5}
\end{equation*}
$$

$$
\begin{equation*}
\sum_{j=1}^{\infty} \sum_{k=1}^{\infty} p_{X, Y}\left(x_{j}, y_{k}\right)=1 \tag{4.6}
\end{equation*}
$$

The marginal probability mass functions :

$$
\begin{align*}
p_{X}\left(x_{j}\right) & =P\left[X=x_{j}\right\rfloor \\
& =P\left[X=x_{j}, Y=\text { anything }\right] \\
& =P\left[\left\{X=x_{j} \text { and } Y=y_{1}\right\} \cup\left\{X=x_{j} \text { and } Y=y_{2}\right\} \cup \cdots\right] \\
& =\sum_{k=1}^{\infty} p_{X, Y}\left(x_{j}, y_{k}\right)  \tag{4.7a}\\
p_{Y}\left(y_{k}\right) & =P\left[Y=y_{k}\right] \\
& =\sum_{j=1}^{\infty} p_{X, Y}\left(x_{j}, y_{k}\right) \tag{4.7b}
\end{align*}
$$

## The Joint cdf of $X$ and $Y$

The joint cumulative distribution function of X and Y is defined as the probability of the product-form event $\left\{X \leq x_{1}\right\} \cap\left\{Y \leq y_{1}\right\}$ ":

$$
\begin{equation*}
F_{X, Y}\left(x_{1}, y_{1}\right)=P\left[X \leq x_{1}, Y \leq y_{1}\right] \tag{4.8}
\end{equation*}
$$

The joint cdf is nondecreasing in the "northeast" direction,
(i) $\quad F_{X, Y}\left(x_{1}, y_{1}\right) \leq F_{X, Y}\left(x_{2}, y_{2}\right)$ if $x_{1} \leq x_{2}$ and $y_{1} \leq y_{2}$,

It is impossible for either $X$ or $Y$ to assume a value less than $-\infty$ therefore
(ii) $\quad F_{X, Y}\left(-\infty, y_{1}\right)=F_{X, Y}\left(x_{2},-\infty\right)=0$

It is certain that $X$ and $Y$ will assume values less than infinity, therefore
(iii) $\quad F_{X, Y}(\infty, \infty)=1$.

If we let one of the variables approach infinity while keeping the other fixed, we obtain the marginal cumulative distribution functions

$$
\text { (iv) } \quad F_{X}(x)=F_{X, Y}(x, \infty)=P[X \leq x, Y \leq \infty]=P[X \leq x]
$$

and

$$
F_{Y}(y)=F_{X, Y}(\infty, y)=P[Y \leq y] .
$$

Recall that the cdf for a single random variable is continuous form the right. It can be shown that the joint cdf is continuous from the "north" and from the "east"

$$
\text { (v) } \quad \lim _{x \rightarrow a^{+}} F_{X, Y}(x, y)=F_{X, Y}(a, y)
$$

and

$$
\lim _{y \rightarrow b^{+}} F_{X, Y}(x, y)=F_{X, Y}(x, b)
$$

## FIGURE 4.4

The joint cumulative distribution function is defined as the probability of the semi-infinite rectangle defined by the point $\left(x_{1}, y_{1}\right)$.


## FIGURE 4.5

The marginal coff's are the probabilities of these halfplanes.

$F_{Y\left(y_{1}\right)}=P\left[X<\infty, Y \leqslant y_{1}\right]$

## The Joint pdf of Two Jointly Continuous Random Variables

We say that the random variables $X$ and $Y$ are jointly continuous if the probabilities of events involving $(X, Y)$ can be expressed as an integral of a pdf. There is a nonnegative function $f_{X, Y}(x, y)$, called the joint probability density function, that is defined on the real plane such that for every event $A$, a subset of the plane,

$$
\begin{equation*}
P[\mathbf{X} \text { in } A]=\int_{A} \int f_{X, Y}\left(x^{\prime}, y^{\prime}\right) d x^{\prime} d y^{\prime}, \tag{4.9}
\end{equation*}
$$

as shown in Fig. 4.7. When a is the entire plane, the integral must equal one :

$$
\begin{equation*}
1=\int_{-\infty}^{\infty} \int_{-\infty}^{\infty} f_{X, Y}\left(x^{\prime}, y^{\prime}\right) d x^{\prime} d y^{\prime} . \tag{4.10}
\end{equation*}
$$

The joint cdf can be obtained in terms of the joint pdf of jointly continuous random variables by integrating over the semi-infinite

FIGURE 4.7
The probability of $A$ is the integral of $f_{x},(x, y)$ over the region defined by $A$.


The marginal pdf's $f_{X}(x)$ and $f_{Y}(y)$ are obtained by taking the derivative of the corresponding marginal cdf's

$$
\begin{align*}
F_{X}(x) & =F_{X, Y}(x, \infty) \\
F_{Y}(y) & =F_{X, Y}(\infty, y) . \\
F_{X}(x) & =\frac{d}{d x} \int_{-\infty}^{x}\left\{\int_{-\infty}^{\infty} f_{X, Y}\left(x^{\prime}, y^{\prime}\right) d y^{\prime}\right\} d x^{\prime} \\
& =\int_{-\infty}^{\infty} f_{X, Y}\left(x, y^{\prime}\right) d y^{\prime} . \tag{4.15a}
\end{align*}
$$

$$
\begin{equation*}
F_{Y}(y)=\int_{-\infty}^{\infty} f_{X, Y}\left(x^{\prime}, y\right) d x^{\prime} \tag{4.15b}
\end{equation*}
$$

## INDEPENDENCE OF TWO RANDOM VARIABLES

$X$ and $Y$ are independent random variables if any event $A_{1}$ defined in terms of $X$ is independent of any event $A_{2}$ defined in terms of $Y$;

$$
\begin{equation*}
P\left[X \text { in } A_{1}, Y \text { in } A_{2}\right]=P\left[X \text { in } A_{1}\right] P\left[Y \text { in } A_{2}\right] \tag{4,17}
\end{equation*}
$$

Suppose that $X$ and $Y$ are a pair of discrete random variables. If we let

$$
A_{1}=\left\{X=x_{j}\right\} \text { and } A_{2}=\left\{Y=y_{k}\right\} \text { then the independence of } X \text { and } Y
$$ implies that

$$
\begin{align*}
p_{X, Y}\left(x_{j}, y_{k}\right) & =P\left[X=x_{j}, Y=y_{k}\right] \\
& =P\left[X=x_{j}\right] P\left[Y=y_{k}\right] \\
& =p_{X}\left(x_{j}\right) p_{Y}\left(y_{k}\right) \quad \text { for all } x_{j} \text { and } y_{k} . \tag{4.18}
\end{align*}
$$

### 4.4 CONDITIONAL PROBABILITY AND CONDITIONAL EXPECTATION

## Conditional Probability

In Section 2.4, we know

$$
\begin{equation*}
P[Y \text { in } A \mid X=x]=\frac{P[Y \text { in } A, X=x]}{P[X=x]} . \tag{4.22}
\end{equation*}
$$

If $X$ is discrete, then Eq. (4.22) can be used to obtain the conditional cdf of $Y$ given $X=x_{k}$ :

$$
\begin{equation*}
F_{Y}\left(y \mid x_{k}\right)=\frac{P\left[Y \leq y, X=x_{k}\right]}{P\left[X=x_{k}\right]}, \text { for } P\left[X=x_{k}\right]>0 \tag{4.23}
\end{equation*}
$$

The conditional pdf of $Y$ given $X=x_{k}$, if the derivative exists, is given by $f_{Y}\left(y \mid x_{k}\right)=\frac{d}{d y} F_{Y}\left(y \mid x_{k}\right)$.

## MULTIPLE RANDOM VARIABLES

## Joint Distributions

The joint cumulative distribution function of $X_{1}, X_{2}, \ldots, X_{n}$ is defined as the probability of an $n$-dimensional semi-infinite rectangle associate with the point $\left(x_{1}, \ldots, x_{\mathrm{n}}\right)$ :

$$
\begin{equation*}
F_{X_{1}, X_{2}, \ldots X_{n}}\left(x_{1}, x_{2}, \ldots x_{n}\right)=P\left[X_{1} \leq x_{1}, X_{2} \leq x_{2}, \ldots, X_{n} \leq x_{n}\right] . \tag{4.38}
\end{equation*}
$$

The joint cdf is defined for discrete, continuous, and random variables of mixed type

## FUNCTIONS OF SEVERAL RANDOM VARIABLES

## One Function of Several Random Variables

Let the random variable $Z$ be defined as a function of several random variables:

$$
\begin{equation*}
Z=g\left(X_{1}, X_{2}, \ldots, X_{n}\right) . \tag{4.51}
\end{equation*}
$$

The cdf of $Z$ is found by first finding the equivalent event of that is, the set $R_{Z}=\left\{\mathrm{x}=\left(x_{1}, \ldots, x_{n}\right)\right.$ such that $\left.g(\mathrm{x}) \leq z\right\}$, then

$$
\begin{align*}
F_{Z}(z) & =P\left[\mathrm{X} \text { in } R_{z}\right] \\
& =\int_{\mathrm{xin} R_{z}}^{\ldots} \int f_{X_{1}, \ldots, X_{n}}\left(x_{1}^{\prime}, \ldots, x_{n}^{\prime}\right) d x_{1}^{\prime} \ldots d x_{n}^{\prime} . \tag{4.52}
\end{align*}
$$

EXAMPLE 4.31 Sum of Two Random Variables
Let $Z=X+Y$. Find $F_{Z}(z)$ and $f_{Z}(z)$ in terms of the joint pdf of $X$ and $Y$.

The cdf of $Z$ is

$$
F_{Z}(z)=\int_{-\infty}^{\infty} \int_{-\infty}^{z-x^{\prime}} f_{X, Y}\left(x^{\prime}, y^{\prime}\right) d y^{\prime} d x^{\prime}
$$

The pdf of $Z$ is

$$
\begin{equation*}
f_{Z}(z)=\frac{d}{d z} F_{Z}(z)=\int_{-\infty}^{\infty} f_{X, Y}\left(x^{\prime}, z-x^{\prime}\right) d x^{\prime} \tag{4.53}
\end{equation*}
$$

Thus the pdf for the sum of two random variables is given by a superposition integral.

If $X$ and $Y$ are
independent random variables, then by Eq. (4.21) the pdf is given by the convolution integral of the margial pdf's of $X$ and $Y$ :

$$
\begin{equation*}
f_{Z}(z)=\int_{-\infty}^{\infty} f_{X}\left(x^{\prime}\right) f_{Y}\left(z-x^{\prime}\right) d x^{\prime} \tag{4.54}
\end{equation*}
$$

## pdf of Linear Transformations

We consider first the linear transformation of two random variables

$$
\begin{aligned}
& V=a X+b Y \\
& W=c X+e Y
\end{aligned}
$$

$$
\left[\begin{array}{l}
V \\
W
\end{array}\right]=\left[\begin{array}{ll}
a & b \\
c & e
\end{array}\right]\left[\begin{array}{l}
X \\
Y
\end{array}\right]
$$

Denote the above matrix by $A$. We will assume $A$ has an inverse, so each point $(v, w)$ has a unique corresponding point $(x, y)$ obtained from

$$
\left[\begin{array}{l}
x  \tag{4.56}\\
y
\end{array}\right]=A^{-1}\left[\begin{array}{l}
v \\
w
\end{array}\right] .
$$

In Fig. 4.15, the infinitesimal rectangle and the parallelogram are equivalent events, so their probabilities must be equal. Thus

$$
f_{X, Y}(x, y) d x d y \cong f_{V, W}(v, w) d P
$$

FIGURE 4.15
Imageoi an infinitesimal rectangle undera linear transomation.

where $d P$ is the area of the parallelogram. The joint pdf of $V$ and $W$ is thus given by

$$
\begin{equation*}
f_{V, W}(v, w)=\frac{f_{X, Y}(x, y)}{\left|\frac{d P}{d x d y}\right|}, \tag{4.57}
\end{equation*}
$$

where $x$ an $y$ are related to ( $v, w$ ) by Eq. (4.56)
It can be shown tbat $=(|a e-b c|) d x d y$, so the "stretch factor" is

$$
\left|\frac{d P}{d x d y}\right|=\frac{|a e-b c|(d x d y)}{(d x d y)}=|a e-b c|=|A|
$$

where $|A|$ is the determinant of $A$.
Let the n -dimensional vector $\mathbf{Z}$ be

$$
\mathbf{Z}=A \mathbf{X}
$$

where $A$ is an $n \times n$ invertible matrix. The joint of $\mathbf{Z}$ is then

## EXPECTED VALUE OF FUNCTIONS OF RANDOM VARIABLES

The expected value of $Z=g(X, Y)$ can be found using the following expressions

$$
E[Z]= \begin{cases}\int_{-\infty}^{\infty} \int_{-\infty}^{\infty} g(x, y) f_{X, Y}(x, y) & X, Y \text { jointly continuous }  \tag{4.64}\\ \sum_{i} \sum_{n} g\left(x_{i}, y_{n}\right) p_{X, Y}\left(x_{i}, y_{n}\right) & X, Y \text { discrete } .\end{cases}
$$

## *Joint Characteristic Function

The joint characteristic function of n random variables is defined as

$$
\begin{align*}
& \Phi_{X_{1}, X_{2}, \ldots X_{n}}\left(w_{1}, w_{2}, \ldots w_{n}\right)=E\left[e^{j\left(w_{1} X_{1}+w_{2} X_{2}+\cdots+w_{n} X_{n}\right)}\right] .  \tag{4.73a}\\
& \Phi_{X, Y}\left(w_{1}, w_{2}\right)=E\left[e^{j\left(w_{1} X+w_{2} Y\right)}\right] \tag{4.73b}
\end{align*}
$$

If $X$ and $Y$ are jointly continuous random variables, then

$$
\begin{equation*}
\Phi_{X, Y}\left(w_{1}, w_{2}\right)=\int_{-\infty}^{\infty} \int_{-\infty}^{\infty} f_{X, Y}(x, y) e^{j\left(w_{1} x+w_{2}, y\right)} d x d y \tag{4.73c}
\end{equation*}
$$

The inversion formula for the Fourier transform implies that the joint pdf is given by

$$
\begin{equation*}
f_{X, Y}(x, y)=\frac{1}{4 \pi^{2}} \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} \Phi_{X, Y}\left(w_{1}, w_{2}\right) e^{j\left(w_{1} x+w_{2} y\right)} d w_{1} d w_{2} . \tag{4.74}
\end{equation*}
$$

## JOINTLY GAUSSIAN RANDOM VARIABLES

The random variables $X$ and $Y$ are said to be jointly Gaussian if their joint pdf has the form

$$
\begin{align*}
& f_{X, Y}(x, y) \\
& =\frac{\exp \left\{\frac{-1}{2\left(1-\rho_{X, Y}^{2}\right)}\left[\left(\frac{x-m_{1}}{\sigma_{1}}\right)^{2}-2 \rho_{X, Y}\left(\frac{x-m_{1}}{\sigma_{1}}\right)\left(\frac{y-m_{2}}{\sigma_{2}}\right)+\left(\frac{y-m_{2}}{\sigma_{2}}\right)^{2}\right]\right\}}{2 \pi \sigma_{1} \sigma_{2} \sqrt{1-\rho_{X, Y}^{2}}} \tag{4.79}
\end{align*}
$$

$-\infty<x<\infty$ and $-\infty<y<\infty$
The pdf is constant for values $x$ and $y$ for which the argument of the exponent is constant

$$
\left[\left(\frac{x-m_{1}}{\sigma_{1}}\right)^{2}-2 \rho_{X, Y}\left(\frac{x-m_{1}}{\sigma_{1}}\right)\left(\frac{y-m_{2}}{\sigma_{2}}\right)+\left(\frac{y-m_{2}}{\sigma_{2}}\right)^{2}\right]=\text { constant }
$$

When $\rho_{X, Y}=0, X$ and $Y$ are independent ; when $\rho_{X, Y} \neq 0$, the major axis of the ellipse is oriented along the angle

$$
\begin{equation*}
\theta=\frac{1}{2} \arctan \left(\frac{2 \rho_{X, Y} \sigma_{1} \sigma_{2}}{\sigma_{1}^{2}-\sigma_{2}^{2}}\right) . \tag{4.80}
\end{equation*}
$$

Note that the angle is $45^{\circ}$ when the variance are equal.
The marginal pdf of $X$ is found by integrating $f_{X, Y}(x, y)$ over all $y$

$$
\begin{equation*}
f_{X}(x)=\frac{e^{-\left(x-m_{1}\right)^{2} / 2 \sigma_{1}^{2}}}{\sqrt{2 \pi} \sigma_{1}} \tag{4.81}
\end{equation*}
$$

that is, $X$ is a Gaussian random variable with mean $m_{1}$ and variance

$$
\sigma_{1}^{2}
$$

## FIGURE 4.19

Orientation of contours of equal value of joint Gaussian pdf for $p_{x, y}>0$.


(c)

$$
o_{1}<o_{2}
$$

$$
\frac{\pi}{4}<\theta<\frac{\pi}{2}
$$

## n Jointly Gaussian Random Variables

The random variables $X_{1}, X_{2}, \ldots, X_{n}$ are said to be jointly Gaussian if their joint pdf is given by

$$
\begin{equation*}
f_{\mathbf{x}}(\mathbf{x}) \equiv f_{X_{1}, X_{2}, \ldots, X_{n}}\left(x_{1}, x_{2}, \ldots x_{n}\right)=\frac{\exp \left\{-\frac{1}{2}(\mathbf{x}-\mathbf{m})^{1} K^{-1}(\mathbf{x}-\mathbf{m})\right\}}{(2 \pi)^{n / 2}|k|^{1 / 2}}, \tag{4.83}
\end{equation*}
$$

where $\mathbf{x}$ and $\mathbf{m}$ are column vectors defined by

$$
\mathbf{x}=\left[\begin{array}{c}
x_{1} \\
x_{2} \\
\vdots \\
x_{n}
\end{array}\right], \quad \mathbf{m}=\left[\begin{array}{c}
m_{1} \\
m_{2} \\
\vdots \\
m_{n}
\end{array}\right]=\left[\begin{array}{c}
E\left[X_{1}\right] \\
E\left[X_{2}\right] \\
E\left[X_{3}\right] \\
E\left[X_{4}\right]
\end{array}\right]
$$

and K is the covariance matrix that is defined by

$$
K=\left[\begin{array}{cccc}
\operatorname{VAR}\left(X_{1}\right) & \operatorname{COV}\left(X_{2}, X_{1}\right) & \cdots & \operatorname{Cov}\left(X_{1}, X_{n}\right)  \tag{4.84}\\
\operatorname{Cov}\left(X_{2}, X_{1}\right) & \operatorname{VAR}\left(X_{2}\right) & \cdots & \operatorname{cov}\left(X_{2}, X_{n}\right) \\
\vdots & \vdots & & \vdots \\
\operatorname{Cov}\left(X_{n}, X_{1}\right) & \cdots & & \operatorname{VAR}\left(X_{n}\right)
\end{array}\right]
$$

## Transformations of Random Vectors

Let $X_{1}, \ldots, X_{n}$ be random variables associate with some experiment, and let the random variables $Z_{1}, \ldots, Z_{n}$ be defined by n functions of $\mathbf{X}=\left(X_{1}, \ldots, X_{n}\right)$ :

$$
Z_{1}=g_{1}(\mathrm{X}) \quad Z_{2}=g_{2}(\mathrm{X}) \quad \ldots \quad Z_{n}=g_{n}(\mathrm{X}) .
$$

The joint cdf of $Z_{1}, \ldots, Z_{n}$ at the point $\mathbf{z}=(z 1, \ldots, z n)$ is equal to the probability of the region of $\mathbf{x}$ where

$$
\begin{align*}
& F_{Z_{1}, \ldots, Z_{n}}\left(z_{1}, \ldots, z_{n}\right)=P\left[g_{1}(\mathrm{X}) \leq z_{1}, \ldots, g_{n}(\mathrm{X}) \leq z_{n}\right] .  \tag{4.55a}\\
& F_{Z_{1}, \ldots, Z_{n}}\left(z_{1}, \ldots, z_{n}\right)=\int_{x^{\prime}: g_{k}\left(\mathrm{x}^{\prime}\right) \leq z_{k}} \cdots \int f_{X_{1}, \ldots, X_{n}}\left(x_{1}^{\prime}, \ldots, x_{n}^{\prime}\right) d x_{1}^{\prime} \cdots d x_{n}^{\prime} . \tag{4.55b}
\end{align*}
$$

## pdf of Linear Transformations

We consider first the linear transformation of two random variables

$$
\begin{aligned}
& V=a X+b Y \\
& W=c X+e Y
\end{aligned}
$$

$$
\left[\begin{array}{l}
V \\
W
\end{array}\right]=\left[\begin{array}{ll}
a & b \\
c & e
\end{array}\right]\left[\begin{array}{l}
X \\
Y
\end{array}\right] .
$$

Denote the above matrix by $A$. We will assume $A$ has an inverse, so each point $(v, w)$ has a unique corresponding point $(x, y)$ obtained from

$$
\left[\begin{array}{l}
x  \tag{4.56}\\
y
\end{array}\right]=A^{-1}\left[\begin{array}{l}
v \\
w
\end{array}\right] .
$$

In Fig. 4.15, the infinitesimal rectangle and the parallelogram are equivalent events, so their probabilities must be equal. Thus

$$
f_{X, Y}(x, y) d x d y \cong f_{V, W}(v, w) d P
$$

FIGURE 4.15
Imageoi an infinitesimal rectangle undera linear transomation.


## Stochastic Processes

Let $\xi$ denote the random outcome of an experiment. To every such outconme suppose a waveform
$X(t, \xi)$ is assigned.
The collection of such waveforms form a stochastic process. The set of $\left\{\xi_{k}\right\}$ and the time index $t$ can be continuous or discrete (countably infinite or finite) as well.
For fixed $\xi_{i} \in S$ (the set of


Fig. 14.1 all experimental outcomes), $X(t, \xi$;s a specific time function.
For fixed $t$,

$$
X_{1}=X\left(t_{1}, \xi_{i}\right)
$$

is a random variable. The ensemble of all such realizations $X(t, \xi)$ over time represents the stochastic
process $X(t)$. (see Fig 14.1). For example

$$
X(t)=a \cos \left(\omega_{0} t+\varphi\right)
$$

If $X(t)$ is a stochastic process, then for fixed $t, X(t)$ represents a random variable. Its distribution function is given by

$$
F_{x}(x, t)=P\{X(t) \leq x\}
$$

Notice that $F_{X}(x, t)$ depends on $t$, since for a different $t$, we obtain a different random variable. Further

$$
f_{x}^{\text {ther }}(x, t)=\frac{d F_{X}(x, t)}{d x}
$$

represents the first-order probability density function of the process $X(t)$.

For $t=t_{1}$ and $t=t_{2}, X(t)$ represents two different random variables $X_{1}=X\left(t_{1}\right)$ and $X_{2}=X\left(t_{2}\right)$ respectively. Their joint distribution is given by

$$
F_{x}\left(x_{1}, x_{2}, t_{1}, t_{2}\right)=P\left\{X\left(t_{1}\right) \leq x_{1}, X\left(t_{2}\right) \leq x_{2}\right\}
$$

and

$$
f_{x}\left(x_{1}, x_{2}, t_{1}, t_{2}\right)=\frac{\partial^{2} F_{x}\left(x_{1}, x_{2}, t_{1}, t_{2}\right)}{\partial x_{1} \partial x_{2}}
$$

represents the second-order density function of the process $X(t)$. Similarly $f_{x}\left(x_{1}, x_{2}, \cdots x_{n}, t_{1}, t_{2} \cdots, t_{\text {伴 }}\right.$ bresents the $\mathrm{n}^{\text {th }}$ order density function of the process $X(t)$. Complete specification of the stochastic process $X(t)$ requires the knowledge of $f_{x}\left(x_{1}, x_{2}, \cdots x_{n}, t_{1}, t_{2} \cdots, t_{n}\right)$ for alt $i_{i}, \quad i=1,2, \cdots$, nand for all $n$. (an almost impossible task in reality).

## Mean of a Stochastic Process:

$$
\mu(t)=E\{X(t)\}=\int_{-\infty}^{+\infty} x f_{x}(x, t) d x
$$

represents the mean value of a process $X(t)$. In general, the mean of a process can depend on the time index $t$.

Autocorrelation function of a process $X(t)$ is defined as

$$
R_{x x}\left(t_{1}, t_{2}\right)=E\left\{X\left(t_{1}\right) X^{*}\left(t_{2}\right)\right\}=\iint x_{1} x_{2}^{*} f_{x}\left(x_{1}, x_{2}, t_{1}, t_{2}\right) d x_{1} d x_{2}
$$

and it represents the interrelationship between the random variables $X_{1}=X\left(t_{1}\right)$ and $X_{2}=X\left(t_{2}\right)$ generated from the process $X(t)$.

## Properties:

1. $R_{x x}\left(t_{1}, t_{2}\right)=R_{x x}^{*}\left(t_{2}, t_{1}\right)=\left[E\left\{X\left(t_{2}\right) X^{*}\left(t_{1}\right)\right\}\right]^{*}$
2. $R_{x x}(t, t)=E\left\{|X(t)|^{2}\right\}>0$.
3. $R_{\text {sefo }}\left(t_{1}, t_{2}\right.$, represents a nonnegative definite function, i.e., for any set'of constants

$$
\begin{equation*}
\left\{a_{i}\right\}_{i=1}^{n} \tag{14-8}
\end{equation*}
$$

Eq. (14-8) follows by noticing that $\sum_{i=1}^{n} \sum_{j=1}^{n} a_{i} a_{j}^{*} R_{x X}\left(t_{i}, t_{j}\right) \geq 0$.
The function

$$
E\left\{|Y|^{2}\right\} \geq 0 \text { for } Y=\sum_{i=1}^{n} a_{i} X\left(t_{i}\right)
$$

represents the autocovariance function of the process $X(t)$.

$$
\begin{aligned}
& \text { Example 14.1 } \left.C_{x x}\left(t_{1}, t_{2}\right)=R_{x x}\left(t_{1}, t_{2}\right)-\mu_{x}\left(t_{1}\right) \mu_{x}^{*}\left(t_{2}\right)\right)
\end{aligned}
$$

Then

$$
\begin{gather*}
z=\int_{-T}^{T} X(t) d t \\
E\left[|z|^{2}\right]= \\
=\int_{-T}^{T} \int_{-T}^{T} E\left\{X\left(t_{1}\right) X^{*}\left(t_{2}\right)\right\} d t_{1} d t_{2}  \tag{14-10}\\
=
\end{gather*}
$$

## Stationary Stochastic Processes

Stationary processes exhibit statistical properties that are invariant to shift in the time index. Thus, for example, second-order stationarity implies that the statistical properties of the pairs $\left\{X\left(t_{1}\right), X\left(t_{2}\right)\right\}$ and $\left\{X\left(t_{1}+c\right), X\left(t_{2}+c\right)\right\}$ are the same for any $c$.
Similarly first-order stationarity implies that the statistical properties of $X\left(t_{i}\right)$ and $X\left(t_{i}+c\right)$ are the same for any $c$.

In strict terms, the statistical properties are governed by the joint probability density function. Hence a process is $\mathrm{n}^{\text {th }}$-order

## Strict-Sense Stationary (S.S.S) if

$$
f_{x}\left(x_{1}, x_{2}, \cdots x_{n}, t_{1}, t_{2} \cdots, t_{n}\right) \equiv f_{x}\left(x_{1}, x_{2}, \cdots x_{n}, t_{1}+c, t_{2}+c \cdots, t_{n}+c\right)
$$

for any $c$, where the left side represents the joint density function of the random variables $X_{1}=X\left(t_{1}\right), X_{2}=X\left(t_{2}\right), \cdots, X_{n}=X\left(t_{n}\right) \quad$ and the right side corresponds to the joint density function of the random variables $X_{1}^{\prime}=X\left(t_{1}+c\right), X_{2}^{\prime}=X\left(t_{2}+c\right), \cdots, X_{n}^{\prime}=X\left(t_{n}+c\right)$.
A process $X(t)$ is said to be strict-sense stationary if $(14-14)$ is true for all $t_{i}, \quad i=1,2, \cdots, n, \quad n=1,2, \cdots$ and any $c$.

For a first-order strict sense stationary process, from (14-14) we have

$$
\begin{equation*}
f_{x}(x, t) \equiv f_{x}(x, t+c) \tag{14-15}
\end{equation*}
$$

for any $c$. In particular $c=-t$ gives

$$
\begin{equation*}
f_{x}(x, t)=f_{x}(x) \tag{14-16}
\end{equation*}
$$

i.e., the first-order density of $X(t)$ is independent of $t$. In that case

Similarly, for a second-order strict-sense stationary process we have from $\left.\underset{E}{(14-14)} X^{4}(t)\right]=\int_{-\infty}^{+\infty} x f(x) d x=\mu$, a constant.
for any $c$. For $c=-t_{2}$ we get

$$
\begin{gather*}
f_{X}\left(x_{1}, x_{2}, t_{1}, t_{2}\right) \equiv f_{x}\left(x_{1}, x_{2}, t_{1}+c, t_{2}+c\right) \\
f_{X}\left(x_{1}, x_{2}, t_{1}, t_{2}\right) \equiv f_{x}\left(x_{1}, x_{2}, t_{1}-t_{2}\right) \tag{14-18}
\end{gather*}
$$

i.e., the second order density function of a strict sense stationary process depends only on the difference of the time indices In that case the autocorrelation function is given by

$$
t_{1}-t_{2}=\tau
$$

$$
\begin{align*}
R_{x x}\left(t_{1}, t_{2}\right) & =E\left\{X\left(t_{1}\right) X^{*}\left(t_{2}\right)\right\} \\
& \pm \iint x_{1} x_{2}^{*} f_{x}\left(x_{1}, x_{2}, \tau=t_{1}-t_{2}\right) d x_{1} d x_{2}  \tag{14-19}\\
& =R_{x x}\left(t_{1}-t_{2}\right)=R_{x x}(\tau)=R_{x x}^{*}(-\tau)
\end{align*}
$$

i.e., the autocorrelation function of a second order strict-sense stationary process depends only on the difference of the time indices $\tau=t_{1}-t_{2}$.
Notice that (14-17) and (14-19) are consequences of the stochastic process being first and second-order strict sense stationary.
On the other hand, the basic conditions for the first and second order stationarity - Eqs. (14-16) and (14-18) - are usually difficult to verify. In that case, we often resort to a looser definition of stationarity, known as Wide-Sense Stationarity (W.S.S), by making use of
(14-17) and (14-19) as the necessary conditions. Thus, a process $X(t)$ is said to be Wide-Sense Stationary if
and $E\{X(t)\}=\mu$

$$
\begin{equation*}
E\left\{X\left(t_{1}\right) X^{*}\left(t_{2}\right)\right\}=R_{x x}\left(t_{1}-t_{2}\right), \tag{14-20}
\end{equation*}
$$

i.e., for wide-sense stationary processes, the mean is a constant and the autocorrelation function depends only on the difference between the time indices. Notice that (14-20)-(14-21) does not say anything about the nature of the probability density functions, and instead deal with the average behavior of the process. Since (14-20)-(14-21) follow from (14-16) and (14-18), strict-sense stationarity always implies wide-sense stationarity. However, the converse is not true in general, the only exception being the Gaussian process.
This follows, since if $X(t)$ is a Gaussian process, then by definition $X_{1}=X\left(t_{1}\right), X_{2}=X\left(t_{2}\right), \cdots, X_{n}=X\left(t_{n}\right)$ are jointly Gaussian random variables for any $t_{1}, t_{2} \cdots, t_{n}$ whose joint characteristic function is given by

$$
\begin{equation*}
\phi_{\underline{x}}\left(\omega_{1}, \omega_{2}, \cdots, \omega_{n}\right)=e^{j \sum_{k=1}^{n} \mu\left(t_{k}\right) \omega_{k}-\sum_{l, k}^{n} \sum C_{x x}\left(t_{i}, t_{k}\right) \omega_{i} \omega_{k} / 2} \tag{14-22}
\end{equation*}
$$

where $C_{X X}\left(t_{i}, t_{k}\right)$ is as defined on (14-9). If $X(t)$ is wide-sense stationary, then using (14-20)-(14-21) in (14-22) we get

$$
\begin{equation*}
\phi_{\underline{X}}\left(\omega_{1}, \omega_{2}, \cdots, \omega_{n}\right)=e^{j \sum_{k=1}^{n} \mu \omega_{k}-\frac{1}{2} \sum_{\mathrm{l}=1}^{n} \sum_{k=1}^{n} C_{x x}\left(t_{i}-t_{k}\right) \omega_{i} \omega_{k}} \tag{14-23}
\end{equation*}
$$

and hence if the set of time indices are shifted by a constant $c$ to generate a new set of jointly Gaussian random variables $X_{1}^{\prime}=X\left(t_{1}+c\right)$, $X_{2}^{\prime}=X\left(t_{2}+c\right), \cdots, X_{n}^{\prime}=X\left(t_{n}+c\right) \quad$ then their joint characteristic function is identical to (14-23). Thus the set of random variables and $\left\{X_{i}^{\prime}\right\}_{i=1}^{n}$ have the same joint probability distribution for all $n$ and $\left\{X_{i}\right\}_{i=1}^{n}$ all $c$, establishing the strict sense stationarity of Gaussian processes from its wide-sense stationarity.

To summarize if $X(t)$ is a Gaussian process, then wide-sense stationarity (w.s.s) $\Rightarrow$ strict-sense stationarity (s.s.s). Notice that since the joint p.d.f of Gaussian random variables depends only on their second order statistics, which is also the basis

## Systems with Stochastic Inputs

A deterministic system ${ }^{1}$ transforms each input waveform $X\left(t, \xi_{i}\right)$ nto an output waveform $Y\left(t, \xi_{i}\right)=T\left[X\left(t, \xi_{i}\right)\right]$ by operating only on the time variable $t$. Thus a set of realizations at the input corresponding to a process $X(t)$ generates a new set of realizations $\{Y(t, \xi)\}$ at the output associated with a new process $Y(t)$.


Fig. 14.3

Our goal is to study the output process statistics in terms of the input process statistics and the system function.
${ }^{1}$ A stochastic system on the other hand operates on both the variables $t$ and $\xi$.

## Linear Systems: $L[\cdot]$ represents a linear system if

$$
\begin{equation*}
L\left\{a_{1} X\left(t_{1}\right)+a_{2} X\left(t_{2}\right)\right\}=a_{1} L\left\{X\left(t_{1}\right)\right\}+a_{2} L\left\{X\left(t_{2}\right)\right\} . \tag{14-28}
\end{equation*}
$$

Let

$$
\begin{equation*}
Y(t)=L\{X(t)\} \tag{14-29}
\end{equation*}
$$

represent the output of a linear system.
Time-Invariant System: $L[\cdot]$ represents a time-invariant system if

$$
\begin{equation*}
Y(t)=L\{X(t)\} \Rightarrow L\left\{X\left(t-t_{0}\right)\right\}=Y\left(t-t_{0}\right) \tag{14-30}
\end{equation*}
$$

i.e., shift in the input results in the same shift in the output also.

If $L[\cdot]$ satisfies both (14-28) and (14-30), then it corresponds to a linear time-invariant (LTI) system.
LTI systems can be uniquely represented in terms of their output to a delta function

then



Fig. 14.6


$$
Y(t)=\int_{-\infty}^{+\infty} h(t-\tau) X(\tau) d \tau
$$

$$
=\int_{-\infty}^{+\infty} h(\tau) X(t-\tau) d \tau
$$

Eq. (14-31) follows by expressing $X(t)$ as

$$
\begin{equation*}
X(t)=\int_{-\infty}^{+\infty} X(\tau) \delta(t-\tau) d \tau \tag{14-32}
\end{equation*}
$$

and applying (14-28) and (14-30) to $Y(t)=L\{X(t)\}$ Thus

$$
\begin{array}{rlrl}
Y(t) & =L\{X(t)\}=L\left\{\int_{-\infty}^{+\infty} X(\tau) \delta(t-\tau) d \tau\right\} \\
& =\int_{-\infty}^{+\infty} L\{X(\tau) \delta(t-\tau) d \tau\} \quad \text { By Linearity } \\
& =\int_{-\infty}^{+\infty} X(\tau) L\{\delta(t-\tau)\} d \tau & \text { By Time-invariance } \\
& =\int_{-\infty}^{+\infty} X(\tau) h(t-\tau) d \tau=\int_{-\infty}^{+\infty} h(\tau) X(t-\tau) d \tau .(14-33)
\end{array}
$$

## Output Statistics: Using (14-33), the mean of the output process

 is given by$$
\begin{align*}
\mu_{\gamma}(t) & =E\{Y(t)\}=\int_{-\infty}^{+\infty} E\{X(\tau) h(t-\tau) d \tau\} \\
& =\int_{-\infty}^{+\infty} \mu_{x}(\tau) h(t-\tau) d \tau=\mu_{x}(t) * h(t) \tag{14-34}
\end{align*}
$$

Similarly the cross-correlation function between the input and output processes is given by

$$
\begin{align*}
R_{X Y}\left(t_{1}, t_{2}\right) & =E\left\{X\left(t_{1}\right) Y^{*}\left(t_{2}\right)\right\} \\
& =E\left\{X\left(t_{1}\right) \int_{-\infty}^{+\infty} X^{*}\left(t_{2}-\alpha\right) h^{*}(\alpha) d \alpha\right\} \\
& =\int_{-\infty}^{+\infty} E\left\{X\left(t_{1}\right) X^{*}\left(t_{2}-\alpha\right)\right\} h^{*}(\alpha) d \alpha \\
& =\int_{-\infty}^{+\infty} R_{X X}\left(t_{1}, t_{2}-\alpha\right) h^{*}(\alpha) d \alpha  \tag{14-35}\\
& =R_{X X}\left(t_{1}, t_{2}\right) * h^{*}\left(t_{2}\right)
\end{align*}
$$

Finally the output autocorrelation function is given by

$$
\begin{align*}
R_{r y}\left(t_{1}, t_{2}\right) & =E\left\{Y\left(t_{1}\right) Y^{*}\left(t_{2}\right)\right\} \\
& =E\left\{\int_{-\infty}^{+\infty} X\left(t_{1}-\beta\right) h(\beta) d \beta Y^{*}\left(t_{2}\right)\right\} \\
& =\int_{-\infty}^{+\infty} E\left\{X\left(t_{1}-\beta\right) Y^{*}\left(t_{2}\right)\right\} h(\beta) d \beta \\
& =\int_{-\infty}^{+\infty} R_{x v}\left(t_{1}-\beta, t_{2}\right) h(\beta) d \beta \\
& =R_{x \gamma}\left(t_{1}, t_{2}\right) * h\left(t_{1}\right), \tag{14-36}
\end{align*}
$$

or

$$
\begin{equation*}
R_{Y Y}\left(t_{1}, t_{2}\right)=R_{X X}\left(t_{1}, t_{2}\right) * h^{*}\left(t_{2}\right) * h\left(t_{1}\right) \tag{14-37}
\end{equation*}
$$


(a)


Fig. 14.7

In particular if $X(t)$ is wide-sense stationary, then we have so that from (14-34)

$$
\begin{equation*}
\mu_{Y}(t)=\mu_{x} \int_{-\infty}^{+\infty} h(\tau) d \tau=\mu_{X} c, \quad a \text { constant } \tag{14-38}
\end{equation*}
$$

$\mu_{X}(t)=\mu_{X}$

Also $R_{x x}\left(t_{1}, t_{2}\right)=R_{x x}\left(t_{1}-t_{2}\right)$ so that (14-35) reduces to

$$
\begin{align*}
R_{X Y}\left(t_{1}, t_{2}\right) & =\int_{-\infty}^{+\infty} R_{X X}\left(t_{1}-t_{2}+\alpha\right) h^{*}(\alpha) d \alpha  \tag{14-39}\\
& =R_{x X}(\tau) * h^{*}(-\tau)=R_{X Y}(\tau), \quad \tau=t_{1}-t_{2}
\end{align*}
$$

Thus $X(t)$ and $Y(t)$ are jointly w.s.s. Further, from (14-36), the output autocorrelation simplifies to

$$
\begin{align*}
R_{Y Y}\left(t_{1}, t_{2}\right) & \Delta \int_{-\infty}^{+\infty} R_{X Y}\left(t_{1}-\beta-t_{2}\right) h(\beta) d \beta, \quad \tau=t_{1}-t_{2} \\
& =R_{X Y}(\tau) * h(\tau)=R_{Y Y}(\tau) . \tag{14-40}
\end{align*}
$$

From (14-37), we obtain

$$
\begin{equation*}
R_{r y}(\tau)=R_{x x}(\tau) * h^{*}(-\tau) * h(\tau) . \tag{14-41}
\end{equation*}
$$

From (14-38)-(14-40), the output process is also wide-sense stationary. This gives rise to the following representation


Fig. 14.8

## Discrete Time Stochastic Processes:

A discrete time stochastic process $X_{n}=X(n T)$ is a sequence of random variables. The mean, autocorrelation and auto-covariance functions of a discrete-time process are gives by

$$
\begin{equation*}
\mu_{n}=E\{X(n T)\} \tag{14-57}
\end{equation*}
$$

and

$$
\begin{equation*}
R\left(n_{1}, n_{2}\right)=E\left\{X\left(n_{1} T\right) X^{*}\left(n_{2} T\right)\right\} \tag{14-58}
\end{equation*}
$$

respectively. As before strict sense stationarity and wide-sense stationarity definitions apply here also.
For example, $X(n T)$ is wide sense stationary if

$$
\begin{equation*}
C\left(n_{1}, n_{2}\right)=R\left(n_{1}, n_{2}\right)-\mu_{n_{1}} \mu_{n_{2}}^{*} \tag{14-59}
\end{equation*}
$$

and

$$
\begin{gather*}
E\{X(n T)\}=\mu, \quad a \text { constant }  \tag{14-60}\\
E\left[X\{(k+n) T\} X^{*}\{(k) T\}\right]=R(n)=r_{n} \triangleq r_{-n}^{*} \tag{14-61}
\end{gather*}
$$

## Power Spectrum

For a deterministic signal $x(\mathrm{t})$, the spectrum is well defined: If $\quad X(\omega)$ represents its Fourier transform, i.e., if

$$
\begin{equation*}
X(\omega)=\int_{-\infty}^{+\infty} x(t) e^{-j \omega t} d t \tag{18-1}
\end{equation*}
$$

then $|X(\omega)|^{2}$ represents its energy spectrum. This follows from Parseval's theorem since the signal energy is given by

$$
\begin{equation*}
\int_{-\infty}^{+\infty} x^{2}(t) d t=\frac{1}{2 \pi} \int_{-\infty}^{+\infty}|X(\omega)|^{2} d \omega=E . \tag{18-2}
\end{equation*}
$$

$(\omega, \omega+\Delta \omega)$
Thus $|X(\omega)|^{2} \Delta \omega$ represents the signal energy in the band (see Fig 18.1).



Fig 18.1

However for stochastic processes, a direct application of (18-1) generates a sequence of random variables for every $\omega$. Moreover, for a stochastic process, $E\left\{|X(t)|^{2}\right\}$ represents the ensemble average power (instantaneous energy) at the instant $t$.

To obtain the spectral distribution of power versus frequency for stochastic processes, it is best to avoid infinite intervals to begin with, and start with a finite interval ( $-T, T$ ) in (18-1). Formally, partial Fourier transform of a process $X(t)$ based on $(-T, T)$ is given by
so that

$$
\begin{equation*}
X_{T}(\omega)=\int_{-T}^{T} X(t) e^{-j \omega t} d t \tag{18-3}
\end{equation*}
$$

represents the power distribution associated with that realization based on ( $-T, T$ ). Notice that (18-4) represents a random variable for every and its ensemble average gives, the average power distribution based on ( $-T, T$ ). Thus

$$
\begin{equation*}
\frac{\left|X_{T}(\omega)\right|^{2}}{2 T}=\frac{1}{2 T}\left|\int_{-T}^{T} X(t) e^{-j \omega t} d t\right|^{2} \quad \omega \tag{18-4}
\end{equation*}
$$

$$
P_{T}(\omega)=E\left\{\frac{\left|X_{T}(\omega)\right|^{2}}{2 T}\right\}=\frac{1}{2 T} \int_{-T}^{T} \int_{-T}^{T} E\left\{X\left(t_{1}\right) X^{*}\left(t_{2}\right)\right\} e^{-j \omega\left(t_{1}-t_{2}\right)} d t_{1} d t_{2}
$$

$$
\begin{equation*}
=\frac{1}{2 T} \int_{-T}^{T} \int_{-T}^{T} R_{X X}\left(t_{1}, t_{2}\right) e^{-j \omega\left(t_{1}-t_{2}\right)} d t_{1} d t_{2} \tag{18-5}
\end{equation*}
$$

represents the power distribution of $X(t)$ based on $(-T, T)$. For wide sense stationary (w.s.s) processes, it is possible to further simplify (18-5). Thus if $X(t)$ is assumed to be w.s.s, then and (18-5) simplifies to

$$
R_{x x}\left(t_{1}, t_{2}\right)=R_{x x}\left(t_{1}-t_{2}\right)
$$

Let $\tau=t_{1}-t_{2}$ and proceeding as in (14-24), we get

$$
P_{T}(\omega)=\frac{1}{2 T} \int_{-T}^{T} \int_{-T}^{T} R_{X X}\left(t_{1}-t_{2}\right) e^{-j \omega\left(t_{1}-t_{2}\right)} d t_{1} d t_{2}
$$

to be the power distribution of the w.s.s. process $X(t)$ based on $(-T, T)$. Finally letting $T \underset{\rightarrow}{\rightarrow} \infty$ in (18-6), we obtain

$$
\begin{align*}
P_{T}(\omega) & =\frac{1}{2 T} \int_{-2 T}^{2 T} R_{X X}(\tau) e^{-j \omega \tau}(2 T-|\tau|) d \tau  \tag{18-6}\\
& =\int_{-2 T}^{2 T} R_{x X}(\tau) e^{-j \omega \tau}\left(1-\frac{|\tau|}{2 T}\right) d \tau \geq 0
\end{align*}
$$

$$
\begin{equation*}
S_{x x}(\omega)=\lim _{T \rightarrow \infty} P_{T}(\omega)=\int_{-\infty}^{+\infty} R_{x x}(\tau) e^{-j \omega \tau} d \tau \geq 0 \tag{18-7}
\end{equation*}
$$

to be the power spectral density of the w.s.s process $X(t)$. Notice that

$$
\begin{equation*}
R_{x x}(\omega) \stackrel{\text { F.T }}{\longleftrightarrow} S_{x x}(\omega) \geq 0 \tag{18-8}
\end{equation*}
$$

i.e., the autocorrelation function and the power spectrum of a w.s.s Process form a Fourier transform pair, a relation known as the Wiener-Khinchin Theorem. From (18-8), the inverse formula gives
and in particular for $\tau=0$, we get

$$
R_{x x}(\tau)=\frac{1}{2 \pi} \int_{-\infty}^{+\infty} S_{x x}(\omega) e^{j \omega \tau} d \omega
$$

From (18-10), the area under $S_{x x}(\omega)$ represents the total power of the process $X(t)$, and hence $S_{X X}(\omega)$ truly represents the power spectrum $_{+\infty}$ (Fig 18.2).

$$
\frac{1}{2 \pi} \int_{-\infty}^{+\infty} S_{x x}(\omega) d \omega=R_{x x}(0)=E\left\{|X(t)|^{2}\right\}=P, \quad \text { the total ppwer }(18-10)
$$

If $X(t)$ is a real w.s.s process, then $R_{x x}(\tau)=R_{x x}(-\tau)$ so that

$$
\begin{aligned}
S_{x X}(\omega) & =\int_{-\infty}^{+\infty} R_{X X}(\tau) e^{-j \omega \tau} d \tau \\
& =\int_{-\infty}^{+\infty} R_{x X}(\tau) \cos \omega \tau d \tau \\
& =2 \int_{0}^{\infty} R_{x X}(\tau) \cos \omega \tau d \tau=S_{x X}(-\omega) \geq 0
\end{aligned}
$$

so that the power spectrum is an even function, (in addition to being real and nonnegative).

