# Lecture Slides for Signals and Systems 

## Edition 5.0 @ $\odot \Theta \Theta$

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## Part 0

## Preface

- This document constitutes a detailed set of lecture slides on signals and systems, covering both the continuous-time and discrete-time cases.
- These slides are organized in such a way as to facilitate the teaching of a course that covers:
- only the continuous-time case, or
- only the discrete-time case, or
- both the continuous-time and discrete-time cases.
- These slides are intended to be used in conjunction with the following textbook:
- M. D. Adams, Signals and Systems, Edition 5.0, Dec. 2022, ISBN 978-1-990707-00-1 (PDF). Available from Google Books, Google Play Books, and author's web site https://www.ece.uvic.ca/~mdadams/sigsys.book.


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## Companion Web Site

- The author of the lecture slides maintains a companion web site for the lecture slides and the associated textbook.
- The most recent version of the textbook and lecture slides can be downloaded from this site.
- Additional information related to the slides is also available from this site, including:
- errata for the slides; and
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## Video Lectures

- The author has prepared video lectures for some of the material covered in these slides and the associated textbook.
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- https://www.youtube.com/iamcanadian1867
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## Companion Git Repository

- These lecture slides have a companion Git repository.
- Numerous code examples are available from this repository.
- This repository is hosted by GitHub.
- The URL of the main repository page on GitHub is:
- https://github.com/mdadams/sigsysbook_companion
- The URL of the actual repository itself is:
- https://github.com/mdadams/sigsysbook_companion.git


## Part 1

## Introduction

- A signal is a function of one or more variables that conveys information about some (usually physical) phenomenon.
- For a function $f$, in the expression $f\left(t_{1}, t_{2}, \ldots, t_{n}\right)$, each of the $\left\{t_{k}\right\}$ is called an independent variable, while the function value itself is referred to as a dependent variable.
- Some examples of signals include:
- a voltage or current in an electronic circuit
- the position, velocity, or acceleration of an object
- a force or torque in a mechanical system
- a flow rate of a liquid or gas in a chemical process
- a digital image, digital video, or digital audio
- a stock market index


## Classification of Signals

- Number of independent variables (i.e., dimensionality):
- A signal with one independent variable is said to be one dimensional (e.g., audio).
- A signal with more than one independent variable is said to be multi-dimensional (e.g., image).
■ Continuous or discrete independent variables:
- A signal with continuous independent variables is said to be continuous time (CT) (e.g., voltage waveform).
- A signal with discrete independent variables is said to be discrete time (DT) (e.g., stock market index).
- Continuous or discrete dependent variable:
- A signal with a continuous dependent variable is said to be continuous valued (e.g., voltage waveform).
- A signal with a discrete dependent variable is said to be discrete valued (e.g., digital image).
- A continuous-valued CT signal is said to be analog (e.g., voltage waveform).
- A discrete-valued DT signal is said to be digital (e.g., digital audio).


## Graphical Representation of Signals



Continuous-Time (CT) Signal


Discrete-Time (DT) Signal

- A system is an entity that processes one or more input signals in order to produce one or more output signals.



## Classification of Systems

- Number of inputs:
- A system with one input is said to be single input (SI).
- A system with more than one input is said to be multiple input (MI).
- Number of outputs:
- A system with one output is said to be single output (SO).
- A system with more than one output is said to be multiple output (MO).
- Types of signals processed:
- A system can be classified in terms of the types of signals that it processes.
- Consequently, terms such as the following (which describe signals) can also be used to describe systems:
- one-dimensional and multi-dimensional,
- continuous-time (CT) and discrete-time (DT), and
- analog and digital.
- For example, a continuous-time (CT) system processes CT signals and a discrete-time (DT) system processes DT signals.

| Input |  | Discrete-Time |  | Discrete-TimeSignal |  | Output Continuous-Time Signal |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| Continuous-Time Signal | Continuous-to-Discrete-Time (C/D) Converter | Signal <br> Before Processing |  | After | Discrete-to-Continuous-Time (D/C) Converter |  |
|  |  |  | Discrete-Time System | Processing |  |  |
|  |  |  |  |  |  |  |

## Processing a Continuous-Time Signal With a Discrete-Time System

| Input | Continuous-Time |  |  | Continuous-Time |  | Output |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | Signal |  | Signal |  |  |
| Discrete-Time Signal | Discrete-to-Continuous-Time (D/C) Converter | Before Processing | Continuous-Time System | After Processing | Continuous-to-Discrete-Time (C/D) Converter | Discrete-Time Signal |
|  |  |  |  | Processing |  |  |

## Processing a Discrete-Time Signal With a Continuous-Time System

## Communication Systems



## General Structure of a Communication System

## Control Systems

Reference


General Structure of a Feedback Control System

- Engineers build systems that process/manipulate signals.
- We need a formal mathematical framework for the study of such systems.
- Such a framework is necessary in order to ensure that a system will meet the required specifications (e.g., performance and safety).
- If a system fails to meet the required specifications or fails to work altogether, negative consequences usually ensue.
- When a system fails to operate as expected, the consequences can sometimes be catastrophic.
- The (original) Tacoma Narrows Bridge was a suspension bridge linking Tacoma and Gig Harbor (WA, USA).
■ This mile-long bridge, with a 2,800-foot main span, was the third largest suspension bridge at the time of opening.
- Construction began in Nov. 1938 and took about 19 months to build at a cost of $\$ 6,400,000$.
- On July 1, 1940, the bridge opened to traffic.
- On Nov. 7, 1940 at approximately 11:00, the bridge collapsed during a moderate (42 miles/hour) wind storm.
- The bridge was supposed to withstand winds of up to 120 miles/hour.
- The collapse was due to wind-induced vibrations and an unstable mechanical system.
- Repair of the bridge was not possible.
- Fortunately, a dog trapped in an abandoned car was the only fatality.

Image of bridge collapse omitted for copyright reasons.
A video of the bridge collapse can be found at
https://youtu.be/j-zczJXSxnw.

## Part 2

## Preliminaries

## Section 2.1

## Functions, Sequences, System Operators, and Transforms

- A rational number is a number of the form $x / y$, where $x$ and $y$ are integers and $y \neq 0$ (i.e., a ratio of integers).
- For example, $-\frac{5}{3}, \frac{17}{11}$, and $0=\frac{0}{1}$ are rational numbers, whereas $\pi$ and $e$ are irrational numbers (i.e., not rational).
- The symbols employed to denote several commonly-used sets are as follows:

| Symbol | Set |
| :--- | :--- |
| $\mathbb{Z}$ | integers |
| $\mathbb{R}$ | real numbers |
| $\mathbb{C}$ | complex numbers |
| $\mathbb{Q}$ | rational numbers |

## Notation for Sets of Consecutive Integers

- For two integers $a$ and $b$, we define the following notation for sets of consecutive integers:

$$
\begin{aligned}
& {[a . . b]=\{x \in \mathbb{Z}: a \leq x \leq b\},} \\
& {[a . . b)=\{x \in \mathbb{Z}: a \leq x<b\},} \\
& (a . . b]=\{x \in \mathbb{Z}: a<x \leq b\}, \quad \text { and } \\
& (a . . b)=\{x \in \mathbb{Z}: a<x<b\} .
\end{aligned}
$$

- In this notation, $a$ and $b$ indicate the endpoints of the range for the set, and the type of brackets used (i.e., parenthesis versus square bracket) indicates whether each endpoint is included in the set.
- For example:
- [0..4] denotes the set of integers $\{0,1,2,3,4\}$;

ㅁ [0..4) denotes the set of integers $\{0,1,2,3\}$; and

- $[0 \ldots N-1]$ and $[0 \ldots N)$ both denote the set of integers $\{0,1,2, \ldots, N-1\}$.


## Notation for Intervals on the Real Line

- For two real numbers $a$ and $b$, we define the following notation for intervals on the real line:

$$
\begin{aligned}
& {[a, b]=\{x \in \mathbb{R}: a \leq x \leq b\},} \\
& (a, b)=\{x \in \mathbb{R}: a<x<b\}, \\
& {[a, b)=\{x \in \mathbb{R}: a \leq x<b\}, \quad \text { and }} \\
& (a, b]=\{x \in \mathbb{R}: a<x \leq b\} .
\end{aligned}
$$

- In this notation, $a$ and $b$ indicate the endpoints of the interval for the set, and the type of brackets used (i.e., parenthesis versus square bracket) indicate whether each endpoint is included in the set.
- For example:
- $[0,100]$ denotes the set of all real numbers from 0 to 100 , including both 0 and 100 ;
- $(-\pi, \pi]$ denotes the set of all real numbers from $-\pi$ to $\pi$, excluding $-\pi$ but including $\pi$; and
- $[-\pi, \pi)$ denotes the set of all real numbers from $-\pi$ to $\pi$, including $-\pi$ but excluding $\pi$.


## Mappings

- A mapping is a relationship involving two sets that associates each element in one set, called the domain, with an element from the other set, called the codomain.
- The notation $f: A \rightarrow B$ denotes a mapping $f$ whose domain is the set $A$ and whose codomain is the set $B$.
- Example:


$$
\begin{gathered}
f: A \rightarrow B \\
A=\{1,2,3,4\} \\
B=\{0,1,2,3\} \\
f(x)= \begin{cases}0 & x \in\{1,2\} \\
1 & x=4 \\
2 & x=3 .\end{cases}
\end{gathered}
$$

- Although many types of mappings exist, the types of most relevance to our study of signals and systems are: functions, sequences, system operators, and transforms.
- A function is a mapping where the domain is a set that is continuous in nature, such as the real numbers or complex numbers.
- In practice, the codomain is typically either the real numbers or complex numbers.
- Functions are also commonly referred to as continuous-time (CT) signalls.
- Example:
- Let $f: \mathbb{R} \rightarrow \mathbb{R}$ such that $f(t)=t^{2}$ (i.e., $f$ is the squaring function).
- The function $f$ maps each real number $t$ to the real number $f(t)=t^{2}$.
- The domain and codomain are the real numbers.
- Note that $f$ is a function, whereas $f(t)$ is a number (namely, the value of the function $f$ evaluated at $t$ ).
- Herein, we will focus almost exclusively on functions of a single independent variable (i.e., one-dimensional functions).


## Sequences

- A sequence is a mapping where the domain is a set that is discrete in nature, such as the integers, or a subset thereof.
- In practice, the codomain is typically either the real numbers or complex numbers.
- Sequences are also commonly referred to as discrete-time (DT) signals.
- Example:
- Let $f: \mathbb{Z}^{+} \rightarrow \mathbb{Z}^{+}$such that $f(n)=n^{2}$, where $\mathbb{Z}^{+}$denotes the set of (strictly) positive integers (i.e., $f$ is the sequence of perfect squares).
- The sequence $f$ maps each (strictly) positive integer $n$ to the (strictly) positive integer $f(n)=n^{2}$.
- The domain and codomain are $\mathbb{Z}^{+}$(i.e., the positive integers).
- Note that $f$ is a sequence, whereas $f(n)$ is a number (namely, the value of the sequence $f$ evaluated at $n$ ).
- As a matter of notation, the $n$th element of a sequence $x$ is denoted as either $x(n)$ or $x_{n}$.
- Herein, we will focus almost exclusively on sequences with a single independent variable (i.e., one-dimensional sequences).


## Remarks on Notation for Functions and Sequences

- For a real-valued function $f$ of a real variable and an arbitrary real number $t$, the expression $f$ denotes the function $f$ itself and the expression $f(t)$ denotes the value of the function $f$ evaluated at $t$.
- That is, $f$ is a function and $f(t)$ is a number.
- Unfortunately, the practice of using $f(t)$ to denote the function $f$ is quite common, although strictly speaking this is an abuse of notation.
- In contexts where imprecise notation may lead to problems, one should be careful to clearly distinguish between a function and its value.
- For the real-valued functions $f$ and $g$ of a real variable and an arbitrary real number $t$ :
- The expression $f+g$ denotes a function, namely, the function formed by adding the functions $f$ and $g$.
- The expression $f(t)+g(t)$ denotes a number, namely, the sum of: 1) the value of the function $f$ evaluated at $t$; and 2 ) the value of the function $g$ evaluated at $t$.
- Similar comments as the ones made above for functions also hold in the case of sequences.
- To express that two functions $f$ and $g$ are equal, we can write either:
$1 f=g$; or
$2 f(t)=g(t)$ for all $t$.
■ Of the preceding two expressions, the first (i.e., $f=g$ ) is usually preferable, as it is less verbose.
- For the functions $f$ and $g$ and an operation $\circ$ that is defined pointwise for functions (such as addition, subtraction, multiplication, and division), the following relationship holds:

$$
(f \circ g)(t)=f(t) \circ g(t)
$$

- Some operations o involving functions (such as convolution, to be discussed later) cannot be defined in a pointwise manner, in which case $(f \circ g)(t)$ is a valid mathematical expression, while $f(t) \circ g(t)$ is not.
- Again, similar comments as the ones made above for functions also hold in the case of sequences.


## System Operators

■ A system operator is a mapping used to represent a system.

- We will focus exclusively on the case of single-input single-output systems.
- A (single-input single-output) system operator maps a function or sequence representing the input of a system to a function or sequence representing the output of the system.
- The domain and codomain of a system operator are sets of functions or sequences, not sets of numbers.
■ Example:
- Let $\mathcal{H}: F \rightarrow F$ such that $\mathcal{H} x(t)=2 x(t)$ (for all $t \in \mathbb{R}$ ) and $F$ is the set of functions mapping $\mathbb{R}$ to $\mathbb{R}$.
- The system $\mathcal{H}$ maps a function to a function.
- In particular, the domain and codomain are each $F$, which is a set of functions.
- The system $\mathcal{H}$ multiplies its input function $x$ by a factor of 2 in order to produce its output function $\mathcal{H} x$.
- Note that $\mathcal{H} x$ is a function, not a number.
- For a system operator $\mathcal{H}$ and a function $x, \mathcal{H} x$ is the function produced as the output of the system $\mathcal{H}$ when the input is the function $x$.
- Brackets around the operand of an operator are often omitted when not required for grouping.
- For example, for an operator $\mathcal{H}$, a function $x$, and a real number $t$, we would normally prefer to write:
$1 \mathcal{H} x$ instead of the equivalent expression $\mathcal{H}(x)$; and
$2 \mathcal{H} x(t)$ instead of the equivalent expression $\mathcal{H}(x)(t)$.
- Also, note that $\mathcal{H} x$ is a function and $\mathcal{H} x(t)$ is a number (namely, the value of the function $\mathcal{H} x$ evaluated at $t)$.
- In the expression $\mathcal{H}\left(x_{1}+x_{2}\right)$, the brackets are needed for grouping, since $\mathcal{H}\left(x_{1}+x_{2}\right) \not \equiv \mathcal{H} x_{1}+x_{2}$ (where " $\not \equiv$ " means "not equivalent").
- When multiple operators are applied, they group from right to left.
- For example, for the operators $\mathcal{H}_{1}$ and $\mathcal{H}_{2}$, and the function $x$, the expression $\mathcal{H}_{2} \mathcal{H}_{1} x$ means $\mathcal{H}_{2}\left[\mathcal{H}_{1}(x)\right]$.
- For a system operator $\mathcal{H}$ and a sequence $x, \mathcal{H} x$ is the sequence produced as the output of the system $\mathcal{H}$ when the input is the sequence $x$.
■ Brackets around the operand of an operator are often omitted when not required for grouping.
- For example, for an operator $\mathcal{H}$, a sequence $x$, and an integer $n$, we would normally prefer to write:
$1 \mathcal{H} x$ instead of the equivalent expression $\mathcal{H}(x)$; and
$2 \mathcal{H} x(n)$ instead of the equivalent expression $\mathcal{H}(x)(n)$.
- Also, note that $\mathcal{H} x$ is a sequence and $\mathcal{H} x(n)$ is a number (namely, the value of the sequence $\mathcal{H} x$ evaluated at $n$ ).
- In the expression $\mathcal{H}\left(x_{1}+x_{2}\right)$, the brackets are needed for grouping, since $\mathcal{H}\left(x_{1}+x_{2}\right) \not \equiv \mathcal{H} x_{1}+x_{2}$ (where " $\not \equiv$ " means "not equivalent").
- When multiple operators are applied, they group from right to left.
- For example, for the operators $\mathcal{H}_{1}$ and $\mathcal{H}_{2}$, and the sequence $x$, the expression $\mathcal{H}_{2} \mathcal{H}_{1} x$ means $\mathcal{H}_{2}\left[\mathcal{H}_{1}(x)\right]$.
- Later, we will be introduced to several types of mappings known as transforms.
- Transforms have a mathematical structure similar to system operators.
- That is, transforms map functions/sequences to functions/sequences.
- Due to this similar structure, many of the earlier comments about system operators also apply to the case of transforms.
- For example, the Fourier transform (introduced later) is denoted as $\mathcal{F}$ and the result of applying the Fourier transform operator to the function/sequence $x$ is denoted as $\mathcal{F} x$.
- Some examples of transforms of interest in the study of signals and systems are listed on the next slide.


## Examples of Transforms

| Name | Domain | Codomain |
| :--- | :--- | :--- |
| CT Fourier Series | $T$-periodic functions <br> (with domain $\mathbb{R}$ ) | sequences <br> (with domain $\mathbb{Z}$ ) |
| CT Fourier Transform | functions <br> (with domain $\mathbb{R}$ ) | functions <br> (with domain $\mathbb{R}$ ) |
| Laplace Transform | functions <br> (with domain $\mathbb{R}$ ) | functions <br> (with domain $\mathbb{C}$ ) |
| DT Fourier Series | $N$-periodic sequences <br> (with domain $\mathbb{Z}$ ) | $N$-periodic sequences <br> (with domain $\mathbb{Z}$ ) |
| DT Fourier Transform | sequences <br> (with domain $\mathbb{Z}$ ) | $2 \pi$-periodic functions <br> (with domain $\mathbb{R}$ ) |
| Z Transform | sequences <br> (with domain $\mathbb{Z}$ ) | functions <br> (with domain $\mathbb{C}$ ) |

## Section 2.2

## Properties of Signals

## Even Symmetry

- A function $x$ is said to be even if it satisfies

$$
x(t)=x(-t) \text { for all } t \text { (where } t \text { is a real number). }
$$

- A sequence $x$ is said to be even if it satisfies

$$
x(n)=x(-n) \quad \text { for all } n \text { (where } n \text { is an integer). }
$$

- Geometrically, the graph of an even signal is symmetric with respect to the vertical axis.
- Some examples of even signals are shown below.




## Odd Symmetry

- A function $x$ is said to be odd if it satisfies

$$
x(t)=-x(-t) \text { for all } t \text { (where } t \text { is a real number). }
$$

- A sequence $x$ is said to be odld if it satisfies

$$
x(n)=-x(-n) \quad \text { for all } n \text { (where } n \text { is an integer). }
$$

- An odd signal $x$ must be such that $x(0)=0$.
- Geometrically, the graph of an odd signal is symmetric with respect to the origin.
- Some examples of odd signals are shown below.




## Conjugate Symmetry

- A function $x$ is said to be conjugate symmetric if it satisfies

$$
x(t)=x^{*}(-t) \text { for all } t \text { (where } t \text { is a real number). }
$$

■ A sequence $x$ is said to be conjugate symmetric if it satisfies

$$
x(n)=x^{*}(-n) \quad \text { for all } n \text { (where } n \text { is an integer). }
$$

- The real part of a conjugate symmetric function or sequence is even.
- The imaginary part of a conjugate symmetric function or sequence is odd.
- An example of a conjugate symmetric function is a complex sinusoid $x(t)=\cos (\omega t)+j \sin (\omega t)$, where $\omega$ is a real constant.
- A function $x$ is said to be periodic with period $T$ (or $T$-periodic) if, for some strictly-positive real constant $T$, the following condition holds:

$$
x(t)=x(t+T) \text { for all } t \text { (where } t \text { is a real number). }
$$

- A sequence $x$ is said to be periodic with period $N$ (or $N$-periodic) if, for some strictly-positive integer constant $N$, the following condition holds:

$$
x(n)=x(n+N) \quad \text { for all } n \text { (where } n \text { is an integer). }
$$

- Some examples of periodic signals are shown below.


- A function/sequence that is not periodic is said to be aperiodic.
- A $T$-periodic function $x$ is said to have frequency $\frac{1}{T}$ and angular frequency $\frac{2 \pi}{T}$.
- An $N$-periodic sequence $x$ is said to have frequency $\frac{1}{N}$ and angular frequency $\frac{2 \pi}{N}$.
- The period of a periodic signal is not unique. That is, a signal that is periodic with period $T$ is also periodic with period $k T$, for every (strictly) positive integer $k$.

- The smallest period with which a signal is periodic is called the fundamental period and its corresponding frequency is called the fundamental frequency.


## Part 3

## Continuous-Time (CT) Signals and Systems

## Section 3.1

## Independent- and Dependent-Variable Transformations

- Time shifting (also called translation) maps the input function $x$ to the output function $y$ as given by

$$
y(t)=x(t-b)
$$

where $b$ is a real number.

- Such a transformation shifts the function (to the left or right) along the time axis.
- If $b>0, y$ is shifted to the right by $|b|$, relative to $x$ (i.e., delayed in time).
- If $b<0, y$ is shifted to the left by $|b|$, relative to $x$ (i.e., advanced in time).

Time Shifting (Translation): Example




- Time reversal (also known as reflection) maps the input function $x$ to the output function $y$ as given by

$$
y(t)=x(-t) .
$$

- Geometrically, the output function $y$ is a reflection of the input function $x$ about the (vertical) line $t=0$.


- Time compression/expansion (also called dilation) maps the input function $x$ to the output function $y$ as given by

$$
y(t)=x(a t),
$$

where $a$ is a strictly positive real number.

- Such a transformation is associated with a compression/expansion along the time axis.
- If $a>1, y$ is compressed along the horizontal axis by a factor of $a$, relative to $x$.
- If $a<1, y$ is expanded (i.e., stretched) along the horizontal axis by a factor of $\frac{1}{a}$, relative to $x$.



- Time scaling maps the input function $x$ to the output function $y$ as given by

$$
y(t)=x(a t),
$$

where $a$ is a nonzero real number.

- Such a transformation is associated with a dilation (i.e., compression/expansion along the time axis) and/or time reversal.
- If $|a|>1$, the function is compressed along the time axis by a factor of $|a|$.
- If $|a|<1$, the function is expanded (i.e., stretched) along the time axis by a factor of $\left|\frac{1}{a}\right|$.
- If $|a|=1$, the function is neither expanded nor compressed.
- If $a<0$, the function is also time reversed.
- Dilation (i.e., expansion/compression) and time reversal commute.
- Time reversal is a special case of time scaling with $a=-1$; and time compression/expansion is a special case of time scaling with $a>0$.


## Time Scaling (Dilation/Reflection): Example






## Combined Time Scaling and Time Shifting

- Consider a transformation that maps the input function $x$ to the output function $y$ as given by

$$
y(t)=x(a t-b),
$$

where $a$ and $b$ are real numbers and $a \neq 0$.
■ The above transformation can be shown to be the combination of a time-scaling operation and time-shifting operation.

- Since time scaling and time shifting do not commute, we must be particularly careful about the order in which these transformations are applied.
- The above transformation has two distinct but equivalent interpretations:

11 first, time shifting $x$ by $b$, and then time scaling the result by $a$;
$\square$ first, time scaling $x$ by $a$, and then time shifting the result by $b / a$.

- Note that the time shift is not by the same amount in both cases.
- In particular, note that when time scaling is applied first followed by time shifting, the time shift is by $b / a$, not $b$.


## Combined Time Scaling and Time Shifting: Example

time shift by 1 and then time scale by 2

Given $x$ as shown below, find

$$
y(t)=x(2 t-1) .
$$




time scale by 2 and then time shift by $\frac{1}{2}$


- A transformation of the independent variable can be viewed in terms of

11 the effect that the transformation has on the function; or
[ the effect that the transformation has on the horizontal axis.

- This distinction is important because such a transformation has opposite effects on the function and horizontal axis.
■ For example, the (time-shifting) transformation that replaces $t$ by $t-b$ (where $b$ is a real number) in $x(t)$ can be viewed as a transformation that

11 shifts the function $x$ right by $b$ units; or
[2 shifts the horizontal axis left by $b$ units.

- In our treatment of independent-variable transformations, we are only interested in the effect that a transformation has on the function.
- If one is not careful to consider that we are interested in the function perspective (as opposed to the axis perspective), many aspects of independent-variable transformations will not make sense.


## Amplitude Scaling

- Amplitude scaling maps the input function $x$ to the output function $y$ as given by

$$
y(t)=a x(t),
$$

where $a$ is a real number.

- Geometrically, the output function $y$ is expanded/compressed in amplitude and/or reflected about the horizontal axis.






## Amplitude Shifting

- Amplitude shifting maps the input function $x$ to the output function $y$ as given by

$$
y(t)=x(t)+b,
$$

where $b$ is a real number.

- Geometrically, amplitude shifting adds a vertical displacement to $x$.




## Combined Amplitude Scaling and Amplitude Shifting

- We can also combine amplitude scaling and amplitude shifting transformations.
- Consider a transformation that maps the input function $x$ to the output function $y$, as given by

$$
y(t)=a x(t)+b
$$

where $a$ and $b$ are real numbers.

- Equivalently, the above transformation can be expressed as

$$
y(t)=a\left[x(t)+\frac{b}{a}\right] .
$$

- The above transformation is equivalent to:

1 first amplitude scaling $x$ by $a$, and then amplitude shifting the resulting function by $b$; or
2 first amplitude shifting $x$ by $b / a$, and then amplitude scaling the resulting function by $a$.

## Section 3.2

## Properties of Functions

- Sums involving even and odd functions have the following properties:
- The sum of two even functions is even.
- The sum of two odd functions is odd.
- The sum of an even function and odd function is neither even nor odd, provided that neither of the functions is identically zero.
- That is, the sum of functions with the same type of symmetry also has the same type of symmetry.
- Products involving even and odd functions have the following properties:
- The product of two even functions is even.
- The product of two odd functions is even.
- The product of an even function and an odd function is odd.
- That is, the product of functions with the same type of symmetry is even, while the product of functions with opposite types of symmetry is odd.


## Decomposition of a Function into Even and Odd Parts

- Every function $x$ has a unique representation of the form

$$
x(t)=x_{\mathrm{e}}(t)+x_{\mathrm{o}}(t),
$$

where the functions $x_{\mathrm{e}}$ and $x_{\mathrm{o}}$ are even and odd, respectively.

- In particular, the functions $x_{\mathrm{e}}$ and $x_{\mathrm{o}}$ are given by

$$
x_{\mathrm{e}}(t)=\frac{1}{2}[x(t)+x(-t)] \quad \text { and } \quad x_{\mathrm{o}}(t)=\frac{1}{2}[x(t)-x(-t)] .
$$

- The functions $x_{\mathrm{e}}$ and $x_{\mathrm{o}}$ are called the even part and odld part of $x$, respectively.
- For convenience, the even and odd parts of $x$ are often denoted as Even $\{x\}$ and $\operatorname{Odd}\{x\}$, respectively.
- Sum of periodic functions. For two periodic functions $x_{1}$ and $x_{2}$ with fundamental periods $T_{1}$ and $T_{2}$, respectively, and the sum $y=x_{1}+x_{2}$ :

1 The sum $y$ is periodic if and only if the ratio $T_{1} / T_{2}$ is a rational number (i.e., the quotient of two integers).
0 If $y$ is periodic, its fundamental period is $r T_{1}$ (or equivalently, $q T_{2}$, since $r T_{1}=q T_{2}$ ), where $T_{1} / T_{2}=q / r$ and $q$ and $r$ are integers and coprime (i.e., have no common factors). (Note that $r T_{1}$ is simply the least common multiple of $T_{1}$ and $T_{2}$.)

- Although the above theorem only directly addresses the case of the sum of two functions, the case of $N$ functions (where $N>2$ ) can be handled by applying the theorem repeatedly $N-1$ times.


## Right-Sided Functions

- A function $x$ is said to be right sided if, for some (finite) real constant $t_{0}$, the following condition holds:

$$
x(t)=0 \quad \text { for all } t<t_{0}
$$

(i.e., $x$ is only potentially nonzero to the right of $t_{0}$ ).

- An example of a right-sided function is shown below.

- A function $x$ is said to be causal if

$$
x(t)=0 \quad \text { for all } t<0
$$

- A causal function is a special case of a right-sided function.
- A causal function is not to be confused with a causal system. In these two contexts, the word "causal" has very different meanings.


## Left-Sided Functions

- A function $x$ is said to be left sided if, for some (finite) real constant $t_{0}$, the following condition holds:

$$
x(t)=0 \quad \text { for all } t>t_{0}
$$

(i.e., $x$ is only potentially nonzero to the left of $t_{0}$ ).

- An example of a left-sided function is shown below.

- Similarly, a function $x$ is said to be anticausal if

$$
x(t)=0 \quad \text { for all } t>0 .
$$

- An anticausal function is a special case of a left-sided function.
- An anticausal function is not to be confused with an anticausal system. In these two contexts, the word "anticausal" has very different meanings.
- A function that is both left sided and right sided is said to be finite duration (or time limited).
- An example of a finite duration function is shown below.

- A function that is neither left sided nor right sided is said to be two sided.
- An example of a two-sided function is shown below.

- A function $x$ is said to be bounded if there exists some (finite) positive real constant $A$ such that

$$
|x(t)| \leq A \quad \text { for all } t
$$

(i.e., $x(t)$ is finite for all $t$ ).

- For example, the sine and cosine functions are bounded, since

$$
|\sin t| \leq 1 \text { for all } t \quad \text { and } \quad|\cos t| \leq 1 \text { for all } t .
$$

- In contrast, the tangent function and any nonconstant polynomial function $p$ (e.g., $p(t)=t^{2}$ ) are unbounded, since

$$
\lim _{t \rightarrow \pi / 2}|\tan t|=\infty \quad \text { and } \quad \lim _{|t| \rightarrow \infty}|p(t)|=\infty .
$$

- The energy $E$ contained in the function $x$ is given by

$$
E=\int_{-\infty}^{\infty}|x(t)|^{2} d t
$$

- A signal with finite energy is said to be an energy signal.
- The average power $P$ contained in the function $x$ is given by

$$
P=\lim _{T \rightarrow \infty} \frac{1}{T} \int_{-T / 2}^{T / 2}|x(t)|^{2} d t
$$

- A signal with (nonzero) finite average power is said to be a power signal.


## Section 3.3

## Elementary Functions

- A real sinusoidal function is a function of the form

$$
x(t)=A \cos (\omega t+\theta),
$$

where $A, \omega$, and $\theta$ are real constants.

- Such a function is periodic with fundamental period $T=\frac{2 \pi}{|\omega|}$ and fundamental frequency $|\omega|$.
- A real sinusoid has a plot resembling that shown below.



## Complex Exponential Functions

- A complex exponential function is a function of the form

$$
x(t)=A e^{\lambda t}
$$

where $A$ and $\lambda$ are complex constants.

- A complex exponential can exhibit one of a number of distinct modes of behavior, depending on the values of its parameters $A$ and $\lambda$.
- For example, as special cases, complex exponentials include real exponentials and complex sinusoids.
- A real exponential function is a special case of a complex exponential $x(t)=A e^{\lambda t}$, where $A$ and $\lambda$ are restricted to be real numbers.
- A real exponential can exhibit one of three distinct modes of behavior, depending on the value of $\lambda$, as illustrated below.
■ If $\lambda>0, x(t)$ increases exponentially as $t$ increases (i.e., a growing exponential).
- If $\lambda<0, x(t)$ decreases exponentially as $t$ increases (i.e., a decaying exponential).
- If $\lambda=0, x(t)$ simply equals the constant $A$.

$\lambda>0$

$\lambda=0$

$\lambda<0$


## Complex Sinusoidal Functions

- A complex sinusoidal function is a special case of a complex exponential $x(t)=A e^{\lambda t}$, where $A$ is complex and $\lambda$ is purely imaginary (i.e., $\operatorname{Re}\{\lambda\}=0)$.
- That is, a complex sinusoidal function is a function of the form

$$
x(t)=A e^{j \omega t}
$$

where $A$ is complex and $\omega$ is real.

- By expressing $A$ in polar form as $A=|A| e^{j \theta}$ (where $\theta$ is real) and using Euler's relation, we can rewrite $x(t)$ as

$$
x(t)=\underbrace{|A| \cos (\omega t+\theta)}_{\operatorname{Re}\{x(t)\}}+j \underbrace{|A| \sin (\omega t+\theta)}_{\operatorname{Im}\{x(t)\}} .
$$

- Thus, $\operatorname{Re}\{x\}$ and $\operatorname{Im}\{x\}$ are the same except for a time shift.
- Also, $x$ is periodic with fundamental period $T=\frac{2 \pi}{|\omega|}$ and fundamental frequency $|\omega|$.


## Complex Sinusoidal Functions (Continued)

- The graphs of $\operatorname{Re}\{x\}$ and $\operatorname{Im}\{x\}$ have the forms shown below.



Plot of $x(t)=e^{j \omega t}$ for $\omega \in\{2 \pi,-2 \pi\}$


$$
\omega=2 \pi
$$


$\omega=-2 \pi$

## General Complex Exponential Functions

■ In the most general case of a complex exponential function $x(t)=A e^{\lambda t}, A$ and $\lambda$ are both complex.
■ Letting $A=|A| e^{j \theta}$ and $\lambda=\sigma+j \omega$ (where $\theta, \sigma$, and $\omega$ are real), and using Euler's relation, we can rewrite $x(t)$ as

$$
x(t)=\underbrace{|A| e^{\sigma t} \cos (\omega t+\theta)}_{\operatorname{Re}\{x(t)\}}+j \underbrace{|A| e^{\sigma t} \sin (\omega t+\theta)}_{\operatorname{Im}\{x(t)\}} .
$$

- Thus, $\operatorname{Re}\{x\}$ and $\operatorname{Im}\{x\}$ are each the product of a real exponential and real sinusoid.
- One of three distinct modes of behavior is exhibited by $x(t)$, depending on the value of $\sigma$.
- If $\sigma=0, \operatorname{Re}\{x\}$ and $\operatorname{Im}\{x\}$ are real sinusoids.
- If $\sigma>0, \operatorname{Re}\{x\}$ and $\operatorname{Im}\{x\}$ are each the product of a real sinusoid and a growing real exponential.
- If $\sigma<0, \operatorname{Re}\{x\}$ and $\operatorname{Im}\{x\}$ are each the product of a real sinusoid and a decaying real exponential.


## General Complex Exponential Functions (Continued)

■ The three modes of behavior for $\operatorname{Re}\{x\}$ and $\operatorname{Im}\{x\}$ are illustrated below.

$\sigma>0$

$\sigma=0$

$\sigma<0$

## Relationship Between Complex Exponentials and Real Sinusoids

- From Euler's relation, a complex sinusoid can be expressed as the sum of two real sinusoids as

$$
A e^{j \omega t}=A \cos (\omega t)+j A \sin (\omega t)
$$

- Moreover, a real sinusoid can be expressed as the sum of two complex sinusoids using the identities

$$
\begin{gathered}
A \cos (\omega t+\theta)=\frac{A}{2}\left[e^{j(\omega t+\theta)}+e^{-j(\omega t+\theta)}\right] \quad \text { and } \\
A \sin (\omega t+\theta)=\frac{A}{2 j}\left[e^{j(\omega t+\theta)}-e^{-j(\omega t+\theta)}\right] .
\end{gathered}
$$

- Note that, above, we are simply restating results from the (appendix) material on complex analysis.
- The unit-step function (also known as the Heaviside function), denoted $u$, is defined as

$$
u(t)= \begin{cases}1 & t \geq 0 \\ 0 & \text { otherwise }\end{cases}
$$

- Due to the manner in which $u$ is used in practice, the actual value of $u(0)$ is unimportant. Sometimes values of 0 and $\frac{1}{2}$ are also used for $u(0)$.
- A plot of this function is shown below.

- The signum function, denoted sgn, is defined as

$$
\operatorname{sgn} t= \begin{cases}1 & t>0 \\ 0 & t=0 \\ -1 & t<0\end{cases}
$$

- From its definition, one can see that the signum function simply computes the sign of a number.
- A plot of this function is shown below.

- The rectangular function (also called the unit-rectangular pulse function), denoted rect, is given by

$$
\operatorname{rect} t= \begin{cases}1 & -\frac{1}{2} \leq t<\frac{1}{2} \\ 0 & \text { otherwise }\end{cases}
$$

- Due to the manner in which the rect function is used in practice, the actual value of rect $t$ at $t= \pm \frac{1}{2}$ is unimportant. Sometimes different values are used from those specified above.
- A plot of this function is shown below.



## Indicator Function

- Functions and sequences that are one over some subset of their domain and zero elsewhere appear very frequently in engineering (e.g., the unit-step function and rectangular function).
- Indicator function notation provides a concise way to denote such functions and sequences.
- The indicator function of a subset $S$ of a set $A$, denoted $\chi_{S}$, is defined as

$$
\chi_{S}(t)= \begin{cases}1 & \text { if } t \in S \\ 0 & \text { otherwise }\end{cases}
$$

- A rectangular pulse (defined on $\mathbb{R}$ ) having an amplitude of 1 , a leading edge at $a$, and falling edge at $b$ is $\chi_{[a, b]}$.
- The unit-step function (defined on $\mathbb{R}$ ) is $\chi_{[0, \infty)}$.
- The unit-rectangular pulse (defined on $\mathbb{R}$ ) is $\chi_{[-1 / 2,1 / 2]}$.
- The unit-step sequence (defined on $\mathbb{Z}$ ) is $\chi_{[0 . . \infty)}$.


## Triangular Function

- The triangular function (also called the unit-triangular pulse function), denoted tri, is defined as

$$
\operatorname{tri} t= \begin{cases}1-2|t| & |t| \leq \frac{1}{2} \\ 0 & \text { otherwise }\end{cases}
$$

- A plot of this function is shown below.



## Cardinal Sine Function

- The cardinal sine function, denoted sinc, is given by

$$
\operatorname{sinc} t=\frac{\sin t}{t}
$$

- By l'Hopital's rule, $\operatorname{sinc} 0=1$.
- A plot of this function for part of the real line is shown below. [Note that the oscillations in $\operatorname{sinc} t$ do not die out for finite $t$.]

- The floor function, denoted $\lfloor\cdot\rfloor$, is a function that maps a real number $x$ to the largest integer not more than $x$.
- In other words, the floor function rounds a real number to the nearest integer in the direction of negative infinity.
- For example,

$$
\left\lfloor-\frac{1}{2}\right\rfloor=-1, \quad\left\lfloor\frac{1}{2}\right\rfloor=0, \quad \text { and } \quad\lfloor 1\rfloor=1 .
$$

- The ceiling function, denoted $\lceil\cdot\rceil$, is a function that maps a real number $x$ to the smallest integer not less than $x$.
- In other words, the ceiling function rounds a real number to the nearest integer in the direction of positive infinity.
- For example,

$$
\left\lceil-\frac{1}{2}\right\rceil=0, \quad\left\lceil\frac{1}{2}\right\rceil=1, \quad \text { and } \quad\lceil 1\rceil=1
$$

- Several useful properties of the floor and ceiling functions include:

$$
\begin{gathered}
\lfloor x+n\rfloor=\lfloor x\rfloor+n \quad \text { for } x \in \mathbb{R} \text { and } n \in \mathbb{Z} ; \\
\lceil x+n\rceil=\lceil x\rceil+n \text { for } x \in \mathbb{R} \text { and } n \in \mathbb{Z} ; \\
\lceil x\rceil=-\lfloor-x\rfloor \quad \text { for } x \in \mathbb{R} ; \\
\lfloor x\rfloor=-\lceil-x\rceil \quad \text { for } x \in \mathbb{R} ; \\
\left\lceil\frac{m}{n}\right\rceil=\left\lfloor\frac{m+n-1}{n}\right\rfloor=\left\lfloor\frac{m-1}{n}\right\rfloor+1 \quad \text { for } m, n \in \mathbb{Z} \text { and } n>0 ; \quad \text { and } \\
\left\lfloor\frac{m}{n}\right\rfloor=\left\lceil\frac{m-n+1}{n}\right\rceil=\left\lceil\frac{m+1}{n}\right\rceil-1 \quad \text { for } m, n \in \mathbb{Z} \text { and } n>0 .
\end{gathered}
$$

## Delta Function

- The delta function (also known as the Dirac delta function or unit-impulse function), denoted $\delta$, is defined as the function with the following two properties:

$$
\begin{gathered}
\delta(t)=0 \quad \text { for } t \neq 0 \quad \text { and } \\
\int_{-\infty}^{\infty} \delta(t) d t=1 .
\end{gathered}
$$

- Technically, $\delta$ is not a function in the ordinary sense. Rather, it is what is known as a generalized function. Consequently, the $\delta$ function sometimes behaves in unusual ways.
- Graphically, the delta function is represented as shown below.




## Delta Function as a Limit

- Consider the function $g_{\varepsilon}$ (where $\varepsilon$ is a real constant) defined by

$$
g_{\varepsilon}(t)= \begin{cases}\frac{1}{|\varepsilon|} & |t|<\frac{|\varepsilon|}{2} \\ 0 & \text { otherwise }\end{cases}
$$

- A plot of $g_{\varepsilon}$ is shown below.

- Clearly, for any choice of $\varepsilon, \int_{-\infty}^{\infty} g_{\varepsilon}(t) d t=1$.
- The function $\delta$ can be obtained as the following limit:

$$
\delta(t)=\lim _{\varepsilon \rightarrow 0} g_{\varepsilon}(t)
$$

- That is, $\delta$ can be viewed as a limiting case of a rectangular pulse where the pulse width becomes infinitesimally small and the pulse height becomes infinitely large in such a way that the integral of the resulting function remains unity.
- Equivalence property. For any continuous function $x$ and any real constant $t_{0}$,

$$
x(t) \delta\left(t-t_{0}\right)=x\left(t_{0}\right) \delta\left(t-t_{0}\right)
$$

- Sifting property. For any continuous function $x$ and any real constant $t_{0}$,

$$
\int_{-\infty}^{\infty} x(t) \delta\left(t-t_{0}\right) d t=x\left(t_{0}\right)
$$

- Scaling property. For any nonzero real constant $a$,

$$
\delta(a t)=\frac{1}{|a|} \delta(t)
$$

- Even property. The $\delta$ function is such that

$$
\delta(t)=\delta(-t)
$$

## Graphical Interpretation of Equivalence Property



Function $x$


Time-Shifted Delta Function


Product

- For real constants $a$ and $b$ where $a \leq b$, consider a function $x$ of the form

$$
x(t)= \begin{cases}1 & a \leq t<b \\ 0 & \text { otherwise }\end{cases}
$$

(i.e., $x$ is a rectangular pulse of height one, with a rising edge at $a$ and falling edge at $b$ ).

- The function $x$ can be equivalently written as

$$
x(t)=u(t-a)-u(t-b)
$$

(i.e., the difference of two time-shifted unit-step functions).

- Unlike the original expression for $x$, this latter expression for $x$ does not involve multiple cases.
- In effect, by using unit-step functions, we have collapsed a formula involving multiple cases into a single expression.
- The idea from the previous slide can be extended to handle any function that is defined in a piecewise manner (i.e., via an expression involving multiple cases).
- That is, by using unit-step functions, we can always collapse a formula involving multiple cases into a single expression.
- Often, simplifying a formula in this way can be quite beneficial.


## Section 3.4

## Continuous-Time (CT) Systems

## CT Systems

- A system with input $x$ and output $y$ can be described by the equation

$$
y=\mathcal{H} x,
$$

where $\mathcal{H}$ denotes an operator (i.e., transformation).

- Note that the operator $\mathcal{H}$ maps a function to a function (not a number to a number).
- Alternatively, we can express the above relationship using the notation

$$
x \xrightarrow{\mathcal{H}} y .
$$

- If clear from the context, the operator $\mathcal{H}$ is often omitted, yielding the abbreviated notation

$$
x \rightarrow y .
$$

- Note that the symbols " $\rightarrow$ " and " $=$ " have very different meanings.

■ The symbol " $\rightarrow$ " should be read as "produces" (not as "equals").

## Block Diagram Representations

- Often, a system defined by the operator $\mathcal{H}$ and having the input $x$ and output $y$ is represented in the form of a block diagram as shown below.



## Interconnection of Systems

- Two basic ways in which systems can be interconnected are shown below.

- A series (or cascade) connection ties the output of one system to the input of the other.
- The overall series-connected system is described by the equation

$$
y=\mathcal{H}_{2} \mathcal{H}_{1} x
$$

- A parallel connection ties the inputs of both systems together and sums their outputs.
- The overall parallel-connected system is described by the equation

$$
y=\mathcal{H}_{1} x+\mathcal{H}_{2} x .
$$

## Section 3.5

## Properties of (CT) Systems

- A system $\mathcal{H}$ is said to be memoryless if, for every real constant $t_{0}, \mathcal{H} x\left(t_{0}\right)$ does not depend on $x(t)$ for some $t \neq t_{0}$.
- In other words, a memoryless system is such that the value of its output at any given point in time can depend on the value of its input at only the same point in time.
- A system that is not memoryless is said to have memory.
- Although simple, a memoryless system is not very flexible, since its current output value cannot rely on past or future values of the input.


## Memory (Continued)

If the system $\mathcal{H}$ is memoryless, the output $\mathcal{H} x$ at $t_{0}$ can depend on the input $x$ only at $t_{0}$.


Consider the calculation of the output $\mathcal{H} x$ at $t_{0}$.

## Causality

- A system $\mathcal{H}$ is said to be causal if, for every real constant $t_{0}, \mathcal{H} x\left(t_{0}\right)$ does not depend on $x(t)$ for some $t>t_{0}$.
- In other words, a causal system is such that the value of its output at any given point in time can depend on the value of its input at only the same or earlier points in time (i.e., not later points in time).
- If the independent variable $t$ represents time, a system must be causal in order to be physically realizable.
- Noncausal systems can sometimes be useful in practice, however, since the independent variable need not always represent time (e.g., the independent variable might represent position).
- A memoryless system is always causal, although the converse is not necessarily true.


## Causality (Continued)

If the system $\mathcal{H}$ is causal, the output $\mathcal{H} x$ at $t_{0}$ can depend on the input $x$ only at points $t \leq t_{0}$.


Consider the calculation of the output $\mathcal{H} x$ at $t_{0}$.

## Invertibility

- The inverse of a system $\mathcal{H}$ (if it exists) is another system $\mathcal{H}^{-1}$ such that, for every function $x$,

$$
\mathcal{H}^{-1} \mathcal{H} x=x
$$

(i.e., the system formed by the cascade interconnection of $\mathcal{H}$ followed by $\mathcal{H}^{-1}$ is a system whose input and output are equal).

- A system is said to be invertible if it has a corresponding inverse system (i.e., its inverse exists).
- Equivalently, a system is invertible if its input can always be uniquely determined from its output.
- An invertible system will always produce distinct outputs from any two distinct inputs (i.e., $x_{1} \neq x_{2} \Rightarrow \mathcal{H}\left(x_{1} \neq \mathcal{H}\left(x_{2}\right)\right.$.
- To show that a system is invertible, we simply find the inverse system.
- To show that a system is not invertible, we find two distinct inputs that result in identical outputs (i.e., $x_{1} \neq x_{2}$ and $\mathcal{H}\left(x_{1}=\mathcal{H}\left(x_{2}\right)\right.$.
- In practical terms, invertible systems are "nice" in the sense that their effects can be undone.


## Invertibility (Continued)

■ A system $\mathcal{H}^{-1}$ being the inverse of $\mathcal{H}$ means that the following two systems are equivalent (i.e., $\mathcal{H}^{-1} \mathcal{H}$ is an identity):


System 1: $y=\mathcal{H}^{-1} \mathcal{H} x$


System 2: $y=x$

## Bounded-Input Bounded-Output (BIBO) Stability

- A system $\mathcal{H}$ is said to be bounded-input bounded-output (BIBO) stable if, for every bounded function $x, \mathcal{H} x$ is bounded (i.e., $|x(t)|<\infty$ for all $t$ implies that $|\mathcal{H} x(t)|<\infty$ for all $t$ ).
- In other words, a BIBO stable system is such that it guarantees to always produce a bounded output as long as its input is bounded.
- To show that a system is BIBO stable, we must show that every bounded input leads to a bounded output.
- To show that a system is not BIBO stable, we only need to find a single bounded input that leads to an unbounded output.
- In practical terms, a BIBO stable system is well behaved in the sense that, as long as the system input is finite everywhere (in its domain), the output will also be finite everywhere.
■ Usually, a system that is not BIBO stable will have serious safety issues.
- For example, a portable music player with a battery input of 3.7 volts and headset output of $\infty$ volts would result in one vaporized human (and likely a big lawsuit as well).

■ A system $\mathcal{H}$ is said to be time invariant (TII) (or shift invariant (SI)) if, for every function $x$ and every real constant $t_{0}$, the following condition holds:

$$
\mathcal{H} x\left(t-t_{0}\right)=\mathcal{H} x^{\prime}(t) \text { for all } t, \quad \text { where } \quad x^{\prime}(t)=x\left(t-t_{0}\right)
$$

(i.e., $\mathcal{H}$ commutes with time shifts).

- In other words, a system is time invariant if a time shift (i.e., advance or delay) in the input always results only in an identical time shift in the output.
- A system that is not time invariant is said to be time varying.
- In simple terms, a time invariant system is a system whose behavior does not change with respect to time.
■ Practically speaking, compared to time-varying systems, time-invariant systems are much easier to design and analyze, since their behavior does not change with respect to time.
- Let $\mathcal{S}_{t_{0}}$ denote an operator that applies a time shift of $t_{0}$ to a function (i.e., $\left.\mathcal{S}_{t_{0}} x(t)=x\left(t-t_{0}\right)\right)$.
■ A system $\mathcal{H}$ is time invariant if and only if the following two systems are equivalent (i.e., $\mathcal{H}$ commutes with $\mathcal{S}_{t_{0}}$ ):


System 1: $y=\mathcal{H S}_{t_{0}} x$
$\left[\begin{array}{c}y(t)=\mathcal{H} x^{\prime}(t) \\ x^{\prime}(t)=\mathcal{S}_{t_{0}} x(t)=x\left(t-t_{0}\right)\end{array}\right]$


System 2: $y=\mathcal{S}_{t_{0}} \mathcal{H} x$ $\left[y(t)=\mathcal{H} x\left(t-t_{0}\right)\right]$

## Additivity, Homogeneity, and Linearity

■ A system $\mathcal{H}$ is said to be additive if, for all functions $x_{1}$ and $x_{2}$, the following condition holds:

$$
\mathcal{H}\left(x_{1}+x_{2}\right)=\mathcal{H} x_{1}+\mathcal{H} x_{2}
$$

(i.e., $\mathcal{H}$ commutes with addition).

- A system $\mathcal{H}$ is said to be homogeneous if, for every function $x$ and every complex constant $a$, the following condition holds:

$$
\mathcal{H}(a x)=a \mathcal{H} x
$$

(i.e., $\mathcal{H}$ commutes with scalar multiplication).

- A system that is both additive and homogeneous is said to be linear.
- In other words, a system $\mathcal{H}$ is linear, if for all functions $x_{1}$ and $x_{2}$ and all complex constants $a_{1}$ and $a_{2}$, the following condition holds:

$$
\mathcal{H}\left(a_{1} x_{1}+a_{2} x_{2}\right)=a_{1} \mathcal{H} x_{1}+a_{2} \mathcal{H} x_{2}
$$

(i.e., $\mathcal{H}$ commutes with linear combinations).

- The linearity property is also referred to as the superposition property.
- Practically speaking, linear systems are much easier to design and analyze than nonlinear systems.


## Additivity, Homogeneity, and Linearity (Continued 1)

- The system $\mathcal{H}$ is additive if and only if the following two systems are equivalent (i.e., $\mathcal{H}$ commutes with addition):



System 2: $y=\mathcal{H} x_{1}+\mathcal{H} x_{2}$

- The system $\mathcal{H}$ is homogeneous if and only if the following two systems are equivalent (i.e., $\mathcal{H}$ commutes with scalar multiplication):


System 1: $y=\mathcal{H}(a x)$


System 2: $y=a \mathcal{H}(x$

## Additivity, Homogeneity, and Linearity (Continued 2)

- The system $\mathcal{H}$ is linear if and only if the following two systems are equivalent (i.e., $\mathcal{H}$ commutes with linear combinations):


System 1: $y=\mathcal{H}\left(a_{1} x_{1}+a_{2} x_{2}\right)$


System 2: $y=a_{1} \mathcal{H} x_{1}+a_{2} \mathcal{H} x_{2}$

- A function $x$ is said to be an eigenfunction of the system $\mathcal{H}$ with the eigenvalue $\lambda$ if

$$
\mathcal{H} x=\lambda x
$$

where $\lambda$ is a complex constant.

- In other words, the system $\mathcal{H}$ acts as an ideal amplifier for each of its eigenfunctions $x$, where the amplifier gain is given by the corresponding eigenvalue $\lambda$.
- Different systems have different eigenfunctions.
- Many of the mathematical tools developed for the study of CT systems have eigenfunctions as their basis.


## Part 4

## Continuous-Time Linear Time-Invariant (LTI) Systems

- In engineering, linear time-invariant (LTI) systems play a very important role.
- Very powerful mathematical tools have been developed for analyzing LTI systems.
- LTI systems are much easier to analyze than systems that are not LTI.
- In practice, systems that are not LTI can be well approximated using LTI models.
■ So, even when dealing with systems that are not LTI, LTI systems still play an important role.


## Section 4.1

## Convolution

- The (CT) convolution of the functions $x$ and $h$, denoted $x * h$, is defined as the function

$$
x * h(t)=\int_{-\infty}^{\infty} x(\tau) h(t-\tau) d \tau .
$$

- The convolution result $x * h$ evaluated at the point $t$ is simply a weighted average of the function $x$, where the weighting is given by $h$ time reversed and shifted by $t$.
■ Herein, the asterisk symbol (i.e., "*") will always be used to denote convolution, not multiplication.
- As we shall see, convolution is used extensively in systems theory.
- In particular, convolution has a special significance in the context of LTI systems.
- To compute the convolution

$$
x * h(t)=\int_{-\infty}^{\infty} x(\tau) h(t-\tau) d \tau
$$

we proceed as follows:
1 Plot $x(\tau)$ and $h(t-\tau)$ as a function of $\tau$.
[2 Initially, consider an arbitrarily large negative value for $t$. This will result in $h(t-\tau)$ being shifted very far to the left on the time axis.
3 Write the mathematical expression for $x * h(t)$.
4 Increase $t$ gradually until the expression for $x * h(t)$ changes form. Record the interval over which the expression for $x * h(t)$ was valid.
5 Repeat steps 3 and 4 until $t$ is an arbitrarily large positive value. This corresponds to $h(t-\tau)$ being shifted very far to the right on the time axis.
6 The results for the various intervals can be combined in order to obtain an expression for $x * h(t)$ for all $t$.

- The convolution operation is commutative. That is, for any two functions $x$ and $h$,

$$
x * h=h * x
$$

- The convolution operation is associative. That is, for any functions $x, h_{1}$, and $h_{2}$,

$$
\left(x * h_{1}\right) * h_{2}=x *\left(h_{1} * h_{2}\right) .
$$

- The convolution operation is distributive with respect to addition. That is, for any functions $x, h_{1}$, and $h_{2}$,

$$
x *\left(h_{1}+h_{2}\right)=x * h_{1}+x * h_{2}
$$

- For any function $x$,

$$
x * \delta(t)=\int_{-\infty}^{\infty} x(\tau) \delta(t-\tau) d \tau=x(t)
$$

- Thus, any function $x$ can be written in terms of an expression involving $\delta$.
- Moreover, $\delta$ is the convolutional identity. That is, for any function $x$,

$$
x * \delta=x
$$

- The convolution of two periodic functions is usually not well defined.
- This motivates an alternative notion of convolution for periodic functions known as periodic convolution.
- The periodic convolution of the $T$-periodic functions $x$ and $h$, denoted $x \circledast h$, is defined as

$$
x \circledast h(t)=\int_{T} x(\tau) h(t-\tau) d \tau
$$

where $\int_{T}$ denotes integration over an interval of length $T$.

- The periodic convolution and (linear) convolution of the $T$-periodic functions $x$ and $h$ are related as follows:

$$
x \circledast h(t)=x_{0} * h(t) \quad \text { where } \quad x(t)=\sum_{k=-\infty}^{\infty} x_{0}(t-k T)
$$

(i.e., $x_{0}(t)$ equals $x(t)$ over a single period of $x$ and is zero elsewhere).

## Section 4.2

## Convolution and LTI Systems

## Impulse Response

- The response $h$ of a system $\mathcal{H}$ to the input $\delta$ is called the impulse response of the system (i.e., $h=\mathcal{H} \delta$ ).
- For any LTI system with input $x$, output $y$, and impulse response $h$, the following relationship holds:

$$
y=x * h .
$$

- In other words, a LTI system simply computes a convolution.

■ Furthermore, a LTI system is completely characterized by its impulse response.

- That is, if the impulse response of a LTI system is known, we can determine the response of the system to any input.
- Since the impulse response of a LTI system is an extremely useful quantity, we often want to determine this quantity in a practical setting.
- Unfortunately, in practice, the impulse response of a system cannot be determined directly from the definition of the impulse response.


## Step Response

- The response $s$ of a system $\mathcal{H}$ to the input $u$ is called the step response of the system (i.e., $s=\mathcal{H} u$ ).
- The impulse response $h$ and step response $s$ of a LTI system are related as

$$
h(t)=\frac{d s(t)}{d t} .
$$

■ Therefore, the impulse response of a system can be determined from its step response by differentiation.

- The step response provides a practical means for determining the impulse response of a system.
- Often, it is convenient to represent a (CT) LTI system in block diagram form.
- Since such systems are completely characterized by their impulse response, we often label a system with its impulse response.
- That is, we represent a system with input $x$, output $y$, and impulse response $h$, as shown below.



## Interconnection of LTI Systems

- The series interconnection of the LTI systems with impulse responses $h_{1}$ and $h_{2}$ is the LTI system with impulse response $h_{1} * h_{2}$. That is, we have the equivalence shown below.

- The parallel interconnection of the LTI systems with impulse responses $h_{1}$ and $h_{2}$ is the LTI system with impulse response $h_{1}+h_{2}$. That is, we have the equivalence shown below.



## Section 4.3

## Properties of LTI Systems

## Memory

- A LTI system with impulse response $h$ is memoryless if and only if

$$
h(t)=0 \quad \text { for all } t \neq 0 .
$$

- That is, a LTI system is memoryless if and only if its impulse response $h$ is of the form

$$
h(t)=K \delta(t)
$$

where $K$ is a complex constant.

- Consequently, every memoryless LTI system with input $x$ and output $y$ is characterized by an equation of the form

$$
y=x *(K \delta)=K x
$$

(i.e., the system is an ideal amplifier).

- For a LTI system, the memoryless constraint is extremely restrictive (as every memoryless LTI system is an ideal amplifier).


## Causality

- A LTI system with impulse response $h$ is causal if and only if

$$
h(t)=0 \quad \text { for all } t<0
$$

(i.e., $h$ is a causal function).

- It is due to the above relationship that we call a function $x$, satisfying

$$
x(t)=0 \quad \text { for all } t<0
$$

a causal function.

## Invertibility

- The inverse of a LTI system, if such a system exists, is a LTI system.

■ Let $h$ and $h_{\text {inv }}$ denote the impulse responses of a LTI system and its (LTI) inverse, respectively. Then,

$$
h * h_{\mathrm{inv}}=\delta
$$

- Consequently, a LTI system with impulse response $h$ is invertible if and only if there exists a function $h_{\text {inv }}$ such that

$$
h * h_{\mathrm{inv}}=\delta
$$

- Except in simple cases, the above condition is often quite difficult to test.
- A LTI system with impulse response $h$ is BIBO stable if and only if

$$
\int_{-\infty}^{\infty}|h(t)| d t<\infty
$$

(i.e., $h$ is absolutely integrable).

## Eigenfunctions of LTI Systems

- As it turns out, every complex exponential is an eigenfunction of all LTI systems.
- For a LTI system $\mathcal{H}$ with impulse response $h$,

$$
\mathcal{H}\left\{e^{s t}\right\}(t)=H(s) e^{s t},
$$

where $s$ is a complex constant and

$$
H(s)=\int_{-\infty}^{\infty} h(t) e^{-s t} d t
$$

- That is, $e^{s t}$ is an eigenfunction of a LTI system and $H(s)$ is the corresponding eigenvalue.
- We refer to $H$ as the system function (or transfer function) of the system $\mathcal{H}$.
- From above, we can see that the response of a LTI system to a complex exponential is the same complex exponential multiplied by the complex factor $H(s)$.

■ Consider a LTI system with input $x$, output $y$, and system function $H$.

- Suppose that the input $x$ can be expressed as the linear combination of complex exponentials

$$
x(t)=\sum_{k} a_{k} e^{s_{k} t}
$$

where the $a_{k}$ and $s_{k}$ are complex constants.

- Using the fact that complex exponentials are eigenfunctions of LTI systems, we can conclude

$$
y(t)=\sum_{k} a_{k} H\left(s_{k}\right) e^{s_{k} t} .
$$

■ Thus, if an input to a LTI system can be expressed as a linear combination of complex exponentials, the output can also be expressed as a linear combination of the same complex exponentials.

- The above formula can be used to determine the output of a LTI system from its input in a way that does not require convolution.


## Part 5

## Continuous-Time Fourier Series (CTFS)

- The (CT) Fourier series is a representation for periodic functions.
- With a Fourier series, a function is represented as a linear combination of complex sinusoids.
- The use of complex sinusoids is desirable due to their numerous attractive properties.
- For example, complex sinusoids are continuous and differentiable. They are also easy to integrate and differentiate.
- Perhaps, most importantly, complex sinusoids are eigenfunctions of LTI systems.

Section 5.1

## Fourier Series

- A set of complex sinusoids is said to be harmonically related if there exists some constant $\omega_{0}$ such that the fundamental frequency of each complex sinusoid is an integer multiple of $\omega_{0}$.
- Consider the set of harmonically-related complex sinusoids given by

$$
\phi_{k}(t)=e^{j k \omega_{0} t} \quad \text { for all integer } k .
$$

- The fundamental frequency of the $k$ th complex sinusoid $\phi_{k}$ is $k \omega_{0}$, an integer multiple of $\omega_{0}$.
- Since the fundamental frequency of each of the harmonically-related complex sinusoids is an integer multiple of $\omega_{0}$, a linear combination of these complex sinusoids must be periodic.
- More specifically, a linear combination of these complex sinusoids is periodic with period $T=\frac{2 \pi}{\omega_{0}}$.


## CT Fourier Series

- A periodic (complex-valued) function $x$ with fundamental period $T$ and fundamental frequency $\omega_{0}=\frac{2 \pi}{T}$ can be represented as a linear combination of harmonically-related complex sinusoids as

$$
x(t)=\sum_{k=-\infty}^{\infty} c_{k} e^{j k \omega_{0} t}
$$

- Such a representation is known as (the complex exponential form of) a (CT) Fourier series, and the $c_{k}$ are called Fourier series coefficients.
- The above formula for $x$ is often referred to as the Fourier series synthesis equation.
- The terms in the summation for $k=K$ and $k=-K$ are called the $K$ th harmonic components, and have the fundamental frequency $K \omega_{0}$.
- To denote that a function $x$ has the Fourier series coefficient sequence $c_{k}$, we write

$$
x(t) \stackrel{\text { CTFS }}{\longleftrightarrow} c_{k} .
$$

## CT Fourier Series (Continued)

- The periodic function $x$ with fundamental period $T$ and fundamental frequency $\omega_{0}=\frac{2 \pi}{T}$ has the Fourier series coefficients $c_{k}$ given by

$$
c_{k}=\frac{1}{T} \int_{T} x(t) e^{-j k \omega_{0} t} d t
$$

where $\int_{T}$ denotes integration over an arbitrary interval of length $T$ (i.e., one period of $x$ ).

- The above equation for $c_{k}$ is often referred to as the Fourier series analysis equation.


## Section 5.2

## Convergence Properties of Fourier Series

## Remarks on Equality of Functions

- The equality of functions can be defined in more than one way.
- Two functions $x$ and $y$ are said to be equal in the pointwise sense if $x(t)=y(t)$ for all $t$ (i.e., $x$ and $y$ are equal at every point).
- Two functions $x$ and $y$ are said to be equal in the mean-squared error (MSE) sense if $\int|x(t)-y(t)|^{2} d t=0$ (i.e., the energy in $x-y$ is zero).
- Pointwise equality is a stronger condition than MSE equality (i.e., pointwise equality implies MSE equality but the converse is not true).
- Consider the functions

$$
\begin{gathered}
x_{1}(t)=1 \text { for all } t, \quad x_{2}(t)=1 \text { for all } t, \quad \text { and } \\
x_{3}(t)= \begin{cases}2 & t=0 \\
1 & \text { otherwise } .\end{cases}
\end{gathered}
$$

- The functions $x_{1}$ and $x_{2}$ are equal in both the pointwise sense and MSE sense.
- The functions $x_{1}$ and $x_{3}$ are equal in the MSE sense, but not in the pointwise sense.


## Convergence of Fourier Series

- Since a Fourier series can have an infinite number of (nonzero) terms, and an infinite sum may or may not converge, we need to consider the issue of convergence.
- That is, when we claim that a periodic function $x$ is equal to the Fourier series $\sum_{k=-\infty}^{\infty} c_{k} e^{j k \omega_{0} t}$, is this claim actually correct?
- Consider a periodic function $x$ that we wish to represent with the Fourier series

$$
\sum_{k=-\infty}^{\infty} c_{k} e^{j k \omega_{0} t}
$$

- Let $x_{N}$ denote the Fourier series truncated after the $N$ th harmonic components as given by

$$
x_{N}(t)=\sum_{k=-N}^{N} c_{k} e^{j k \omega_{0} t}
$$

- Here, we are interested in whether $\lim _{N \rightarrow \infty} x_{N}$ is equal (in some sense) to $x$.


## Convergence of Fourier Series (Continued)

- Again, let $x_{N}$ denote the Fourier series for the periodic function $x$ truncated after the $N$ th harmonic components as given by

$$
x_{N}(t)=\sum_{k=-N}^{N} c_{k} e^{j k \omega_{0} t}
$$

- If $\lim _{N \rightarrow \infty} x_{N}(t)=x(t)$ for all $t$ (i.e., $\lim _{N \rightarrow \infty} x_{N}$ is equal to $x$ in the pointwise sense), the Fourier series is said to converge pointwise to $x$.
- If convergence is pointwise and the rate of convergence is the same everywhere, the convergence is said to be uniform.
- If $\lim _{N \rightarrow \infty} \frac{1}{T} \int_{T}\left|x_{N}(t)-x(t)\right|^{2} d t=0$ (i.e., $\lim _{N \rightarrow \infty} x_{N}$ is equal to $x$ in the MSE sense), the Fourier series is said to converge to $x$ in the MSE sense.
- Pointwise convergence is a stronger condition than MSE convergence (i.e., pointwise convergence implies MSE convergence, but the converse is not true).


## Convergence of Fourier Series: Continuous Case

- If a periodic function $x$ is continuous and its Fourier series coefficients $c_{k}$ are absolutely summable (i.e., $\sum_{k=-\infty}^{\infty}\left|c_{k}\right|<\infty$ ), then the Fourier series representation of $x$ converges uniformly (i.e., pointwise at the same rate everywhere).
- Since, in practice, we often encounter functions with discontinuities (e.g., a square wave), the above result is of somewhat limited value.


## Convergence of Fourier Series: Finite-Energy Case

- If a periodic function $x$ has finite energy in a single period (i.e., $\left.\int_{T}|x(t)|^{2} d t<\infty\right)$, the Fourier series converges in the MSE sense.
- Since, in situations of practical interest, the finite-energy condition in the above theorem is typically satisfied, the theorem is usually applicable.
- It is important to note, however, that MSE convergence (i.e., $E=0$ ) does not necessarily imply pointwise convergence (i.e., $\tilde{x}(t)=x(t)$ for all $t$ ).
- Thus, the above convergence theorem does not provide much useful information regarding the value of $\tilde{x}(t)$ for specific values of $t$.
- Consequently, the above theorem is typically most useful for simply determining if the Fourier series converges.


## Dirichlet Conditions

- The Dirichlet conditions for the periodic function $x$ are as follows:

1 over a single period, $x$ is absolutely integrable (i.e., $\int_{T}|x(t)| d t<\infty$ );
2 over a single period, $x$ has a finite number of maxima and minima (i.e., $x$ is of bounded variation); and
3 over any finite interval, $x$ has a finite number of discontinuities, each of which is finite.

- Examples of functions violating the Dirichlet conditions are shown below.





## Convergence of Fourier Series: Dirichlet Case

- If a periodic function $x$ satisfies the Dirichlet conditions, then:

1 the Fourier series converges pointwise everywhere to $x$, except at the points of discontinuity of $x$; and
2 at each point $t_{a}$ of discontinuity of $x$, the Fourier series $\tilde{x}$ converges to

$$
\tilde{x}\left(t_{a}\right)=\frac{1}{2}\left[x\left(t_{a}^{-}\right)+x\left(t_{a}^{+}\right)\right],
$$

where $x\left(t_{a}^{-}\right)$and $x\left(t_{a}^{+}\right)$denote the values of the function $x$ on the left- and right-hand sides of the discontinuity, respectively.

- Since most functions tend to satisfy the Dirichlet conditions and the above convergence result specifies the value of the Fourier series at every point, this result is often very useful in practice.


## Gibbs Phenomenon

- In practice, we frequently encounter functions with discontinuities.
- When a function $x$ has discontinuities, the Fourier series representation of $x$ does not converge uniformly (i.e., at the same rate everywhere).
- The rate of convergence is much slower at points in the vicinity of a discontinuity.
- Furthermore, in the vicinity of a discontinuity, the truncated Fourier series $x_{N}$ exhibits ripples, where the peak amplitude of the ripples does not seem to decrease with increasing $N$.
- As it turns out, as $N$ increases, the ripples get compressed towards discontinuity, but, for any finite $N$, the peak amplitude of the ripples remains approximately constant.
- This behavior is known as Gibbs phenomenon.
- The above behavior is one of the weaknesses of Fourier series (i.e., Fourier series converge very slowly near discontinuities).


## Gibbs Phenomenon: Periodic Square Wave Example



Fourier series truncated after the 3rd harmonic components


Fourier series truncated after the 11th harmonic components


Fourier series truncated after the 7th harmonic components


Fourier series truncated after the 101st harmonic components

## Section 5.3

## Properties of Fourier Series

$$
x(t) \stackrel{\text { cTFs }}{\leftrightarrows} a_{k} \quad \text { and } \quad y(t) \stackrel{\text { cTFs }}{\longleftrightarrow} b_{k}
$$

| Property | Time Domain | Fourier Domain |
| :--- | :--- | :--- |
| Linearity | $\alpha x(t)+\beta y(t)$ | $\alpha a_{k}+\beta b_{k}$ |
| Translation | $x\left(t-t_{0}\right)$ | $e^{-j k(2 \pi / T) t_{0}} a_{k}$ |
| Modulation | $e^{j M(2 \pi / T) t} x(t)$ | $a_{k-M}$ |
| Reflection | $x(-t)$ | $a_{-k}$ |
| Conjugation | $x^{*}(t)$ | $a_{-k}^{*}$ |
| Periodic Convolution | $x \circledast y(t)$ | $T a_{k} b_{k}$ |
| Multiplication | $x(t) y(t)$ | $\sum_{n=-\infty}^{\infty} a_{n} b_{k-n}$ |


| Property |  |
| :--- | :--- |
| Parseval's Relation | $\frac{1}{T} \int_{T}\|x(t)\|^{2} d t=\sum_{k=-\infty}^{\infty}\left\|a_{k}\right\|^{2}$ |
| Even Symmetry | $x$ is even $\Leftrightarrow a$ is even |
| Odd Symmetry | $x$ is odd $\Leftrightarrow a$ is odd |
| Real / Conjugate Symmetry | $x$ is real $\Leftrightarrow a$ is conjugate symmetric |

- Let $x$ and $y$ be two periodic functions with the same period. If $x(t) \stackrel{\text { cTFS }}{\longleftrightarrow} a_{k}$ and $y(t) \stackrel{\text { CTFS }}{\longleftrightarrow} b_{k}$, then

$$
\alpha x(t)+\beta y(t) \stackrel{\text { CTFS }}{\longleftrightarrow} \alpha a_{k}+\beta b_{k},
$$

where $\alpha$ and $\beta$ are complex constants.

- That is, a linear combination of functions produces the same linear combination of their Fourier series coefficients.
- Let $x$ denote a periodic function with period $T$ and the corresponding frequency $\omega_{0}=2 \pi / T$. If $x(t) \stackrel{\text { cTFs }}{\leftrightarrows} c_{k}$, then

$$
x\left(t-t_{0}\right) \stackrel{\text { CTFS }}{\longleftrightarrow} e^{-j k \omega_{0} t_{0}} c_{k}=e^{-j k(2 \pi / T) t_{0}} c_{k},
$$

where $t_{0}$ is a real constant.

- In other words, time shifting a periodic function changes the argument (but not magnitude) of its Fourier series coefficients.
- Let $x$ denote a periodic function with period $T$ and the corresponding frequency $\omega_{0}=2 \pi / T$. If $x(t) \stackrel{\text { cTFS }}{\longleftrightarrow} c_{k}$, then

$$
e^{j M(2 \pi / T) t} x(t)=e^{j M \omega_{0} t} x(t) \stackrel{\text { cTFS }}{\longleftrightarrow} c_{k-M}
$$

where $M$ is an integer constant.

- In other words, multiplying a periodic function by $e^{j M \omega_{0} t}$ shifts the Fourier-series coefficient sequence.
- Let $x$ denote a periodic function with period $T$ and the corresponding frequency $\omega_{0}=2 \pi / T$. If $x(t) \stackrel{\text { CTFS }}{\leftrightarrows} c_{k}$, then

$$
x(-t) \stackrel{\text { cTFs }}{\longleftrightarrow} c_{-k} .
$$

- That is, time reversal of a function results in a time reversal of its Fourier series coefficients.


## Conjugation

- For a $T$-periodic function $x$ with Fourier series coefficient sequence $c$, the following property holds:

$$
x^{*}(t) \stackrel{\text { CTFS }}{\longleftrightarrow} c_{-k}^{*}
$$

- In other words, conjugating a function has the effect of time reversing and conjugating the Fourier series coefficient sequence.
- Let $x$ and $y$ be two periodic functions with the same period $T$. If $x(t) \stackrel{\text { CTFS }}{\longleftrightarrow} a_{k}$ and $y(t) \stackrel{\text { CTFS }}{\longleftrightarrow} b_{k}$, then

$$
x \circledast y(t) \stackrel{\text { CTFS }}{\longleftrightarrow} T a_{k} b_{k} .
$$

- In other words, periodic convolution of two functions corresponds to the multiplication (up to a scale factor) of their Fourier-series coefficient sequences.
- Let $x$ and $y$ be two periodic functions with the same period. If $x(t) \stackrel{\text { cTFS }}{\longleftrightarrow} a_{k}$ and $y(t) \stackrel{\text { CTFS }}{\longleftrightarrow} b_{k}$, then

$$
x(t) y(t) \stackrel{\text { CTFS }}{\longleftrightarrow} \sum_{n=-\infty}^{\infty} a_{n} b_{k-n}
$$

- As we shall see later, the above summation is the DT convolution of $a$ and $b$.
- In other words, the multiplication of two periodic functions corresponds to the DT convolution of their corresponding Fourier-series coefficient sequences.
- A function $x$ and its Fourier series coefficient sequence $a$ satisfy the following relationship:

$$
\frac{1}{T} \int_{T}|x(t)|^{2} d t=\sum_{k=-\infty}^{\infty}\left|a_{k}\right|^{2}
$$

- The above relationship is simply stating that the amount of energy in $x$ (i.e., $\frac{1}{T} \int_{T}|x(t)|^{2} d t$ ) and the amount of energy in the Fourier series coefficient sequence $a$ (i.e., $\sum_{k=-\infty}^{\infty}\left|a_{k}\right|^{2}$ ) are equal.
- In other words, the transformation between a function and its Fourier series coefficient sequence preserves energy.
- For a periodic function $x$ with Fourier series coefficient sequence $c$, the following properties hold:

$$
\begin{gathered}
x \text { is even } \Leftrightarrow c \text { is even; and } \\
x \text { is odd } \Leftrightarrow c \text { is odd. }
\end{gathered}
$$

- In other words, the even/odd symmetry properties of $x$ and $c$ always match.
- A function $x$ is real if and only if its Fourier series coefficient sequence $c$ satisfies

$$
c_{k}=c_{-k}^{*} \text { for all } k
$$

(i.e., $c$ is conjugate symmetric).

- Thus, for a real-valued function, the negative-indexed Fourier series coefficients are redundant, as they are completely determined by the nonnegative-indexed coefficients.
- From properties of complex numbers, one can show that $c_{k}=c_{-k}^{*}$ is equivalent to

$$
\left|c_{k}\right|=\left|c_{-k}\right| \quad \text { and } \quad \arg c_{k}=-\arg c_{-k}
$$

(i.e., $\left|c_{k}\right|$ is even and $\arg c_{k}$ is odd).

- Note that $x$ being real does not necessarily imply that $c$ is real.


## Trigonometric Forms of a Fourier Series

- Consider the periodic function $x$ with the Fourier series coefficients $c_{k}$.
- If $x$ is real, then its Fourier series can be rewritten in two other forms, known as the combined trigonometric and trigonometric forms.
- The combined trigonometric form of a Fourier series has the appearance

$$
x(t)=c_{0}+2 \sum_{k=1}^{\infty}\left|c_{k}\right| \cos \left(k \omega_{0} t+\theta_{k}\right)
$$

where $\theta_{k}=\arg c_{k}$.

- The trigonometric form of a Fourier series has the appearance

$$
x(t)=c_{0}+\sum_{k=1}^{\infty}\left[\alpha_{k} \cos \left(k \omega_{0} t\right)+\beta_{k} \sin \left(k \omega_{0} t\right)\right]
$$

where $\alpha_{k}=2 \operatorname{Re} c_{k}$ and $\beta_{k}=-2 \operatorname{Im} c_{k}$.

- Note that the trigonometric forms contain only real quantities.


## Other Properties of Fourier Series

- For a $T$-periodic function $x$ with Fourier-series coefficient sequence $c$, the following properties hold:
$\square c_{0}$ is the average value of $x$ over a single period $T$;
$\boxed{x}$ is real and even $\Leftrightarrow c$ is real and even; and
$3 x$ is real and odd $\Leftrightarrow c$ is purely imaginary and odd.


## Section 5.4

## Fourier Series and Frequency Spectra

## A New Perspective on Functions: The Frequency Domain

- The Fourier series provides us with an entirely new way to view functions.
- Instead of viewing a function as having information distributed with respect to time (i.e., a function whose domain is time), we view a function as having information distributed with respect to frequency (i.e., a function whose domain is frequency).
- This so called frequency-domain perspective is of fundamental importance in engineering.
- Many engineering problems can be solved much more easily using the frequency domain than the time domain.
- The Fourier series coefficients of a function $x$ provide a means to quantify how much information $x$ has at different frequencies.
- The distribution of information in a function over different frequencies is referred to as the frequency spectrum of the function.


## Motivating Example

- Consider the real 1-periodic function $x$ having the Fourier series representation

$$
\begin{gathered}
x(t)=-\frac{j}{10} e^{-j 14 \pi t}-\frac{2 j}{10} e^{-j 10 \pi t}-\frac{4 j}{10} e^{-j 6 \pi t}-\frac{13 j}{10} e^{-j 2 \pi t} \\
+\frac{13 j}{10} e^{j 2 \pi t}+\frac{4 j}{10} e^{j 6 \pi t}+\frac{2 j}{10} e^{j 10 \pi t}+\frac{j}{10} e^{j 14 \pi t}
\end{gathered}
$$

- A plot of $x$ is shown below.

- The terms that make the most dominant contribution to the overall sum are the ones with the largest magnitude coefficients.
- To illustrate this, we consider the problem of determining the best approximation of $x$ that keeps only 4 of the 8 terms in the Fourier series.


## Motivating Example (Continued)



Approximation using the 4 terms with the largest magnitude coefficients


Approximation using the 4 terms with the smallest magnitude nonzero coefficients

- To gain further insight into the role played by the Fourier series coefficients $c_{k}$ in the context of the frequency spectrum of the function $x$, it is helpful to write the Fourier series with the $c_{k}$ expressed in polar form as follows:

$$
x(t)=\sum_{k=-\infty}^{\infty} c_{k} e^{j k \omega_{0} t}=\sum_{k=-\infty}^{\infty}\left|c_{k}\right| e^{j\left(k \omega_{0} t+\arg c_{k}\right)}
$$

- Clearly, the $k$ th term in the summation corresponds to a complex sinusoid with fundamental frequency $k \omega_{0}$ that has been amplitude scaled by a factor of $\left|c_{k}\right|$ and time shifted by an amount that depends on $\arg c_{k}$.
- For a given $k$, the larger $\left|c_{k}\right|$ is, the larger is the amplitude of its corresponding complex sinusoid $e^{j k \omega_{0} t}$, and therefore the larger the contribution the $k$ th term (which is associated with frequency $k \omega_{0}$ ) will make to the overall summation.
- In this way, we can use $\left|c_{k}\right|$ as a measure of how much information a function $x$ has at the frequency $k \omega_{0}$.
- The Fourier series coefficients $c_{k}$ are referred to as the frequency spectrum of $x$.
■ The magnitudes $\left|c_{k}\right|$ of the Fourier series coefficients are referred to as the magnitude spectrum of $x$.
- The arguments $\arg c_{k}$ of the Fourier series coefficients are referred to as the phase spectrum of $x$.
- Normally, the spectrum of a function is plotted against frequency $k \omega_{0}$ instead of $k$.
- Since the Fourier series only has frequency components at integer multiples of the fundamental frequency, the frequency spectrum is discrete in the independent variable (i.e., frequency).
- Due to the general appearance of frequency-spectrum plot (i.e., a number of vertical lines at various frequencies), we refer to such spectra as lime spectra.


## Frequency Spectra of Real Functions

- Recall that, for a real function $x$, the Fourier series coefficient sequence $c$ satisfies

$$
c_{k}=c_{-k}^{*}
$$

(i.e., $c$ is conjugate symmetric), which is equivalent to

$$
\left|c_{k}\right|=\left|c_{-k}\right| \quad \text { and } \quad \arg c_{k}=-\arg c_{-k}
$$

■ Since $\left|c_{k}\right|=\left|c_{-k}\right|$, the magnitude spectrum of a real function is always even.

- Similarly, since $\arg c_{k}=-\arg c_{-k}$, the phase spectrum of a real function is always odd.
- Due to the symmetry in the frequency spectra of real functions, we typically ignore negative frequencies when dealing with such functions.
- In the case of functions that are complex but not real, frequency spectra do not possess the above symmetry, and negative frequencies become important.


## Section 5.5

## Fourier Series and LTI Systems

- Recall that a LTI system $\mathcal{H}$ with impulse response $h$ is such that $\mathcal{H}\left\{e^{s t}\right\}(t)=H_{\mathrm{L}}(s) e^{s t}$, where $H_{\mathrm{L}}(s)=\int_{-\infty}^{\infty} h(t) e^{-s t} d t$. (That is, complex exponentials are eigenfunctions of LTI systems.)
- Since a complex sinusoid is a special case of a complex exponential, we can reuse the above result for the special case of complex sinusoids.
- For a LTI system $\mathcal{H}$ with impulse response $h$,

$$
\mathcal{H}\left\{e^{j \omega t}\right\}(t)=H(\omega) e^{j \omega t}
$$

where $\omega$ is a real constant and

$$
H(\omega)=\int_{-\infty}^{\infty} h(t) e^{-j \omega t} d t
$$

- That is, $e^{j \omega t}$ is an eigenfunction of a LTI system and $H(\boldsymbol{\omega})$ is the corresponding eigenvalue.
- We refer to $H$ as the frequency response of the system $\mathcal{H}$.
- Consider a LTI system with input $x$, output $y$, and frequency response $H$.
- Suppose that the $T$-periodic input $x$ is expressed as the Fourier series

$$
x(t)=\sum_{k=-\infty}^{\infty} c_{k} e^{j k \omega_{0} t}, \quad \text { where } \omega_{0}=\frac{2 \pi}{T} .
$$

- Using our knowledge about the eigenfunctions of LTI systems, we can conclude

$$
y(t)=\sum_{k=-\infty}^{\infty} c_{k} H\left(k \omega_{0}\right) e^{j k \omega_{0} t} .
$$

- Thus, if the input $x$ to a LTI system is a Fourier series, the output $y$ is also a Fourier series. More specifically, if $x(t) \stackrel{\text { CTFS }}{\longleftrightarrow} c_{k}$ then $y(t) \stackrel{\text { CTFS }}{\longleftrightarrow} H\left(k \omega_{0}\right) c_{k}$.
- The above formula can be used to determine the output of a LTI system from its input in a way that does not require convolution.
- In many applications, we want to modify the spectrum of a function by either amplifying or attenuating certain frequency components.
- This process of modifying the frequency spectrum of a function is called filtering.
- A system that performs a filtering operation is called a filter.
- Many types of filters exist.
- Frequency selective filters pass some frequencies with little or no distortion, while significantly attenuating other frequencies.
- Several basic types of frequency-selective filters include: lowpass, highpass, and bandpass.


## Ideal Lowpass Filter

- An ideal lowpass filter eliminates all frequency components with a frequency whose magnitude is greater than some cutoff frequency, while leaving the remaining frequency components unaffected.
- Such a filter has a frequency response of the form

$$
H(\omega)= \begin{cases}1 & |\omega| \leq \omega_{c} \\ 0 & \text { otherwise }\end{cases}
$$

where $\omega_{c}$ is the cutoff frequency.

- A plot of this frequency response is given below.



## Ideal Highpass Filter

- An ideal highpass filter eliminates all frequency components with a frequency whose magnitude is less than some cutoff frequency, while leaving the remaining frequency components unaffected.
- Such a filter has a frequency response of the form

$$
H(\omega)= \begin{cases}1 & |\omega| \geq \omega_{c} \\ 0 & \text { otherwise }\end{cases}
$$

where $\omega_{c}$ is the cutoff frequency.

- A plot of this frequency response is given below.



## Ideal Bandpass Filter

- An ideal bandpass filter eliminates all frequency components with a frequency whose magnitude does not lie in a particular range, while leaving the remaining frequency components unaffected.
- Such a filter has a frequency response of the form

$$
H(\omega)= \begin{cases}1 & \omega_{c 1} \leq|\omega| \leq \omega_{c 2} \\ 0 & \text { otherwise }\end{cases}
$$

where the limits of the passband are $\omega_{c 1}$ and $\omega_{c 2}$.

- A plot of this frequency response is given below.



## Part 6

## Continuous-Time Fourier Transform (CTFT)

- The (CT) Fourier series provide an extremely useful representation for periodic functions.
- Often, however, we need to deal with functions that are not periodic.
- A more general tool than the Fourier series is needed in this case.
- The (CT) Fourier transform can be used to represent both periodic and aperiodic functions.
- Since the Fourier transform is essentially derived from Fourier series through a limiting process, the Fourier transform has many similarities with Fourier series.

Section 6.1

## Fourier Transform

## Development of the Fourier Transform [Aperiodic Case]

- The (CT) Fourier series is an extremely useful function representation.
- Unfortunately, this function representation can only be used for periodic functions, since a Fourier series is inherently periodic.
- Many functions are not periodic, however.
- Rather than abandoning Fourier series, one might wonder if we can somehow use Fourier series to develop a representation that can be applied to aperiodic functions.
- By viewing an aperiodic function as the limiting case of a $T$-periodic function where $T \rightarrow \infty$, we can use the Fourier series to develop a function representation that can be used for aperiodic functions, known as the Fourier transform.
- Recall that the Fourier series representation of a $T$-periodic function $x$ is given by

$$
x(t)=\sum_{k=-\infty}^{\infty} \underbrace{\left(\frac{1}{T} \int_{-T / 2}^{T / 2} x(\tau) e^{-j k(2 \pi / T) \tau} d \tau\right)}_{c_{k}} e^{j k(2 \pi / T) t}
$$

- In the above representation, if we take the limit as $T \rightarrow \infty$, we obtain

$$
x(t)=\frac{1}{2 \pi} \int_{-\infty}^{\infty} \underbrace{\left(\int_{-\infty}^{\infty} x(\tau) e^{-j \omega \tau} d \tau\right)}_{X(\omega)} e^{j \omega t} d \omega
$$

(i.e., as $T \rightarrow \infty$, the outer summation becomes an integral, $\frac{1}{T}$ becomes $\frac{1}{2 \pi} d \omega$, and $\left(\frac{2 \pi}{T}\right) k$ becomes $\omega$ ).

- This representation for aperiodic functions is known as the Fourier transform representation.
- The classical Fourier transform for aperiodic functions does not exist (i.e., $\int_{-\infty}^{\infty} x(t) e^{-j \omega t} d t$ fails to converge) for some functions of great practical interest, such as:
- a nonzero constant function;
- a periodic function (e.g., a real or complex sinusoid);
- the unit-step function (i.e., $u$ ); and
- the signum function (i.e., sgn).
- Fortunately, the Fourier transform can be extended to handle such functions, resulting in what is known as the generalized Fourier transform.
- For our purposes, we can think of the classical and generalized Fourier transforms as being defined by the same formulas.
- Therefore, in what follows, we will not typically make a distinction between the classical and generalized Fourier transforms.


## CT Fourier Transform (CTFT)

- The (CT) Fourier transform of the function $x$, denoted $\mathcal{F} x$ or $X$, is given by

$$
\mathcal{F} x(\omega)=X(\omega)=\int_{-\infty}^{\infty} x(t) e^{-j \omega t} d t
$$

- The preceding equation is sometimes referred to as Fourier transform analysis equation (or forward Fourier transform equation).
- The inverse Fourier transform of $X$, denoted $\mathcal{F}^{-1} X$ or $x$, is given by

$$
\mathcal{F}^{-1} X(t)=x(t)=\frac{1}{2 \pi} \int_{-\infty}^{\infty} X(\omega) e^{j \omega t} d \omega
$$

- The preceding equation is sometimes referred to as the Fourier transform synthesis equation (or inverse Fourier transform equation).
- As a matter of notation, to denote that a function $x$ has the Fourier transform $X$, we write $x(t) \stackrel{\text { CTFT }}{\longleftrightarrow} X(\omega)$.
- A function $x$ and its Fourier transform $X$ constitute what is called a Fourier transform pair.
- For a function $x$, the Fourier transform of $x$ is denoted using operator notation as $\mathcal{F} x$.
- The Fourier transform of $x$ evaluated at $\omega$ is denoted $\mathcal{F} x(\omega)$.
- Note that $\mathcal{F} x$ is a function, whereas $\mathcal{F} x(\omega)$ is a number.
- Similarly, for a function $X$, the inverse Fourier transform of $X$ is denoted using operator notation as $\mathcal{F}^{-1} X$.
- The inverse Fourier transform of $X$ evaluated at $t$ is denoted $\mathcal{F}^{-1} X(t)$.
- Note that $\mathcal{F}^{-1} X$ is a function, whereas $\mathcal{F}^{-1} X(t)$ is a number.
- With the above said, engineers often abuse notation, and use expressions like those above to mean things different from their proper meanings.
- Since such notational abuse can lead to problems, it is strongly recommended that one refrain from doing this.
- Often, we would like to write an expression for the Fourier transform of a function without explicitly naming the function.
- For example, consider writing an expression for the Fourier transform of the function $v(t)=x(5 t-3)$ but without using the name " $v$ ".
- It would be incorrect to write " $\mathcal{F} x(5 t-3)$ " as this is the function $\mathcal{F} x$ evaluated at $5 t-3$, which is not the meaning that we wish to convey.
- Also, strictly speaking, it would be incorrect to write "F$\{x(5 t-3)\}$ " as the operand of the Fourier transform operator must be a function, and $x(5 t-3)$ is a number (i.e., the function $x$ evaluated at $5 t-3$ ).
- Using dot notation, we can write the following strictly-correct expression for the desired Fourier transform: $\mathcal{F} x(5 \cdot-3)$.
- In many cases, however, it is probably advisable to avoid employing anonymous (i.e., unnamed) functions, as their use tends to be more error prone in some contexts.
- Since dot notation is less frequently used by engineers, the author has elected to minimize its use herein.
- To avoid ambiguous notation, the following conventions are followed:

11 in the expression for the operand of a Fourier transform operator, the independent variable is assumed to be the variable named " $t$ " unless otherwise indicated (i.e., in terms of dot notation, each " $t$ " is treated as if it were a ".")
2 in the expression for the operand of the inverse Fourier transform operator, the independent variable is assumed to be the variable named " $\omega$ " unless otherwise indicated (i.e., in terms of dot notation, each " $\omega$ " is treated as if it were a ".").

- For example, with these conventions:
- "F $\{\cos (t-\tau)\}$ " denotes the function that is the Fourier transform of the function $v(t)=\cos (t-\tau)$ (not the Fourier transform of the function $v(\tau)=\cos (t-\tau))$.
- " $\mathcal{F}^{-1}\{\delta(3 \omega-\lambda)\}$ " denotes the function that is the inverse Fourier transform of the function $V(\omega)=\delta(3 \omega-\lambda)$ (not the inverse Fourier transform of the function $V(\lambda)=\delta(3 \omega-\lambda))$.


## Section 6.2

## Convergence Properties of the Fourier Transform

## Convergence of the Fourier Transform

- Consider an arbitrary function $x$.
- The function $x$ has the Fourier transform representation $\tilde{x}$ given by

$$
\tilde{x}(t)=\frac{1}{2 \pi} \int_{-\infty}^{\infty} X(\omega) e^{j \omega t} d \omega, \quad \text { where } \quad X(\omega)=\int_{-\infty}^{\infty} x(t) e^{-j \omega t} d t
$$

- Now, we need to concern ourselves with the convergence properties of this representation.
- In other words, we want to know when $\tilde{x}$ is a valid representation of $x$.
- Since the Fourier transform is essentially derived from Fourier series, the convergence properties of the Fourier transform are closely related to the convergence properties of Fourier series.


## Convergence of the Fourier Transform: Continuous Case

- If a function $x$ is continuous and absolutely integrable (i.e., $\left.\int_{-\infty}^{\infty}|x(t)| d t<\infty\right)$ and the Fourier transform $X$ of $x$ is absolutely integrable (i.e., $\int_{-\infty}^{\infty}|X(\omega)| d \omega<\infty$ ), then the Fourier transform representation of $x$ converges pointwise (i.e., $x(t)=\frac{1}{2 \pi} \int_{-\infty}^{\infty}\left[\int_{-\infty}^{\infty} x(t) e^{-j \omega t} d t\right] e^{j \omega t} d \omega$ for all $t$ ).
- Since, in practice, we often encounter functions with discontinuities (e.g., a rectangular pulse), the above result is sometimes of limited value.


## Convergence of the Fourier Transform: Finite-Energy Case

- If a function $x$ is of finite energy (i.e., $\int_{-\infty}^{\infty}|x(t)|^{2} d t<\infty$ ), then its Fourier transform representation converges in the MSE sense.
- In other words, if $x$ is of finite energy, then the energy $E$ in the difference function $\tilde{x}-x$ is zero; that is,

$$
E=\int_{-\infty}^{\infty}|\tilde{x}(t)-x(t)|^{2} d t=0
$$

- Since, in situations of practical interest, the finite-energy condition in the above theorem is often satisfied, the theorem is frequently applicable.
- It is important to note, however, that the condition $E=0$ does not necessarily imply $\tilde{x}(t)=x(t)$ for all $t$.
- Thus, the above convergence result does not provide much useful information regarding the value of $\tilde{x}(t)$ at specific values of $t$.
- Consequently, the above theorem is typically most useful for simply determining if the Fourier transform representation converges.


## Dirichlet Conditions

- The Dirichlet conditions for the function $x$ are as follows:

1 the function $x$ is absolutely integrable (i.e., $\int_{-\infty}^{\infty}|x(t)| d t<\infty$ );
2 on any finite interval, $x$ has a finite number of maxima and minima (i.e., $x$ is of bounded variation); and
3 on any finite interval, $x$ has a finite number of discontinuities and each discontinuity is itself finite.
■ Examples of functions violating the Dirichlet conditions are shown below.




## Convergence of the Fourier Transform: Dirichlet Case

- If a function $x$ satisfies the Dirichlet conditions, then:

1 the Fourier transform representation $\tilde{x}$ converges pointwise everywhere to $x$, except at the points of discontinuity of $x$; and
2 at each point $t_{a}$ of discontinuity of $x$, the Fourier transform representation $\tilde{x}$ converges to

$$
\tilde{x}\left(t_{a}\right)=\frac{1}{2}\left[x\left(t_{a}^{-}\right)+x\left(t_{a}^{+}\right)\right]
$$

where $x\left(t_{a}^{-}\right)$and $x\left(t_{a}^{+}\right)$denote the values of the function $x$ on the left- and right-hand sides of the discontinuity, respectively.

- Since most functions tend to satisfy the Dirichlet conditions and the above convergence result specifies the value of the Fourier transform representation at every point, this result is often very useful in practice.


## Section 6.3

## Properties of the Fourier Transform

| Property | Time Domain | Frequency Domain |
| :--- | :--- | :--- |
| Linearity | $a_{1} x_{1}(t)+a_{2} x_{2}(t)$ | $a_{1} X_{1}(\omega)+a_{2} X_{2}(\omega)$ |
| Time-Domain Shifting | $x\left(t-t_{0}\right)$ | $e^{-j \omega t_{0}} X(\omega)$ |
| Frequency-Domain Shifting | $e^{j \omega_{0} t} x(t)$ | $X\left(\omega-\omega_{0}\right)$ |
| Time/Frequency-Domain Scaling | $x(a t)$ | $\frac{1}{\|a\|} X\left(\frac{\omega}{a}\right)$ |
| Conjugation | $x^{*}(t)$ | $X^{*}(-\omega)$ |
| Duality | $X(t)$ | $2 \pi x(-\omega)$ |
| Time-Domain Convolution | $x_{1} * x_{2}(t)$ | $X_{1}(\omega) X_{2}(\omega)$ |
| Time-Domain Multiplication | $x_{1}(t) x_{2}(t)$ | $\frac{1}{2 \pi} X_{1} * X_{2}(\omega)$ |
| Time-Domain Differentiation | $\frac{d}{d t} x(t)$ | $j \omega X(\omega)$ |
| Frequency-Domain Differentiation | $t x(t)$ | $j \frac{d}{d \omega} X(\omega)$ |
| Time-Domain Integration | $\int_{-\infty}^{t} x(\tau) d \tau$ | $\frac{1}{j \omega} X(\omega)+\pi X(0) \delta(\omega)$ |


| Property |  |
| :--- | :--- |
| Parseval's Relation | $\int_{-\infty}^{\infty}\|x(t)\|^{2} d t=\frac{1}{2 \pi} \int_{-\infty}^{\infty}\|X(\omega)\|^{2} d \omega$ |
| Even Symmetry | $x$ is even $\Leftrightarrow X$ is even |
| Odd Symmetry | $x$ is odd $\Leftrightarrow X$ is odd |
| Real / Conjugate Symmetry | $x$ is real $\Leftrightarrow X$ is conjugate symmetric |

## (CT) Fourier Transform Pairs

| Pair | $x(t)$ | $X(\omega)$ |
| :--- | :--- | :--- |
| 1 | $\delta(t)$ | 1 |
| 2 | $u(t)$ | $\pi \delta(\omega)+\frac{1}{j \omega}$ |
| 3 | 1 | $2 \pi \delta(\omega)$ |
| 4 | $\operatorname{sgn}(t)$ | $\frac{2}{j \omega}$ |
| 5 | $e^{j \omega_{0} t}$ | $2 \pi \delta\left(\omega-\omega_{0}\right)$ |
| 6 | $\cos \left(\omega_{0} t\right)$ | $\pi\left[\delta\left(\omega-\omega_{0}\right)+\delta\left(\omega+\omega_{0}\right)\right]$ |
| 7 | $\sin \left(\omega_{0} t\right)$ | $\frac{\pi}{j}\left[\delta\left(\omega-\omega_{0}\right)-\delta\left(\omega+\omega_{0}\right)\right]$ |
| 8 | $\operatorname{rect}(t / T)$ | $\|T\| \operatorname{sinc}(T \omega / 2)$ |
| 9 | $\operatorname{sinc}(B t)$ | $\frac{\pi}{\|B\|} \operatorname{rect}\left(\frac{\omega}{2 B}\right)$ |
| 10 | $e^{-a t} u(t), \operatorname{Re}\{a\}>0$ | $\frac{1}{a+j \omega}$ |
| 11 | $t^{n-1} e^{-a t} u(t), \operatorname{Re}\{a\}>0$ | $\frac{(n-1)!}{(a+j \omega)^{n}}$ |
| 12 | $\operatorname{tri}(t / T)$ | $\frac{\|T\|}{2} \operatorname{sinc}^{2}(T \omega / 4)$ |

■ If $x_{1}(t) \stackrel{\text { cTrT }}{\longleftrightarrow} X_{1}(\omega)$ and $x_{2}(t) \stackrel{\text { ctrT }}{\longleftrightarrow} X_{2}(\omega)$, then

$$
a_{1} x_{1}(t)+a_{2} x_{2}(t) \stackrel{\text { cTFT }}{\longleftrightarrow} a_{1} X_{1}(\omega)+a_{2} X_{2}(\omega),
$$

where $a_{1}$ and $a_{2}$ are arbitrary complex constants.

- This is known as the linearity property of the Fourier transform.


## Time-Domain Shifting (Translation)

- If $x(t) \stackrel{\text { cTrT }}{\longleftrightarrow} X(\omega)$, then

$$
x\left(t-t_{0}\right) \stackrel{\text { cTFT }}{\longleftrightarrow} e^{-j \omega t_{0}} X(\omega),
$$

where $t_{0}$ is an arbitrary real constant.

- This is known as the translation (or time-domain shifting) property of the Fourier transform.
- If $x(t) \stackrel{\text { CTFT }}{\longleftrightarrow} X(\omega)$, then

$$
e^{j \omega_{0} t} x(t) \stackrel{\mathrm{CTFT}}{\longleftrightarrow} X\left(\omega-\omega_{0}\right)
$$

where $\omega_{0}$ is an arbitrary real constant.

- This is known as the modulation (or frequency-domain shifting) property of the Fourier transform.
- If $x(t) \stackrel{\text { CTFT }}{\longleftrightarrow} X(\omega)$, then

$$
x(a t) \stackrel{\text { cTFT }}{\longleftrightarrow} \frac{1}{|a|} X\left(\frac{\omega}{a}\right),
$$

where $a$ is an arbitrary nonzero real constant.

- This is known as the dilation (or time/frequency-domain scaling) property of the Fourier transform.


## Conjugation

- If $x(t) \stackrel{\text { CTFT }}{\longleftrightarrow} X(\omega)$, then

$$
x^{*}(t) \stackrel{\text { CTFT }}{\longleftrightarrow} X^{*}(-\omega) .
$$

- This is known as the conjugation property of the Fourier transform.


## Duality

- If $x(t) \stackrel{\text { CTFT }}{\longleftrightarrow} X(\omega)$, then

$$
X(t) \stackrel{\text { CTFT }}{\longleftrightarrow} 2 \pi x(-\omega)
$$

- This is known as the duality property of the Fourier transform.
- This property follows from the high degree of symmetry in the forward and inverse Fourier transform equations, which are respectively given by

$$
X(\lambda)=\int_{-\infty}^{\infty} x(\theta) e^{-j \theta \lambda} d \theta \quad \text { and } \quad x(\lambda)=\frac{1}{2 \pi} \int_{-\infty}^{\infty} X(\theta) e^{j \theta \lambda} d \theta
$$

- That is, the forward and inverse Fourier transform equations are identical except for a factor of $2 \pi$ and different sign in the parameter for the exponential function.
- Although the relationship $x(t) \stackrel{\text { CTFT }}{\longleftrightarrow} X(\omega)$ only directly provides us with the Fourier transform of $x(t)$, the duality property allows us to indirectly infer the Fourier transform of $X(t)$. Consequently, the duality property can be used to effectively double the number of Fourier transform pairs that we know.
- If $x_{1}(t) \stackrel{\text { cTFT }}{\longleftrightarrow} X_{1}(\omega)$ and $x_{2}(t) \stackrel{\text { CTFT }}{\longleftrightarrow} X_{2}(\omega)$, then

$$
x_{1} * x_{2}(t) \stackrel{\text { cTFT }}{\longleftrightarrow} X_{1}(\omega) X_{2}(\omega) .
$$

- This is known as the convolution (or time-domain convolution) property of the Fourier transform.
- In other words, a convolution in the time domain becomes a multiplication in the frequency domain.
- This suggests that the Fourier transform can be used to avoid having to deal with convolution operations.

■ If $x_{1}(t) \stackrel{\text { CTFT }}{\longleftrightarrow} X_{1}(\omega)$ and $x_{2}(t) \stackrel{\text { CTFT }}{\longleftrightarrow} X_{2}(\omega)$, then

$$
x_{1}(t) x_{2}(t) \stackrel{\text { СтनT }}{\longleftrightarrow} \frac{1}{2 \pi} X_{1} * X_{2}(\omega)=\frac{1}{2 \pi} \int_{-\infty}^{\infty} X_{1}(\theta) X_{2}(\omega-\theta) d \theta
$$

■ This is known as the (time-domain) multiplication (or frequency-domain convolution) property of the Fourier transform.

- In other words, multiplication in the time domain becomes convolution in the frequency domain (up to a scale factor of $2 \pi$ ).
- Do not forget the factor of $\frac{1}{2 \pi}$ in the above formula!
- This property of the Fourier transform is often tedious to apply (in the forward direction) as it turns a multiplication into a convolution.
- If $x(t) \stackrel{\text { CTFT }}{\longleftrightarrow} X(\omega)$, then

$$
\frac{d x(t)}{d t} \stackrel{\text { cтғT }}{\longleftrightarrow} j \omega X(\omega) .
$$

- This is known as the (time-domain) differentiation property of the Fourier transform.
- Differentiation in the time domain becomes multiplication by $j \omega$ in the frequency domain.
- Of course, by repeated application of the above property, we have that $\left(\frac{d}{d t}\right)^{n} x(t) \stackrel{\text { CTFT }}{\longleftrightarrow}(j \omega)^{n} X(\omega)$.
- The above suggests that the Fourier transform might be a useful tool when working with differential (or integro-differential) equations.
- If $x(t) \stackrel{\text { CTFT }}{\longleftrightarrow} X(\omega)$, then

$$
t x(t) \stackrel{\text { cTFT }}{\longleftrightarrow} j \frac{d}{d \omega} X(\omega) .
$$

- This is known as the frequency-domain differentiation property of the Fourier transform.


## Time-Domain Integration

- If $x(t) \stackrel{\text { CTFT }}{\longleftrightarrow} X(\omega)$, then

$$
\int_{-\infty}^{t} x(\tau) d \tau \stackrel{\text { стет }}{\longleftrightarrow} \frac{1}{j \omega} X(\omega)+\pi X(0) \delta(\omega) .
$$

- This is known as the (time-domain) integration property of the Fourier transform.
- Whereas differentiation in the time domain corresponds to multiplication by $j \omega$ in the frequency domain, integration in the time domain is associated with division by $j \omega$ in the frequency domain.
- Since integration in the time domain becomes division by $j \omega$ in the frequency domain, integration can be easier to handle in the frequency domain.
- The above property suggests that the Fourier transform might be a useful tool when working with integral (or integro-differential) equations.
- Recall that the energy of a function $x$ is given by $\int_{-\infty}^{\infty}|x(t)|^{2} d t$.
- If $x(t) \stackrel{\text { CTFT }}{\longleftrightarrow} X(\omega)$, then

$$
\int_{-\infty}^{\infty}|x(t)|^{2} d t=\frac{1}{2 \pi} \int_{-\infty}^{\infty}|X(\omega)|^{2} d \omega
$$

(i.e., the energy of $x$ and energy of $X$ are equal up to a factor of $2 \pi$ ).

- This relationship is known as Parseval's relation.
- Since energy is often a quantity of great significance in engineering applications, it is extremely helpful to know that the Fourier transform preserves energy (up to a scale factor).
- For a function $x$ with Fourier transform $X$, the following assertions hold:

$$
\begin{gathered}
x \text { is even } \Leftrightarrow X \text { is even; and } \\
x \text { is odd } \Leftrightarrow X \text { is odd. }
\end{gathered}
$$

- In other words, the forward and inverse Fourier transforms preserve even/odd symmetry.
- A function $x$ is real if and only if its Fourier transform $X$ satisfies

$$
X(\omega)=X^{*}(-\omega) \text { for all } \omega
$$

(i.e., $X$ is conjugate symmetric).

- Thus, for a real-valued function, the portion of the graph of $X(\omega)$ for $\omega<0$ is completely redundant, as it is determined by symmetry.
- From properties of complex numbers, one can show that $X(\omega)=X^{*}(-\omega)$ is equivalent to

$$
|X(\omega)|=|X(-\omega)| \quad \text { and } \quad \arg X(\omega)=-\arg X(-\omega)
$$

(i.e., $|X(\omega)|$ is even and $\arg X(\omega)$ is odd).

- Note that $x$ being real does not necessarily imply that $X$ is real.


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Section 6.4

## Fourier Transform of Periodic Functions

- The Fourier transform can be generalized to also handle periodic functions.
- Consider a periodic function $x$ with period $T$ and frequency $\omega_{0}=\frac{2 \pi}{T}$.
- Define the function $x_{T}$ as

$$
x_{T}(t)= \begin{cases}x(t) & -\frac{T}{2} \leq t<\frac{T}{2} \\ 0 & \text { otherwise }\end{cases}
$$

(i.e., $x_{T}(t)$ is equal to $x(t)$ over a single period and zero elsewhere).

- Let $a$ denote the Fourier series coefficient sequence of $x$.
- Let $X$ and $X_{T}$ denote the Fourier transforms of $x$ and $x_{T}$, respectively.
- The following relationships can be shown to hold:

$$
\begin{gathered}
X(\omega)=\sum_{k=-\infty}^{\infty} \omega_{0} X_{T}\left(k \omega_{0}\right) \delta\left(\omega-k \omega_{0}\right), \\
a_{k}=\frac{1}{T} X_{T}\left(k \omega_{0}\right), \quad \text { and } \quad X(\omega)=\sum_{k=-\infty}^{\infty} 2 \pi a_{k} \delta\left(\omega-k \omega_{0}\right) .
\end{gathered}
$$

- The Fourier transform $X$ of a periodic function is a series of impulses that occur at integer multiples of the fundamental frequency $\omega_{0}$ (i.e., $\left.X(\omega)=\sum_{k=-\infty}^{\infty} 2 \pi a_{k} \delta\left(\omega-k \omega_{0}\right)\right)$.
- Due to the preceding fact, the Fourier transform of a periodic function can only be nonzero at integer multiples of the fundamental frequency.
- The Fourier series coefficient sequence $a$ is produced by sampling $X_{T}$ at integer multiples of the fundamental frequency $\omega_{0}$ and scaling the resulting sequence by $\frac{1}{T}$ (i.e., $a_{k}=\frac{1}{T} X_{T}\left(k \omega_{0}\right)$ ).


## Section 6.5

## Fourier Transform and Frequency Spectra of Functions

- Like Fourier series, the Fourier transform also provides us with a frequency-domain perspective on functions.
- That is, instead of viewing a function as having information distributed with respect to time (i.e., a function whose domain is time), we view a function as having information distributed with respect to frequency (i.e., a function whose domain is frequency).
- The Fourier transform of a function $x$ provides a means to quantify how much information $x$ has at different frequencies.
- The distribution of information in a function over different frequencies is referred to as the frequency spectrum of the function.


## Fourier Transform and Frequency Spectra

- To gain further insight into the role played by the Fourier transform $X$ in the context of the frequency spectrum of $x$, it is helpful to write the Fourier transform representation of $x$ with $X(\omega)$ expressed in polar form as follows:

$$
x(t)=\frac{1}{2 \pi} \int_{-\infty}^{\infty} X(\omega) e^{j \omega t} d \omega=\frac{1}{2 \pi} \int_{-\infty}^{\infty}|X(\omega)| e^{j[\omega t+\arg X(\omega)]} d \omega .
$$

- In effect, the quantity $|X(\omega)|$ is a weight that determines how much the complex sinusoid at frequency $\omega$ contributes to the integration result $x$.
- The quantity $\arg X(\omega)$ determines how the complex sinusoid at frequency $\omega$ is shifted related to complex sinusoids at other frequencies.
- Perhaps, this can be more easily seen if we express the above integral as the limit of a sum, derived from an approximation of the integral using the areas of rectangles, as shown on the next slide. [Recall that $\left.\int_{-\infty}^{\infty} f(x) d x=\lim _{\Delta x \rightarrow 0} \sum_{k=-\infty}^{\infty} \Delta x f(k \Delta x).\right]$


## Fourier Transform and Frequency Spectra (Continued 1)

■ Expressing the integral (from the previous slide) as the limit of a sum, we obtain

$$
x(t)=\lim _{\Delta \omega \rightarrow 0} \frac{1}{2 \pi} \sum_{k=-\infty}^{\infty} \Delta \omega|X(\omega)| e^{j[\omega t+\arg X(\omega)]}
$$

where $\omega=k \Delta \omega$.

- In the above equation, the $k$ th term in the summation corresponds to a complex sinusoid with fundamental frequency $\omega=k \Delta \omega$ that has had its amplitude scaled by a factor of $|X(\omega)|$ and has been time shifted by an amount that depends on $\arg X(\omega)$.
- For a given $\omega=k \Delta \omega$ (which is associated with the $k$ th term in the summation), the larger $|X(\omega)|$ is, the larger the amplitude of its corresponding complex sinusoid $e^{j \omega t}$ will be, and therefore the larger the contribution the $k$ th term will make to the overall summation.
- In this way, we can use $|X(\omega)|$ as a measure of how much information a function $x$ has at the frequency $\omega$.


## Fourier Transform and Frequency Spectra (Continued 2)

- The Fourier transform $X$ of the function $x$ is referred to as the frequency spectrum of $x$.
- The magnitude $|X(\omega)|$ of the Fourier transform $X$ is referred to as the magnitude spectrum of $x$.
- The argument $\arg X(\omega)$ of the Fourier transform $X$ is referred to as the phase spectrum of $x$.
- Since the Fourier transform is a function of a real variable, a function can potentially have information at any real frequency.
- Since the Fourier transform $X$ of a periodic function $x$ with fundamental frequency $\omega_{0}$ and the Fourier series coefficient sequence $a$ is given by $X(\omega)=\sum_{k=-\infty}^{\infty} 2 \pi a_{k} \delta\left(\omega-k \omega_{0}\right)$, the Fourier transform and Fourier series give consistent results for the frequency spectrum of a periodic function.
- Since the frequency spectrum is complex (in the general case), it is usually represented using two plots, one showing the magnitude spectrum and one showing the phase spectrum.


## Frequency Spectra of Real Functions

- Recall that, for a real function $x$, the Fourier transform $X$ of $x$ satisfies

$$
X(\omega)=X^{*}(-\omega)
$$

(i.e., $X$ is conjugate symmetric), which is equivalent to

$$
|X(\omega)|=|X(-\omega)| \quad \text { and } \quad \arg X(\omega)=-\arg X(-\omega) .
$$

- Since $|X(\omega)|=|X(-\omega)|$, the magnitude spectrum of a real function is always even.
- Similarly, since $\arg X(\omega)=-\arg X(-\omega)$, the phase spectrum of a real function is always odd.
- Due to the symmetry in the frequency spectra of real functions, we typically ignore negative frequencies when dealing with such functions.
- In the case of functions that are complex but not real, frequency spectra do not possess the above symmetry, and negative frequencies become important.


## Bandwidth

- A function with the Fourier transform $X$ is said to be bandllimited if, for some (finite) nonnegative real constant $B$, the following condition holds:

$$
X(\omega)=0 \text { for all } \omega \text { satisfying }|\omega|>B
$$

- The bandwidth $B$ of a function with the Fourier transform $X$ is defined as $B=\omega_{1}-\omega_{0}$, where $X(\omega)=0$ for all $\omega \notin\left[\omega_{0}, \omega_{1}\right]$.
- In the case of real-valued functions, however, this definition of bandwidth is usually amended to consider only nonnegative frequencies.
- The real-valued function $x_{1}$ and complex-valued function $x_{2}$ with the respective Fourier transforms $X_{1}$ and $X_{2}$ shown below each have bandwidth $B$ (where only nonnegative frequencies are considered in the case of $x_{1}$ ).


- One can show that a function cannot be both time limited and bandlimited.


## Energy-Density Spectra

- By Parseval's relation, the energy $E$ in a function $x$ with Fourier transform $X$ is given by

$$
E=\frac{1}{2 \pi} \int_{-\infty}^{\infty} E_{x}(\omega) d \omega
$$

where

$$
E_{x}(\omega)=|X(\omega)|^{2}
$$

- We refer to $E_{x}$ as the energy-density spectrum of the function $x$.
- The function $E_{x}$ indicates how the energy in $x$ is distributed with respect to frequency.
- For example, the energy contributed by frequencies in the range $\left[\omega_{1}, \omega_{2}\right.$ ] is given by

$$
\frac{1}{2 \pi} \int_{\omega_{1}}^{\omega_{2}} E_{x}(\omega) d \omega
$$

## Section 6.6

## Fourier Transform and LTI Systems

- Consider a LTI system with input $x$, output $y$, and impulse response $h$, and let $X, Y$, and $H$ denote the Fourier transforms of $x, y$, and $h$, respectively.
- Since $y(t)=x * h(t)$, we have that

$$
Y(\omega)=X(\omega) H(\omega)
$$

- The function $H$ is called the frequency response of the system.
- A LTI system is completely characterized by its frequency response $H$.
- The above equation provides an alternative way of viewing the behavior of a LTI system. That is, we can view the system as operating in the frequency domain on the Fourier transforms of the input and output functions.
- The frequency spectrum of the output is the product of the frequency spectrum of the input and the frequency response of the system.
- In the general case, the frequency response $H$ is a complex-valued function.
■ Often, we represent $H(\omega)$ in terms of its magnitude $|H(\omega)|$ and argument $\arg H(\omega)$.
- The quantity $|H(\omega)|$ is called the magnitude response of the system.
- The quantity $\arg H(\omega)$ is called the phase response of the system.
- Since $Y(\omega)=X(\omega) H(\omega)$, we trivially have that

$$
|Y(\omega)|=|X(\omega)||H(\omega)| \quad \text { and } \quad \arg Y(\omega)=\arg X(\omega)+\arg H(\omega)
$$

- The magnitude spectrum of the output equals the magnitude spectrum of the input times the magnitude response of the system.
- The phase spectrum of the output equals the phase spectrum of the input plus the phase response of the system.
- Since the frequency response $H$ is simply the frequency spectrum of the impulse response $h$, if $h$ is real, then

$$
|H(\omega)|=|H(-\omega)| \quad \text { and } \quad \arg H(\omega)=-\arg H(-\omega)
$$

(i.e., the magnitude response $|H(\omega)|$ is even and the phase response $\arg H(\omega)$ is odd $)$.

## Unwrapped Phase

- For many types of analysis, restricting the range of a phase function to an interval of length $2 \pi$ (such as $(-\pi, \pi]$ ), often unnecessarily introduces discontinuities into the function.
- This motivates the notion of unwrapped phase.
- The unwrapped phase is simply the phase defined in such a way so as not to restrict the phase to an interval of length $2 \pi$ and to keep the phase function continuous to the greatest extent possible.
- For example, the function $H(\omega)=e^{j \pi \omega}$ has the unwrapped phase $\Theta(\omega)=\pi \omega$.


Phase


Unwrapped Phase

## Interpretation of Magnitude and Phase Response

- Recall that a LTI system $\mathcal{H}$ with frequency response $H$ is such that

$$
\mathcal{H}\left\{e^{j \omega t}\right\}(t)=H(\omega) e^{j \omega t}
$$

- Expressing $H(\omega)$ in polar form, we have

$$
\begin{aligned}
\mathcal{H}\left\{e^{j \omega t}\right\}(t) & =|H(\omega)| e^{j \arg H(\omega)} e^{j \omega t} \\
& =|H(\omega)| e^{j[\omega t+\arg H(\omega)]} \\
& =|H(\omega)| e^{j \omega(t+\arg [H(\omega)] / \omega)}
\end{aligned}
$$

- Thus, the response of the system to the function $e^{j \omega t}$ is produced by applying two transformations to this function:
- (amplitude) scaling by $|H(\omega)|$; and
- translating by $-\frac{\arg H(\omega)}{\omega}$.
- Therefore, the magnitude response determines how different complex sinusoids are scaled (in amplitude) by the system.
- Similarly, the phase response determines how different complex sinusoids are translated (i.e., delayed/advanced) by the system.
- Recall that a LTI system $\mathcal{H}$ with frequency response $H$ is such that

$$
\mathcal{H}\left\{e^{j \omega t}\right\}(t)=|H(\omega)| e^{j \omega(t+\arg [H(\omega)] / \omega)} .
$$

- If $|H(\omega)|$ is a constant (for all $\omega$ ), every complex sinusoid is scaled by the same amount when passing through the system.
- A system for which $|H(\omega)|=1$ (for all $\omega$ ) is said to be allpass.
- In the case of an allpass system, the magnitude spectra of the system's input and output are identical.
- If $|H(\omega)|$ is not a constant, different complex sinusoids are scaled by different amounts, resulting in what is known as magnitude distortion.

■ Recall that a LTI system $\mathcal{H}$ with frequency response $H$ is such that

$$
\mathcal{H}\left\{e^{j \omega t}\right\}(t)=|H(\omega)| e^{j \omega(t+\arg [H(\omega)] / \omega)}
$$

- The preceding equation can be rewritten as

$$
\mathcal{H}\left\{e^{j \omega t}\right\}(t)=|H(\omega)| e^{j \omega\left[t-\tau_{\mathrm{p}}(\omega)\right]} \quad \text { where } \quad \tau_{\mathrm{p}}(\omega)=-\frac{\arg H(\omega)}{\omega}
$$

- The function $\tau_{\mathrm{p}}$ is known as the phase delay of the system.
- If $\tau_{\mathrm{p}}(\omega)=t_{\mathrm{d}}$ (where $t_{\mathrm{d}}$ is a constant), the system shifts all complex sinusoids by the same amount $t_{\mathrm{d}}$.
- Since $\tau_{\mathrm{p}}(\omega)=t_{\mathrm{d}}$ is equivalent to the (unwrapped) phase response being of the form $\arg H(\omega)=-t_{\mathrm{d}} \omega$ (which is a linear function with a zero constant term), a system with a constant phase delay is said to have linear phase.
- In the case that $\tau_{\mathrm{p}}(\omega)=0$, the system is said to have zero phase.
- If $\tau_{\mathrm{p}}(\omega)$ is not a constant, different complex sinusoids are shifted by different amounts, resulting in what is known as phase distortion.


## Distortionless Transmission

- Consider a LTI system $\mathcal{H}$ with input $x$ and output $y$ given by

$$
y(t)=x\left(t-t_{0}\right)
$$

where $t_{0}$ is a real constant.

- That is, the output of the system is simply the input delayed by $t_{0}$.
- This type of behavior is the ideal for which we strive in real-world communication systems (i.e., the received signal $y$ equals a delayed version of the transmitted signal $x$ ).
- Taking the Fourier transform of the preceding equation, we have

$$
Y(\omega)=e^{-j \omega t_{0}} X(\omega)
$$

- Thus, the system has the frequency response $H$ given by

$$
H(\omega)=e^{-j \omega t_{0}}
$$

- Since the phase delay of the system is $\tau_{\mathrm{p}}(\omega)=-\left(\frac{-\omega t_{0}}{\omega}\right)=t_{0}$, the phase delay is constant and the system has linear phase.


## Magnitude and Phase Distortion in Audio

- The relative importance of the magnitude spectrum and phase spectrum is highly dependent on the particular application of interest.
- Consider the case of the human auditory system (i.e., human hearing).
- The human auditory system tends to be quite sensitive to changes in the magnitude spectrum of an audio signal.
- That is, a significant change in the magnitude spectrum of an audio signal is very likely to lead to a noticeable difference in the perceived sound.
- On the other hand, the human auditory system tends to be much less sensitive to changes in the phase spectrum of an audio signal.
- In other words, changes to the phase spectrum of an audio signal are often only barely perceptible or not perceptible at all.
- For the above reasons, in applications involving the human auditory system, magnitude distortion often tends to be more of a concern than phase distortion.


## Magnitude and Phase Distortion in Images

- Consider the case of the human visual system.
- The human visual system tends to be quite sensitive to changes in the phase spectrum of an image.
- That is, a significant change in the phase spectrum of an image is likely to lead to a very substantial difference in how the image is perceived.
- The phase spectrum of an image tends to capture information about the location of the edges in the image, and edges are play a crucial role in how humans perceive images.
- On the other hand, the human visual system tends to be somewhat less sensitive to changes in the magnitude spectrum of an image.
- For the above reasons, phase distortion is usually deemed highly undesirable in systems that process images, when the image data is to be consumed by humans.


## Example: Magnitude and Phase Distortion in Images (1)



Image A


Magnitude Spectrum from Image $B$ and Phase Spectrum from Image A


Image B


Magnitude Spectrum from Image A and
Phase Spectrum from Image B

## Example: Magnitude and Phase Distortion in Images (2)



Image A


Magnitude Spectrum from Image $B$ and Phase Spectrum from Image A


Image B (White Noise)


Magnitude Spectrum from Image A and
Phase Spectrum from Image B

- Consider a LTI system with input $x$, output $y$, and impulse response $h$, and let $X, Y$, and $H$ denote the Fourier transforms of $x, y$, and $h$, respectively.
- Often, it is convenient to represent such a system in block diagram form in the frequency domain as shown below.


■ Since a LTI system is completely characterized by its frequency response, we typically label the system with this quantity.

- The series interconnection of the LTI systems with frequency responses $H_{1}$ and $H_{2}$ is the LTI system with frequency response $H_{1} H_{2}$. That is, we have the equivalence shown below.

- The parallel interconnection of the LTI systems with frequency responses $H_{1}$ and $H_{2}$ is the LTI system with the frequency response $H_{1}+H_{2}$. That is, we have the equivalence shown below.



## LTI Systems and Differential Equations

■ Many LTI systems of practical interest can be represented using an Nth-order linear differential equation with constant coefficients.

- Consider a system with input $x$ and output $y$ that is characterized by an equation of the form

$$
\sum_{k=0}^{N} b_{k}\left(\frac{d}{d t}\right)^{k} y(t)=\sum_{k=0}^{M} a_{k}\left(\frac{d}{d t}\right)^{k} x(t)
$$

where the $a_{k}$ and $b_{k}$ are complex constants and $M \leq N$.

- Let $h$ denote the impulse response of the system, and let $X, Y$, and $H$ denote the Fourier transforms of $x, y$, and $h$, respectively.
- One can show that $H$ is given by

$$
H(\omega)=\frac{Y(\omega)}{X(\omega)}=\frac{\sum_{k=0}^{M} a_{k} j^{k} \omega^{k}}{\sum_{k=0}^{N} b_{k} j^{k} \omega^{k}}
$$

- Observe that, for a system of the form considered above, the frequency response is a rational function.


## Section 6.7

## Application: Filtering

- In many applications, we want to modify the spectrum of a function by either amplifying or attenuating certain frequency components.
- This process of modifying the frequency spectrum of a function is called filtering.
- A system that performs a filtering operation is called a filter.
- Many types of filters exist.
- Frequency selective filters pass some frequencies with little or no distortion, while significantly attenuating other frequencies.
- Several basic types of frequency-selective filters include: lowpass, highpass, and bandpass.


## Ideal Lowpass Filter

- An ideal lowpass filter eliminates all frequency components with a frequency whose magnitude is greater than some cutoff frequency, while leaving the remaining frequency components unaffected.
- Such a filter has a frequency response $H$ of the form

$$
H(\omega)= \begin{cases}1 & |\omega| \leq \omega_{c} \\ 0 & \text { otherwise }\end{cases}
$$

where $\omega_{c}$ is the cutoff frequency.

- A plot of this frequency response is given below.



## Ideal Highpass Filter

- An ideal highpass filter eliminates all frequency components with a frequency whose magnitude is less than some cutoff frequency, while leaving the remaining frequency components unaffected.
- Such a filter has a frequency response $H$ of the form

$$
H(\omega)= \begin{cases}1 & |\omega| \geq \omega_{c} \\ 0 & \text { otherwise }\end{cases}
$$

where $\omega_{c}$ is the cutoff frequency.

- A plot of this frequency response is given below.



## Ideal Bandpass Filter

- An ideal bandpass filter eliminates all frequency components with a frequency whose magnitude does not lie in a particular range, while leaving the remaining frequency components unaffected.
- Such a filter has a frequency response $H$ of the form

$$
H(\omega)= \begin{cases}1 & \omega_{c 1} \leq|\omega| \leq \omega_{c 2} \\ 0 & \text { otherwise }\end{cases}
$$

where the limits of the passband are $\omega_{c 1}$ and $\omega_{c 2}$.

- A plot of this frequency response is given below.



## Section 6.8

## Application: Equalization

- Often, we find ourselves faced with a situation where we have a system with a particular frequency response that is undesirable for the application at hand.

■ As a result, we would like to change the frequency response of the system to be something more desirable.

- This process of modifying the frequency response in this way is referred to as equalization. [Essentially, equalization is just a filtering operation.]
- Equalization is used in many applications.
- In real-world communication systems, equalization is used to eliminate or minimize the distortion introduced when a signal is sent over a (nonideal) communication channel.
- In audio applications, equalization can be employed to emphasize or de-emphasize certain ranges of frequencies. For example, equalization can be used to boost the bass (i.e., emphasize the low frequencies) in the audio output of a stereo.

Equalization (Continued)


Original System


New System with Equalization

■ Let $H_{\text {orig }}$ denote the frequency response of original system (i.e., without equalization).

- Let $H_{\mathrm{d}}$ denote the desired frequency response.
- Let $H_{\text {eq }}$ denote the frequency response of the equalizer.
- The new system with equalization has frequency response

$$
H_{\mathrm{new}}(\omega)=H_{\mathrm{eq}}(\omega) H_{\text {orig }}(\omega)
$$

- By choosing $H_{\text {eq }}(\omega)=H_{\mathrm{d}}(\omega) / H_{\text {orig }}(\omega)$, the new system with equalization will have the frequency response

$$
H_{\text {new }}(\omega)=\left[H_{\mathrm{d}}(\omega) / H_{\text {orig }}(\omega)\right] H_{\text {orig }}(\omega)=H_{\mathrm{d}}(\omega) .
$$

- In effect, by using an equalizer, we can obtain a new system with the frequency response that we desire.


## Section 6.9

## Application: Circuit Analysis

- An electronic circuit is a network of one or more interconnected circuit elements.
- The three most basic types of circuit elements are:

11 resistors;
2 inductors; and
${ }^{3}$ capacitors.

- Two fundamental quantities of interest in electronic circuits are current and voltage.
- Current is the rate at which electric charge flows through some part of a circuit, such as a circuit element, and is measured in units of amperes (A).
- Voltage is the difference in electric potential between two points in a circuit, such as across a circuit element, and is measured in units of volts (V).
- Voltage is essentially a force that makes electric charge (or current) flow.
- A resistor is a circuit element that opposes the flow of current.
- A resistor is characterized by an equation of the form

$$
\left.v(t)=\operatorname{Ri}(t) \quad \text { (or equivalently, } i(t)=\frac{1}{R} v(t)\right),
$$

where $R$ is a nonnegative real constant, and $v$ and $i$ respectively denote the voltage across and current through the resistor as a function of time.

- As a matter of terminology, the quantity $R$ is known as the resistance of the resistor.
- Resistance is measured in units of ohms ( $\Omega$ ).
- In circuit diagrams, a resistor is denoted by the symbol shown below.



## Inductors

- An inductor is a circuit element that converts an electric current into a magnetic field and vice versa.
- An inductor uses the energy stored in a magnetic field in order to oppose changes in current (through the inductor).
- An inductor is characterized by an equation of the form

$$
\left.v(t)=L \frac{d}{d t} i(t) \quad \text { (or equivalently, } i(t)=\frac{1}{L} \int_{-\infty}^{t} v(\tau) d \tau\right)
$$

where $L$ is a nonnegative real constant, and $v$ and $i$ respectively denote the voltage across and current through the inductor as a function of time.

- As a matter of terminology, the quantity $L$ is known as the inductance of the inductor.
- Inductance is measured in units of henrys $(\mathrm{H})$.
- In circuit diagrams, an inductor is denoted by the symbol shown below.



## Capacitors

- A capacitor is a circuit element that stores electric charge.
- A capacitor uses the energy stored in an electric field in order to oppose changes in voltage (across the capacitor).
- A capacitor is characterized by an equation of the form

$$
\left.v(t)=\frac{1}{C} \int_{-\infty}^{t} i(\tau) d \tau \quad \text { (or equivalently, } i(t)=C \frac{d}{d t} v(t)\right)
$$

where $C$ is a nonnegative real constant, and $v$ and $i$ respectively denote the voltage across and current through the capacitor as a function of time.

- As a matter of terminology, the quantity $C$ is known as the capacitance of the capacitor.
- Capacitance is measured in units of farads ( F ).
- In circuit diagrams, a capacitor is denoted by the symbol shown below.



## Circuit Analysis with the Fourier Transform

- The Fourier transform is a very useful tool for circuit analysis.
- The utility of the Fourier transform is partly due to the fact that the differential/integral equations that describe inductors and capacitors are much simpler to express in the Fourier domain than in the time domain.
- Let $v$ and $i$ denote the voltage across and current through a circuit element, and let $V$ and $I$ denote the Fourier transforms of $v$ and $i$, respectively.
- In the frequency domain, the equations characterizing a resistor, an inductor, and a capacitor respectively become:

$$
\begin{gathered}
V(\omega)=R I(\omega) \quad \text { (or equivalently, } I(\omega)=\frac{1}{R} V(\omega) \text { ); } \\
V(\omega)=j \omega L I(\omega) \quad \text { (or equivalently, } I(\omega)=\frac{1}{j \omega L} V(\omega) \text {; and } \\
\left.V(\omega)=\frac{1}{j \omega C} I(\omega) \quad \text { (or equivalently, } I(\omega)=j \omega C V(\omega)\right) .
\end{gathered}
$$

- Note the absence of differentiation and integration in the above equations for an inductor and a capacitor.


## Section 6.10

## Application: Amplitude Modulation (AM)

## Motivation for Amplitude Modulation (AM)

- In communication systems, we often need to transmit a signal using a frequency range that is different from that of the original signal.
- For example, voice/audio signals typically have information in the range of 0 to 22 kHz .
- Often, it is not practical to transmit such a signal using its original frequency range.
- Two potential problems with such an approach are:

I interference; and
[ constraints on antenna length.

- Since many signals are broadcast over the airwaves, we need to ensure that no two transmitters use the same frequency bands in order to avoid interference.
- Also, in the case of transmission via electromagnetic waves (e.g., radio waves), the length of antenna required becomes impractically large for the transmission of relatively low frequency signals.
- For the preceding reasons, we often need to change the frequency range associated with a signal before transmission.


Transmitter


Receiver

- The transmitter is characterized by

$$
y(t)=e^{j \omega_{c} t} x(t) \quad \Longleftrightarrow \quad Y(\omega)=X\left(\omega-\omega_{c}\right) .
$$

- The receiver is characterized by

$$
\hat{x}(t)=e^{-j \omega_{c} t} y(t) \quad \Longleftrightarrow \quad \hat{X}(\omega)=Y\left(\omega+\omega_{c}\right) .
$$

■ Clearly, $\hat{x}(t)=e^{j \omega_{c} t} e^{-j \omega_{c} t} x(t)=x(t)$.

## Trivial Amplitude Modulation (AM) System: Example



Transmitter Input


Transmitter Output


Receiver Output

## Double-Sideband Suppressed-Carrier (DSB-SC) AM



Transmitter


Receiver

■ Let $X=\mathcal{F} x, Y=\mathcal{F} y$, and $\hat{X}=\mathcal{F} \hat{x}$.

- Suppose that $X(\omega)=0$ for all $\omega \notin\left[-\omega_{b}, \omega_{b}\right]$.
- The transmitter is characterized by

$$
Y(\omega)=\frac{1}{2}\left[X\left(\omega+\omega_{c}\right)+X\left(\omega-\omega_{c}\right)\right] .
$$

- The receiver is characterized by

$$
\hat{X}(\omega)=\left[Y\left(\omega+\omega_{c}\right)+Y\left(\omega-\omega_{c}\right)\right] \operatorname{rect}\left(\frac{\omega}{2 \omega_{c_{0}}}\right) .
$$

- If $\omega_{b}<\omega_{c_{0}}<2 \omega_{c}-\omega_{b}$, we have $\hat{X}(\omega)=X(\omega)$ (implying $\hat{x}(t)=x(t)$ ).


$$
\begin{gathered}
y(t)=\cos \left(\omega_{c} t\right) x(t) \\
X=\mathcal{F} x, \quad Y=\mathcal{F} y \\
Y(\omega)=\mathcal{F}\left\{\cos \left(\omega_{c} t\right) x(t)\right\}(\omega) \\
=\mathcal{F}\left\{\frac{1}{2}\left(e^{j \omega_{c} t}+e^{-j \omega_{c} t}\right) x(t)\right\}(\omega) \\
=\frac{1}{2}\left[\mathcal{F}\left\{e^{j \omega_{c} t} x(t)\right\}(\omega)+\mathcal{F}\left\{e^{-j \omega_{c} t} x(t)\right\}(\omega)\right] \\
=\frac{1}{2}\left[X\left(\omega-\omega_{c}\right)+X\left(\omega+\omega_{c}\right)\right]
\end{gathered}
$$



$$
\begin{aligned}
& v(t)=\cos \left(\omega_{c} t\right) y(t), \quad h(t)=\frac{2 \omega_{c} 0}{\pi} \operatorname{sinc}\left(\omega_{c} t\right), \quad \hat{x}(t)=v * h(t) \\
& Y=\mathcal{F} y, \quad V=\mathcal{F} v, \quad H=\mathcal{F} h, \quad \hat{X}=\mathcal{F} \hat{x} \\
& V(\omega)=F \mathcal{F}\left\{\cos \left(\omega_{c} t\right) y(t)\right\}(\omega) \\
& =\mathcal{F}\left\{\frac{1}{2}\left(e^{j \omega_{c} t}+e^{-j \omega_{c} t}\right) y(t)\right\}(\omega) \\
& =\frac{1}{2}\left[\mathcal{F}\left\{e^{j \omega_{c} t} y(t)\right\}(\omega)+\mathcal{F}\left\{e^{-j \omega_{c} t} y(t)\right\}(\omega)\right] \\
& =\frac{1}{2}\left[Y\left(\omega-\omega_{c}\right)+Y\left(\omega+\omega_{c}\right)\right] \\
& H(\omega)=\mathcal{F}\left\{\frac{2 \omega_{c 0}}{\pi} \operatorname{sinc}\left(\omega_{c 0} t\right)\right\}(\omega) \\
& =2 \operatorname{rect}\left(\frac{\omega}{2 \omega_{c 0}}\right) \\
& \hat{X}(\omega)=H(\omega) V(\omega)
\end{aligned}
$$

## DSB-SC AM: Complete System



$$
\begin{gathered}
Y(\omega)=\frac{1}{2}\left[X\left(\omega-\omega_{c}\right)+X\left(\omega+\omega_{c}\right)\right] \\
V(\omega)=\frac{1}{2}\left[Y\left(\omega-\omega_{c}\right)+Y\left(\omega+\omega_{c}\right)\right] \\
=\frac{1}{2}\left[\frac{1}{2}\left[X\left(\left[\omega-\omega_{c}\right]-\omega_{c}\right)+X\left(\left[\omega-\omega_{c}\right]+\omega_{c}\right)\right]+\right. \\
\left.\frac{1}{2}\left[X\left(\left[\omega+\omega_{c}\right]-\omega_{c}\right)+X\left(\left[\omega+\omega_{c}\right]+\omega_{c}\right)\right]\right] \\
=\frac{1}{2} X(\omega)+\frac{1}{4} X\left(\omega-2 \omega_{c}\right)+\frac{1}{4} X\left(\omega+2 \omega_{c}\right)
\end{gathered}
$$

$$
\begin{aligned}
\hat{X}(\omega) & =H(\omega) V(\omega) \\
& =H(\omega)\left[\frac{1}{2} X(\omega)+\frac{1}{4} X\left(\omega-2 \omega_{c}\right)+\frac{1}{4} X\left(\omega+2 \omega_{c}\right)\right] \\
& =\frac{1}{2} H(\omega) X(\omega)+\frac{1}{4} H(\omega) X\left(\omega-2 \omega_{c}\right)+\frac{1}{4} H(\omega) X\left(\omega+2 \omega_{c}\right) \\
& =\frac{1}{2}[2 X(\omega)]+\frac{1}{4}(0)+\frac{1}{4}(0) \\
& =X(\omega)
\end{aligned}
$$

## DSB-SC AM: Example



Transmitter Input





Transmitter


Receiver

- The basic analysis of the SSB-SC AM system is similar to the DSB-SC AM system.
- SSB-SC AM requires half as much bandwidth for the transmitted signal as DSB-SC AM.


## SSB-SC AM: Example










## Section 6.11

## Application: Sampling and Interpolation

- Often, we want to be able to transform a continuous-time signal (i.e., a function) into a discrete-time signal (i.e., a sequence) and vice versa.
- This is accomplished through processes known as sampling and interpolation.
■ Sampling, which is performed by a continuous-time to discrete-time (C/D) converter shown below, transforms a function $x$ to a sequence $y$.

- Interpolation, which is performed by a discrete-time to continuous-time (D/C) converter shown below, transforms a sequence $y$ to a function $x$.

- Note that, unless very special conditions are met, the sampling process loses information (i.e., is not invertible).


## Periodic Sampling

- Although sampling can be performed in many different ways, the most commonly used scheme is periodic sampling.
- With this scheme, a sequence $y$ of samples is obtained from a function $x$ according to the relation

$$
y(n)=x(T n) \text { for all integer } n,
$$

where $T$ is a (strictly) positive real constant.

- As a matter of terminology, we refer to $T$ as the sampling period, and $\omega_{s}=\frac{2 \pi}{T}$ as the (angular) sampling frequency.
- An example of periodic sampling is shown below, where the function $x$ has been sampled with sampling period $T=10$, yielding the sequence $y$.



Sequence Produced by Sampling

## Invertibility of Sampling

- Unless constraints are placed on the functions being sampled, the sampling process is not invertible.
- In other words, in the absence of any constraints, a function cannot be uniquely determined from a sequence of its equally-spaced samples.
- Consider, for example, the functions $x_{1}$ and $x_{2}$ given by

$$
x_{1}(t)=0 \quad \text { and } \quad x_{2}(t)=\sin (2 \pi t) .
$$

- Sampling $x_{1}$ and $x_{2}$ with the sampling period $T=1$ yields the respective sequences

$$
\begin{gathered}
y_{1}(n)=x_{1}(T n)=x_{1}(n)=0 \quad \text { and } \\
y_{2}(n)=x_{2}(T n)=\sin (2 \pi n)=0 .
\end{gathered}
$$

- So, although $x_{1}$ and $x_{2}$ are distinct, $y_{1}$ and $y_{2}$ are identical.
- Given the sequence $y$ where $y=y_{1}=y_{2}$, it is impossible to determine which function was sampled to produce $y$.
- Only by imposing a carefully chosen set of constraints on the functions being sampled can we ensure that a function can be exactly recovered from only its samples.


## Model of Sampling

- An impulse train is a function of the form $v(t)=\sum_{k=-\infty}^{\infty} c_{k} \delta(t-k T)$, where $c_{k}$ and $T$ are real constants.
■ For the purposes of analysis, sampling with sampling period $T$ and frequency $\omega_{s}=\frac{2 \pi}{T}$ can be modelled as shown below.

- The sampling of a function $x$ to produce a sequence $y$ consists of the following two steps (in order):

1 Multiply the function $x$ to be sampled by a periodic impulse train $p$, yielding the impulse train $s(t)=\sum_{n=-\infty}^{\infty} x(n T) \delta(t-n T)$.
2 Convert the impulse train $s$ to a sequence $y$ by forming $y$ from the weights of successive impulses in $s$ so that $y(n)=x(n T)$.

## Model of Sampling: Various Signals




Periodic Impulse Train


Impulse-Sampled Function (Continuous-Time)


Output Sequence (Discrete-Time)


- Since sampling is not invertible and our model of sampling consists of only two steps, at least one of these two steps must not be invertible.
- Recall the two steps in our model of sampling are as follows (in order):

$$
\begin{aligned}
& \text { ■ } x \longrightarrow s(t)=x(t) p(t)=\sum_{n=-\infty}^{\infty} x(n T) \delta(t-n T) \text {; and } \\
& \text { ® } s(t)=\sum_{n=-\infty}^{\infty} x(n T) \delta(t-n T) \longrightarrow y(n)=x(n T) \text {. }
\end{aligned}
$$

- Step 1 cannot be undone (unless we somehow restrict which functions $x$ can be sampled).
- Step 2 is always invertible.
- Therefore, the fact that sampling is not invertible is entirely due to step 1.


## Model of Sampling: Characterization



- In the time domain, the impulse-sampled function $s$ is given by

$$
s(t)=x(t) p(t) \quad \text { where } \quad p(t)=\sum_{k=-\infty}^{\infty} \delta(t-k T)
$$

- In the Fourier domain, the preceding equation becomes

$$
S(\omega)=\frac{\omega_{s}}{2 \pi} \sum_{k=-\infty}^{\infty} X\left(\omega-k \omega_{s}\right) \quad\left(\text { where } \omega_{s}=\frac{2 \pi}{T}\right)
$$

- Thus, the spectrum of the impulse-sampled function $s$ is a scaled sum of an infinite number of shifted copies of the spectrum of the original function $x$.

$$
\begin{aligned}
& p(t)=\sum_{k=-\infty}^{\infty} \delta(t-k T), \quad \omega_{s}=\frac{2 \pi}{T} \\
& p(t)=\sum_{k=-\infty}^{\infty} c_{k} e^{j k \omega_{s} t} \\
& c_{k}=\frac{1}{T} \int_{-T / 2}^{T / 2} p(t) e^{-j k \omega_{s} t} d t \\
&=\frac{1}{T} \int_{-T / 2}^{T / 2} \delta(t) e^{-j k \omega_{s} t} d t \\
&=\frac{1}{T} \int_{-\infty}^{\infty} \delta(t) e^{-j k \omega_{s} t} d t \\
&=\frac{1}{T} \\
&=\frac{\omega_{s}}{2 \pi} \\
& p(t)=\frac{\omega_{s}}{2 \pi} \sum_{k=-\infty}^{\infty} e^{j k \omega_{s} t}
\end{aligned}
$$

## Sampling: Multiplication by a Periodic Impulse Train

$$
\begin{aligned}
& \text { ideal C/D converter } \\
& s(t)=p(t) x(t), \quad p(t)=\sum_{k=-\infty}^{\infty} \delta(t-k T), \quad \omega_{s}=\frac{2 \pi}{T} \\
& p(t)=\frac{\omega_{s}}{2 \pi} \sum_{k=-\infty}^{\infty} e^{j k \omega_{s} t} \\
& s(t)=\frac{\omega_{s}}{2 \pi} \sum_{k=-\infty}^{\infty} e^{j k \omega_{s} t} x(t) \\
& X=\mathcal{F} x, \quad S=\mathcal{F} s \\
& S(\omega)=\frac{\omega_{s}}{2 \pi} \sum_{k=-\infty}^{\infty} X\left(\omega-k \omega_{s}\right)
\end{aligned}
$$

- Consider frequency spectrum $S$ of the impulse-sampled function $s$ given by

$$
S(\omega)=\frac{\omega_{s}}{2 \pi} \sum_{k=-\infty}^{\infty} X\left(\omega-k \omega_{s}\right)
$$

- The function $S$ is a scaled sum of an infinite number of shifted copies of $X$.
- Two distinct behaviors can result in this summation, depending on $\omega_{s}$ and the bandwidth of $x$.
- In particular, the nonzero portions of the different shifted copies of $X$ can either:

II overlap; or
2 not overlap.

- In the case where overlap occurs, the various shifted copies of $X$ add together in such a way that the original shape of $X$ is lost. This phenomenon is known as aliasing.
- When aliasing occurs, the original function $x$ cannot be recovered from its samples in $y$.


## Model of Sampling: Aliasing (Continued)



Spectrum of Input Function<br>(Bandwidth $\omega_{m}$ )

Spectrum of ImpulseSampled Function: No Aliasing Case $\left(\omega_{s}>2 \omega_{m}\right)$

Spectrum of ImpulseSampled Function: Aliasing Case $\left(\omega_{s} \leq 2 \omega_{m}\right)$

## Model of Interpolation

- For the purposes of analysis, interpolation can be modelled as shown below.
ideal D/C converter

- The reconstruction of a function $x$ from its sequence $y$ of samples (i.e., bandlimited interpolation) consists of the following two steps (in order):

1 Convert the sequence $y$ to the impulse train $s$ by using the samples in $y$ as the weights of successive impulses in $s$ so that $s(t)=\sum_{n=-\infty}^{\infty} y(n) \delta(t-T n)$.
2 Apply the lowpass filter with impulse response $h$ to $s$ to produce $\hat{x}$ so that $\hat{x}(t)=s * h(t)=\sum_{n=-\infty}^{\infty} y(n) \operatorname{sinc}\left[\frac{\pi}{T}(t-T n)\right]$.

- The lowpass filter is used to eliminate the extra copies of the originally-sampled function's spectrum present in the spectrum of $s$.
- Sampling Theorem. Let $x$ be a function with Fourier transform $X$, and suppose that $|X(\omega)|=0$ for all $\omega$ satisfying $|\omega|>\omega_{M}$ (i.e., $x$ is bandlimited to frequencies $\left[-\omega_{M}, \omega_{M}\right]$ ). Then, $x$ is uniquely determined by its samples $y(n)=x(T n)$ for all integer $n$, if

$$
\omega_{s}>2 \omega_{M}
$$

where $\omega_{s}=\frac{2 \pi}{T}$. The preceding inequality is known as the Nyquist
condition. If this condition is satisfied, we have that

$$
x(t)=\sum_{n=-\infty}^{\infty} y(n) \operatorname{sinc}\left[\frac{\pi}{T}(t-T n)\right]
$$

or equivalently (i.e., rewritten in terms of $\omega_{s}$ instead of $T$ ),

$$
x(t)=\sum_{n=-\infty}^{\infty} y(n) \operatorname{sinc}\left(\frac{\omega_{s}}{2} t-\pi n\right) .
$$

- We call $\frac{\omega_{s}}{2}$ the Nyquist frequency and $2 \omega_{M}$ the Nyquist rate.


## Part 7

## Laplace Transform (LT)

- Another important mathematical tool in the study of signals and systems is known as the Laplace transform.
- The Laplace transform can be viewed as a generalization of the (classical) Fourier transform.
- Due to its more general nature, the Laplace transform has a number of advantages over the (classical) Fourier transform.
- First, the Laplace transform representation exists for some functions that do not have a Fourier transform representation. So, we can handle some functions with the Laplace transform that cannot be handled with the Fourier transform.
- Second, since the Laplace transform is a more general tool, it can provide additional insights beyond those facilitated by the Fourier transform.
- Earlier, we saw that complex exponentials are eigenfunctions of LTI systems.
- In particular, for a LTI system $\mathcal{H}$ with impulse response $h$, we have that

$$
\mathcal{H}\left\{e^{s t}\right\}(t)=H(s) e^{s t} \quad \text { where } \quad H(s)=\int_{-\infty}^{\infty} h(t) e^{-s t} d t
$$

- Previously, we referred to $H$ as the system function.
- As it turns out, $H$ is the Laplace transform of $h$.
- Since the Laplace transform has already appeared earlier in the context of LTI systems, it is clearly a useful tool.
- Furthermore, as we will see, the Laplace transform has many additional uses.


## Section 7.1

## Laplace Transform

## (Bilateral) Laplace Transform

- The (bilateral) Laplace transform of the function $x$, denoted $\mathcal{L} x$ or $X$, is defined as

$$
\mathcal{L} x(s)=X(s)=\int_{-\infty}^{\infty} x(t) e^{-s t} d t
$$

- The inverse Laplace transform of $X$, denoted $\mathcal{L}^{-1} X$ or $x$, is then given by

$$
\mathcal{L}^{-1} X(t)=x(t)=\frac{1}{2 \pi j} \int_{\sigma-j \infty}^{\sigma+j \infty} X(s) e^{s t} d s
$$

where $\operatorname{Re}(s)=\sigma$ is in the ROC of $X$. (Note that this is a contour integration, since $s$ is complex.)

- We refer to $x$ and $X$ as a Laplace transform pair and denote this relationship as

$$
x(t) \stackrel{\longleftrightarrow T}{\longleftrightarrow} X(s) .
$$

- In practice, we do not usually compute the inverse Laplace transform by directly using the formula from above. Instead, we resort to other means (to be discussed later).


## Bilateral and Unilateral Laplace Transforms

- Two different versions of the Laplace transform are commonly used:

1 the bilateral (or two-sided) Laplace transform; and
2 the unilateral (or one-sided) Laplace transform.

- The unilateral Laplace transform is most frequently used to solve systems of linear differential equations with nonzero initial conditions.
- As it turns out, the only difference between the definitions of the bilateral and unilateral Laplace transforms is in the lower limit of integration.
- In the bilateral case, the lower limit is $-\infty$, whereas in the unilateral case, the lower limit is 0 (i.e., $\int_{-\infty}^{\infty} x(t) e^{-s t} d t$ versus $\int_{0^{-}}^{\infty} x(t) e^{-s t} d t$ ).
- For the most part, we will focus our attention primarily on the bilateral Laplace transform.
- We will, however, briefly introduce the unilateral Laplace transform as a tool for solving differential equations.
- Unless otherwise noted, all subsequent references to the Laplace transform should be understood to mean bilateral Laplace transform.
- For a function $x$, the Laplace transform of $x$ is denoted using operator notation as $\mathcal{L} x$.
- The Laplace transform of $x$ evaluated at $s$ is denoted $\mathcal{L} x(s)$.
- Note that $\mathcal{L} x$ is a function, whereas $\mathcal{L} x(s)$ is a number.
- Similarly, for a function $X$, the inverse Laplace transform of $X$ is denoted using operator notation as $\mathcal{L}^{-1} X$.
- The inverse Laplace transform of $X$ evaluated at $t$ is denoted $\mathcal{L}^{-1} X(t)$.
- Note that $\mathcal{L}^{-1} X$ is a function, whereas $\mathcal{L}^{-1} X(t)$ is a number.
- With the above said, engineers often abuse notation, and use expressions like those above to mean things different from their proper meanings.
- Since such notational abuse can lead to problems, it is strongly recommended that one refrain from doing this.
- Often, we would like to write an expression for the Laplace transform of a function without explicitly naming the function.
- For example, consider writing an expression for the Laplace transform of the function $v(t)=x(5 t-3)$ but without using the name " $v$ ".
- It would be incorrect to write " $\mathcal{L} x(5 t-3)$ " as this is the function $\mathcal{L} x$ evaluated at $5 t-3$, which is not the meaning that we wish to convey.
- Also, strictly speaking, it would be incorrect to write " $\mathcal{L}\{x(5 t-3)\}$ " as the operand of the Laplace transform operator must be a function, and $x(5 t-3)$ is a number (i.e., the function $x$ evaluated at $5 t-3$ ).
- Using dot notation, we can write the following strictly-correct expression for the desired Laplace transform: $\mathcal{L}\{x(5 \cdot-3)\}$.
- In many cases, however, it is probably advisable to avoid employing anonymous (i.e., unnamed) functions, as their use tends to be more error prone in some contexts.


## Remarks on Notational Conventions

- Since dot notation is less frequently used by engineers, the author has elected to minimize its use herein.
- To avoid ambiguous notation, the following conventions are followed:

11 in the expression for the operand of a Laplace transform operator, the independent variable is assumed to be the variable named " $t$ " unless otherwise indicated (i.e., in terms of dot notation, each " $t$ " is treated as if it were a ".")
■ in the expression for the operand of the inverse Laplace transform operator, the independent variable is assumed to be the variable named " $s$ " unless otherwise indicated (i.e., in terms of dot notation, each " $s$ " is treated as if it were a ".").

- For example, with these conventions:
- " $\mathcal{L}\{(t-\tau) u(t-\tau)\}$ " denotes the function that is the Laplace transform of the function $v(t)=(t-\tau) u(t-\tau)$ (not the Laplace transform of the function $v(\tau)=(t-\tau) u(t-\tau))$.
- " $\mathcal{L}^{-1}\left\{\frac{1}{s^{2}-\lambda}\right\}$ " denotes the function that is the inverse Laplace transform of the function $V(s)=\left\{\frac{1}{s^{2}-\lambda}\right\}$ (not the inverse Laplace transform of the function $\left.V(\lambda)=\left\{\frac{1}{s^{2}-\lambda}\right\}\right)$.


## Relationship Between Laplace and Fourier Transforms

- Let $X$ and $X_{\mathrm{F}}$ denote the Laplace and (CT) Fourier transforms of $x$, respectively.
- The function $X$ evaluated at $j \omega$ (where $\omega$ is real) yields $X_{\mathrm{F}}(\omega)$. That is,

$$
X(j \omega)=X_{\mathrm{F}}(\omega) .
$$

- Due to the preceding relationship, the Fourier transform of $x$ is sometimes written as $X(j \omega)$.
- The function $X$ evaluated at an arbitrary complex value $s=\sigma+j \omega$ (where $\sigma=\operatorname{Re}(s)$ and $\omega=\operatorname{Im}(s))$ can also be expressed in terms of a Fourier transform involving $x$. In particular, we have

$$
X(\sigma+j \omega)=X_{\mathrm{F}}^{\prime}(\omega),
$$

where $X_{\mathrm{F}}^{\prime}$ is the (CT) Fourier transform of $x^{\prime}(t)=e^{-\sigma t} x(t)$.

- So, in general, the Laplace transform of $x$ is the Fourier transform of an exponentially-weighted version of $x$.
- Due to this weighting, the Laplace transform of a function may exist when the Fourier transform of the same function does not.


## Laplace Transform Examples

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## Section 7.2

## Region of Convergence (ROC)

- The set $R$ of all complex numbers $s$ satisfying

$$
\operatorname{Re}(s)<a
$$

for some real constant $a$ is said to be a left-half plane (LHP).

- Some examples of LHPs are shown below.


- The set $R$ of all complex numbers $s$ satisfying

$$
\operatorname{Re}(s)>a
$$

for some real constant $a$ is said to be a right-half plane (RHP).
■ Some examples of RHPs are shown below.



## Intersection of Sets

■ For two sets $A$ and $B$, the intersection of $A$ and $B$, denoted $A \cap B$, is the set of all points that are in both $A$ and $B$.

- An illustrative example of set intersection is shown below.

$R_{1}$

$R_{2}$

$R_{1} \cap R_{2}$


## Adding a Scalar to a Set

- For a set $S$ and a scalar constant $a, S+a$ denotes the set given by

$$
S+a=\{z+a: z \in S\}
$$

(i.e., $S+a$ is the set formed by adding $a$ to each element of $S$ ).

- Effectively, adding a scalar to a set applies a translation (i.e., shift) to the region associated with the set.
- An illustrative example is given below.


$R+1$


## Multiplying a Set by a Scalar

- For a set $S$ and a scalar constant $a, a S$ denotes the set given by

$$
a S=\{a z: z \in S\}
$$

(i.e., $a S$ is the set formed by multiplying each element of $S$ by $a$ ).

- Multiplying $z$ by $a$ affects $z$ by: scaling by $|a|$ and rotating about the origin by $\arg a$.
- So, effectively, multiplying a set by a scalar applies a scaling and/or rotation to the region associated with the set.
- An illustrative example is given below.


- As we saw earlier, for a function $x$, the complete specification of its Laplace transform $X$ requires not only an algebraic expression for $X$, but also the ROC associated with $X$.
- Two very different functions can have the same algebraic expressions for $X$.
- On the slides that follow, we will examine a number of key properties of the ROC of the Laplace transform.
- The ROC of a Laplace transform consists of strips parallel to the imaginary axis in the complex plane.
- That is, if a point $s_{0}$ is in the ROC, then the vertical line through $s_{0}$ (i.e., $\left.\operatorname{Re}(s)=\operatorname{Re}\left(s_{0}\right)\right)$ is also in the ROC.
- Some examples of sets that would be either valid or invalid as ROCs are shown below.


Valid


Valid


Invalid

- If a Laplace transform $X$ is a rational function, the ROC of $X$ does not contain any poles and is bounded by poles or extends to infinity.
- Some examples of sets that would be either valid or invalid as ROCs of rational Laplace transforms are shown below.


Valid


Valid


Invalid

- If a function $x$ is finite duration and its Laplace transform $X$ converges for at least one point, then $X$ converges for all points in the complex plane (i.e., the ROC is the entire complex plane).
- Some examples of sets that would be either valid or invalid as ROCs for $X$, if $x$ is finite duration, are shown below.


Valid


Invalid


Invalid

- If a function $x$ is right sided and the (vertical) line $\operatorname{Re}(s)=\sigma_{0}$ is in the ROC of the Laplace transform $X=\mathcal{L} x$, then all values of $s$ for which $\operatorname{Re}(s)>\sigma_{0}$ must also be in the ROC (i.e., the ROC includes a RHP containing $\left.\operatorname{Re}(s)=\sigma_{0}\right)$.
- Thus, if $x$ is right sided but not left sided, the ROC of $X$ is a RHP.
- Some examples of sets that would be either valid or invalid as ROCs for $X$, if $x$ is right sided but not left sided, are shown below.


Valid


Invalid


Invalid

- If a function $x$ is left sided and the (vertical) line $\operatorname{Re}(s)=\sigma_{0}$ is in the ROC of the Laplace transform $X=\mathcal{L} x$, then all values of $s$ for which $\operatorname{Re}(s)<\sigma_{0}$ must also be in the ROC (i.e., the ROC includes a LHP containing $\left.\operatorname{Re}(s)=\sigma_{0}\right)$.
- Thus, if $x$ is left sided but not right sided, the ROC of $X$ is a LHP.
- Some examples of sets that would be either valid or invalid as ROCs for $X$, if $x$ is left sided but not right sided, are shown below.


Valid


Invalid


Invalid

- If a function $x$ is two sided and the (vertical) line $\operatorname{Re}(s)=\sigma_{0}$ is in the ROC of the Laplace transform $X=\mathcal{L} x$, then the ROC will consist of a strip in the complex plane that includes the line $\operatorname{Re}(s)=\sigma_{0}$.
- Some examples of sets that would be either valid or invalid as ROCs for $X$, if $x$ is two sided, are shown below.


Valid


Invalid


Invalid

- If the Laplace transform $X$ of a function $x$ is rational (with at least one pole), then:

1 If $x$ is right sided, the ROC of $X$ is to the right of the rightmost pole of $X$ (i.e., the RHP to the right of the rightmost pole).

2 If $x$ is left sided, the ROC of $X$ is to the left of the leftmost pole of $X$ (i.e., the LHP to the left of the leftmost pole).

- This property is implied by properties $1,2,4$, and 5 .
- Some examples of sets that would be either valid or invalid as ROCs for $X$, if $X$ is rational and $x$ is left/right sided, are given below.


Valid


Invalid


Valid


Invalid

## General Form of the ROC

- To summarize the results of properties $3,4,5$, and 6 , if the Laplace transform $X$ of the function $x$ exists, the ROC of $X$ depends on the leftand right-sidedness of $x$ as follows:

| $x$ |  |  |
| :--- | :--- | :--- |
| left sided | right sided | ROC of $X$ |
| no | no | strip |
| no | yes | RHP |
| yes | no | LHP |
| yes | yes | everywhere |

- Thus, we can infer that, if $X$ exists, its ROC can only be of the form of a LHP, a RHP, a vertical strip, or the entire complex plane.
- For example, the sets shown below would not be valid as ROCs.



## Section 7.3

## Properties of the Laplace Transform

## Properties of the Laplace Transform

| Property | Time Domain | Laplace Domain | ROC |
| :--- | :--- | :--- | :--- |
| Linearity | $a_{1} x_{1}(t)+a_{2} x_{2}(t)$ | $a_{1} X_{1}(s)+a_{2} X_{2}(s)$ | At least $R_{1} \cap R_{2}$ |
| Time-Domain Shifting | $x\left(t-t_{0}\right)$ | $e^{-s t_{0}} X(s)$ | $R$ |
| Laplace-Domain Shifting | $e^{s_{0} t} x(t)$ | $X\left(s-s_{0}\right)$ | $R+\operatorname{Re}\left(s_{0}\right)$ |
| Time/Laplace-Domain Scaling | $x(a t)$ | $\frac{1}{\|a\|} X\left(\frac{s}{a}\right)$ | $a R$ |
| Conjugation | $x^{*}(t)$ | $X^{*}\left(s^{*}\right)$ | $R$ |
| Time-Domain Convolution | $x_{1} * x_{2}(t)$ | $X_{1}(s) X_{2}(s)$ | At least $R_{1} \cap R_{2}$ |
| Time-Domain Differentiation | $\frac{d}{d t} x(t)$ | $s X(s)$ | At least $R$ |
| Laplace-Domain Differentiation | $-t x(t)$ | $\frac{d}{d s} X(s)$ | $R$ |
| Time-Domain Integration | $\int_{-\infty}^{t} x(\tau) d \tau$ | $\frac{1}{s} X(s)$ | At least $R \cap\{\operatorname{Re}(s)>0\}$ |


| Property |  |
| :--- | :--- |
| Initial Value Theorem | $x\left(0^{+}\right)=\lim _{s \rightarrow \infty} s X(s)$ |
| Final Value Theorem | $\lim _{t \rightarrow \infty} x(t)=\lim _{s \rightarrow 0} s X(s)$ |

## Laplace Transform Pairs

| Pair | $x(t)$ | $X(s)$ | $\operatorname{ROC}$ |
| :--- | :--- | :--- | :--- |
| 1 | $\delta(t)$ | 1 | $\mathrm{All} s$ |
| 2 | $u(t)$ | $\frac{1}{s}$ | $\operatorname{Re}(s)>0$ |
| 3 | $-u(-t)$ | $\frac{1}{s}$ | $\operatorname{Re}(s)<0$ |
| 4 | $t^{n} u(t)$ | $\frac{n!}{s^{n+1}}$ | $\operatorname{Re}(s)>0$ |
| 5 | $-t^{n} u(-t)$ | $\frac{n!}{s^{n+1}}$ | $\operatorname{Re}(s)<0$ |
| 6 | $e^{-a t} u(t)$ | $\frac{1}{s+a}$ | $\operatorname{Re}(s)>-a$ |
| 7 | $-e^{-a t} u(-t)$ | $\frac{1}{s+a}$ | $\operatorname{Re}(s)<-a$ |
| 8 | $t^{n} e^{-a t} u(t)$ | $\frac{n!}{(s+a)^{n+1}}$ | $\operatorname{Re}(s)>-a$ |
| 9 | $-t^{n} e^{-a t} u(-t)$ | $\frac{n!}{(s+a)^{n+1}}$ | $\operatorname{Re}(s)<-a$ |
| 10 | $\cos \left(\omega_{0} t\right) u(t)$ | $\frac{s}{s^{2}+\omega_{0}^{2}}$ | $\operatorname{Re}(s)>0$ |
| 11 | $\sin \left(\omega_{0} t\right) u(t)$ | $\frac{\omega_{0}}{s^{2}+\omega_{0}^{2}}$ | $\operatorname{Re}(s)>0$ |
| 12 | $e^{-a t} \cos \left(\omega_{0} t\right) u(t)$ | $\frac{s+a}{(s+a)^{2}+\omega_{0}^{2}}$ | $\operatorname{Re}(s)>-a$ |
| 13 | $e^{-a t} \sin \left(\omega_{0} t\right) u(t)$ | $\frac{\omega_{0}}{(s+a)^{2}+\omega_{0}^{2}}$ | $\operatorname{Re}(s)>-a$ |

■ If $x_{1}(t) \stackrel{\longleftrightarrow T}{\longleftrightarrow} X_{1}(s)$ with ROC $R_{1}$ and $x_{2}(t) \stackrel{\longleftrightarrow T}{\longleftrightarrow} X_{2}(s)$ with ROC $R_{2}$, then $a_{1} x_{1}(t)+a_{2} x_{2}(t) \stackrel{\longleftrightarrow T}{\longleftrightarrow} a_{1} X_{1}(s)+a_{2} X_{2}(s)$ with ROC $R$ containing $R_{1} \cap R_{2}$, where $a_{1}$ and $a_{2}$ are arbitrary complex constants.

- This is known as the linearity property of the Laplace transform.
- The ROC $R$ always contains $R_{1} \cap R_{2}$ but can be larger (in the case that pole-zero cancellation occurs).

■ If $x(t) \stackrel{\text { LT }}{\longleftrightarrow} X(s)$ with ROC $R$, then

$$
x\left(t-t_{0}\right) \stackrel{\llcorner T}{\longleftrightarrow} e^{-s t_{0}} X(s) \text { with ROC } R,
$$

where $t_{0}$ is an arbitrary real constant.

- This is known as the time-domain shifting property of the Laplace transform.


## Laplace-Domain Shifting

- If $x(t) \stackrel{\text { LT }}{\longleftrightarrow} X(s)$ with ROC $R$, then

$$
e^{s_{0} t} x(t) \stackrel{L T}{\longleftrightarrow} X\left(s-s_{0}\right) \text { with ROC } R^{\prime}=R+\operatorname{Re}\left(s_{0}\right),
$$

where $s_{0}$ is an arbitrary complex constant.
■ This is known as the Laplace-domain shifting property of the Laplace transform.

- As illustrated below, the ROC $R$ is shifted right by $\operatorname{Re}\left(s_{0}\right)$.


R

$R^{\prime}=R+\operatorname{Re}\left(s_{0}\right)$

■ If $x(t) \stackrel{L T}{\longleftrightarrow} X(s)$ with ROC $R$, then

$$
x(a t) \stackrel{L T}{\longleftrightarrow} \frac{1}{|a|} X\left(\frac{s}{a}\right) \text { with ROC } R^{\prime}=a R,
$$

where $a$ is a nonzero real constant.

- This is known as the (time-domain/Laplace-domain) scaling property of the Laplace transform.
- As illustrated below, the ROC $R$ is scaled and possibly flipped left to right.


R

$R^{\prime}=a R, a>0$


$$
R^{\prime}=a R, a<0
$$

## Conjugation

- If $x(t) \stackrel{L T}{\longleftrightarrow} X(s)$ with ROC $R$, then

$$
x^{*}(t) \stackrel{L T}{\longleftrightarrow} X^{*}\left(s^{*}\right) \text { with ROC } R .
$$

- This is known as the conjugation property of the Laplace transform.

■ If $x_{1}(t) \stackrel{\longleftrightarrow T}{\longleftrightarrow} X_{1}(s)$ with ROC $R_{1}$ and $x_{2}(t) \stackrel{\longleftrightarrow T}{\longleftrightarrow} X_{2}(s)$ with ROC $R_{2}$, then $x_{1} * x_{2}(t) \stackrel{L T}{\longleftrightarrow} X_{1}(s) X_{2}(s)$ with ROC $R$ containing $R_{1} \cap R_{2}$.

- This is known as the time-domain convolution property of the Laplace transform.
- The ROC $R$ always contains $R_{1} \cap R_{2}$ but can be larger than this intersection (if pole-zero cancellation occurs).
- Convolution in the time domain becomes multiplication in the Laplace domain.
- Consequently, it is often much easier to work with LTI systems in the Laplace domain, rather than the time domain.

■ If $x(t) \stackrel{L T}{\longleftrightarrow} X(s)$ with ROC $R$, then

$$
\frac{d x(t)}{d t} \stackrel{\longleftrightarrow}{\longleftrightarrow} s X(s) \text { with ROC } R^{\prime} \text { containing } R .
$$

- This is known as the time-domain differentiation property of the Laplace transform.
- The ROC $R^{\prime}$ always contains $R$ but can be larger than $R$ (if pole-zero cancellation occurs).
- Differentiation in the time domain becomes multiplication by $s$ in the Laplace domain.
- Consequently, it can often be much easier to work with differential equations in the Laplace domain, rather than the time domain.


## Laplace-Domain Differentiation

■ If $x(t) \stackrel{L T}{\longleftrightarrow} X(s)$ with ROC $R$, then

$$
-t x(t) \stackrel{\longleftrightarrow}{\longleftrightarrow} \frac{d X(s)}{d s} \text { with ROC } R \text {. }
$$

- This is known as the Laplace-domain differentiation property of the Laplace transform.

■ If $x(t) \stackrel{L T}{\longleftrightarrow} X(s)$ with ROC $R$, then

$$
\int_{-\infty}^{t} x(\tau) d \tau \stackrel{\longleftrightarrow}{\longleftrightarrow} \frac{1}{s} X(s) \text { with ROC } R^{\prime} \text { containing } R \cap\{\operatorname{Re}(s)>0\} .
$$

- This is known as the time-domain integration property of the Laplace transform.
- The ROC $R^{\prime}$ always contains at least $R \cap\{\operatorname{Re}(s)>0\}$ but can be larger (if pole-zero cancellation occurs).
- Integration in the time domain becomes division by $s$ in the Laplace domain.
- Consequently, it is often much easier to work with integral equations in the Laplace domain, rather than the time domain.
- For a function $x$ with Laplace transform $X$, if $x$ is causal and contains no impulses or higher order singularities at the origin, then

$$
x\left(0^{+}\right)=\lim _{s \rightarrow \infty} s X(s),
$$

where $x\left(0^{+}\right)$denotes the limit of $x(t)$ as $t$ approaches zero from positive values of $t$.

- This result is known as the initial value theorem.
- In situations where $X$ is known but $x$ is not, the initial value theorem eliminates the need to explicitly find $x$ by an inverse Laplace transform calculation in order to evaluate $x\left(0^{+}\right)$.
- In practice, the values of functions at the origin are frequently of interest, as such values often convey information about the initial state of systems.
- The initial value theorem can sometimes also be helpful in checking for errors in Laplace transform calculations.
- For a function $x$ with Laplace transform $X$, if $x$ is causal and $x(t)$ has a finite limit as $t \rightarrow \infty$, then

$$
\lim _{t \rightarrow \infty} x(t)=\lim _{s \rightarrow 0} s X(s)
$$

- This result is known as the finall value theorem.
- In situations where $X$ is known but $x$ is not, the final value theorem eliminates the need to explicitly find $x$ by an inverse Laplace transform calculation in order to evaluate $\lim _{t \rightarrow \infty} x(t)$.
- In practice, the values of functions at infinity are frequently of interest, as such values often convey information about the steady-state behavior of systems.
- The final value theorem can sometimes also be helpful in checking for errors in Laplace transform calculations.


## More Laplace Transform Examples

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## Section 7.4

## Determination of Inverse Laplace Transform

- Recall that the inverse Laplace transform $x$ of $X$ is given by

$$
x(t)=\frac{1}{2 \pi j} \int_{\sigma-j \infty}^{\sigma+j \infty} X(s) e^{s t} d s,
$$

where $\operatorname{Re}(s)=\sigma$ is in the ROC of $X$.

- Unfortunately, the above contour integration can often be quite tedious to compute.
- Consequently, we do not usually compute the inverse Laplace transform directly using the above equation.
- For rational functions, the inverse Laplace transform can be more easily computed using partial fraction expansions.
- Using a partial fraction expansion, we can express a rational function as a sum of lower-order rational functions whose inverse Laplace transforms can typically be found in tables.


## Section 7.5

## Laplace Transform and LTI Systems

■ Consider a LTI system with input $x$, output $y$, and impulse response $h$. Let $X, Y$, and $H$ denote the Laplace transforms of $x, y$, and $h$, respectively.

- Since $y(t)=x * h(t)$, the system is characterized in the Laplace domain by

$$
Y(s)=X(s) H(s)
$$

- As a matter of terminology, we refer to $H$ as the system function (or transfer function) of the system (i.e., the system function is the Laplace transform of the impulse response).
- A LTI system is completely characterized by its system function $H$.
- When viewed in the Laplace domain, a LTI system forms its output by multiplying its input with its system function.
- If the ROC of $H$ includes the imaginary axis, then $H(j \omega)$ is the frequency response of the LTI system.
- Consider a LTI system with input $x$, output $y$, and impulse response $h$, and let $X, Y$, and $H$ denote the Laplace transforms of $x, y$, and $h$, respectively.
- Often, it is convenient to represent such a system in block diagram form in the Laplace domain as shown below.

- Since a LTI system is completely characterized by its system function, we typically label the system with this quantity.


## Interconnection of LTI Systems

- The series interconnection of the LTI systems with system functions $H_{1}$ and $H_{2}$ is the LTI system with system function $H_{1} H_{2}$. That is, we have the equivalence shown below.

- The parallel interconnection of the LTI systems with system functions $H_{1}$ and $H_{2}$ is the LTI system with the system function $H_{1}+H_{2}$. That is, we have the equivalence shown below.



## Causality

- If a LTI system is causal, its impulse response is causal, and therefore right sided. From this, we have the result below.
- Theorem. The ROC associated with the system function of a causal LTI system is a RHP or the entire complex plane.
- In general, the converse of the above theorem is not necessarily true. That is, if the ROC of the system function is a RHP or the entire complex plane, it is not necessarily true that the system is causal.
- If the system function is rational, however, we have that the converse does hold, as indicated by the theorem below.
- Theorem. For a LTI system with a rational system function $H$, causality of the system is equivalent to the ROC of $H$ being the RHP to the right of the rightmost pole or, if $H$ has no poles, the entire complex plane.
- Whether or not a system is BIBO stable depends on the ROC of its system function.
- Theorem. A LTI system is BIBO stable if and only if the ROC of its system function $H$ contains the imaginary axis (i.e., $\operatorname{Re}(s)=0$ ).
- Theorem. A causal LTI system with a (proper) rational system function $H$ is BIBO stable if and only if all of the poles of $H$ lie in the left half of the plane (i.e., all of the poles have negative real parts).


## Invertibility

- A LTI system $\mathcal{H}$ with system function $H$ is invertible if and only if there exists another LTI system with system function $H_{\text {inv }}$ such that

$$
H(s) H_{\mathrm{inv}}(s)=1
$$

in which case $H_{\text {inv }}$ is the system function of $\mathcal{H}^{-1}$ and

$$
H_{\mathrm{inv}}(s)=\frac{1}{H(s)}
$$

- Since distinct systems can have identical system functions (but with differing ROCs), the inverse of a LTI system is not necessarily unique.
- In practice, however, we often desire a stable and/or causal system. So, although multiple inverse systems may exist, we are frequently only interested in one specific choice of inverse system (due to these additional constraints of stability and/or causality).


## LTI Systems and Differential Equations

■ Many LTI systems of practical interest can be represented using an Nth-order linear differential equation with constant coefficients.

- Consider a system with input $x$ and output $y$ that is characterized by an equation of the form

$$
\sum_{k=0}^{N} b_{k}\left(\frac{d}{d t}\right)^{k} y(t)=\sum_{k=0}^{M} a_{k}\left(\frac{d}{d t}\right)^{k} x(t)
$$

where the $a_{k}$ and $b_{k}$ are complex constants and $M \leq N$.

- Let $h$ denote the impulse response of the system, and let $X, Y$, and $H$ denote the Laplace transforms of $x, y$, and $h$, respectively.
- One can show that $H$ is given by

$$
H(s)=\frac{Y(s)}{X(s)}=\frac{\sum_{k=0}^{M} a_{k} s^{k}}{\sum_{k=0}^{N} b_{k} s^{k}} .
$$

- Observe that, for a system of the form considered above, the system function is always rational.


## Section 7.6

## Application: Circuit Analysis

- An electronic circuit is a network of one or more interconnected circuit elements.
- The three most basic types of circuit elements are:

11 resistors;
2 inductors; and
${ }^{3}$ capacitors.

- Two fundamental quantities of interest in electronic circuits are current and voltage.
- Current is the rate at which electric charge flows through some part of a circuit, such as a circuit element, and is measured in units of amperes (A).
- Voltage is the difference in electric potential between two points in a circuit, such as across a circuit element, and is measured in units of volts (V).
- Voltage is essentially a force that makes electric charge (or current) flow.
- A resistor is a circuit element that opposes the flow of current.
- A resistor is characterized by an equation of the form

$$
\left.v(t)=\operatorname{Ri}(t) \quad \text { (or equivalently, } i(t)=\frac{1}{R} v(t)\right),
$$

where $R$ is a nonnegative real constant, and $v$ and $i$ respectively denote the voltage across and current through the resistor as a function of time.

- As a matter of terminology, the quantity $R$ is known as the resistance of the resistor.
- Resistance is measured in units of ohms ( $\Omega$ ).
- In circuit diagrams, a resistor is denoted by the symbol shown below.



## Inductors

- An inductor is a circuit element that converts an electric current into a magnetic field and vice versa.
- An inductor uses the energy stored in a magnetic field in order to oppose changes in current (through the inductor).
- An inductor is characterized by an equation of the form

$$
\left.v(t)=L \frac{d}{d t} i(t) \quad \text { (or equivalently, } i(t)=\frac{1}{L} \int_{-\infty}^{t} v(\tau) d \tau\right)
$$

where $L$ is a nonnegative real constant, and $v$ and $i$ respectively denote the voltage across and current through the inductor as a function of time.

- As a matter of terminology, the quantity $L$ is known as the inductance of the inductor.
- Inductance is measured in units of henrys $(\mathrm{H})$.
- In circuit diagrams, an inductor is denoted by the symbol shown below.



## Capacitors

- A capacitor is a circuit element that stores electric charge.
- A capacitor uses the energy stored in an electric field in order to oppose changes in voltage (across the capacitor).
- A capacitor is characterized by an equation of the form

$$
\left.v(t)=\frac{1}{C} \int_{-\infty}^{t} i(\tau) d \tau \quad \text { (or equivalently, } i(t)=C \frac{d}{d t} v(t)\right)
$$

where $C$ is a nonnegative real constant, and $v$ and $i$ respectively denote the voltage across and current through the capacitor as a function of time.

- As a matter of terminology, the quantity $C$ is known as the capacitance of the capacitor.
- Capacitance is measured in units of farads ( F ).
- In circuit diagrams, a capacitor is denoted by the symbol shown below.



## Circuit Analysis with the Laplace Transform

- The Laplace transform is a very useful tool for circuit analysis.
- The utility of the Laplace transform is partly due to the fact that the differential/integral equations that describe inductors and capacitors are much simpler to express in the Laplace domain than in the time domain.
- Let $v$ and $i$ denote the voltage across and current through a circuit element, and let $V$ and $I$ denote the Laplace transforms of $v$ and $i$, respectively.
- In the Laplace domain, the equations characterizing a resistor, an inductor, and a capacitor respectively become:

$$
\begin{gathered}
\left.V(s)=R I(s) \quad \text { (or equivalently, } I(s)=\frac{1}{R} V(s)\right) ; \\
\left.V(s)=s L I(s) \quad \text { (or equivalently, } I(s)=\frac{1}{s L} V(s)\right) ; \quad \text { and } \\
\left.V(s)=\frac{1}{s C} I(s) \quad \text { (or equivalently, } I(s)=s C V(s)\right)
\end{gathered}
$$

- Note the absence of differentiation and integration in the above equations for an inductor and a capacitor.


## Section 7.7

## Application: Design and Analysis of Control Systems

## Control Systems

- A control system manages the behavior of one or more other systems with some specific goal.
- Typically, the goal is to force one or more physical quantities to assume particular desired values, where such quantities might include: positions, velocities, accelerations, forces, torques, temperatures, or pressures.
- The desired values of the quantities being controlled are collectively viewed as the input of the control system.
- The actual values of the quantities being controlled are collectively viewed as the output of the control system.
- A control system whose behavior is not influenced by the actual values of the quantities being controlled is called an open loop (or non-feedlback) system.
- A control system whose behavior is influenced by the actual values of the quantities being controlled is called a closed loop (or feedlback) system.
- An example of a simple control system would be a thermostat system, which controls the temperature in a room or building.

Reference


Signal

- input: desired value of the quantity to be controlled
- output: actual value of the quantity to be controlled
- error: difference between the desired and actual values
- plant: system to be controlled
- sensor: device used to measure the actual output
- controller: device that monitors the error and changes the input of the plant with the goal of forcing the error to zero
- Often, we want to ensure that a system is BIBO stable.
- The BIBO stability property is more easily characterized in the Laplace domain than in the time domain.
- Therefore, the Laplace domain is extremely useful for the stability analysis of systems.


## Stabilization Example: Unstable Plant

■ causal LTI plant:


$$
P(s)=\frac{10}{s-1}
$$

- ROC of $P$ :

- system is not BIBO stable


## Stabilization Example: Using Pole-Zero Cancellation

- system formed by series interconnection of plant and causal LTI compensator:

- system function $H$ of overall system:

$$
H(s)=W(s) P(s)=\left(\frac{s-1}{10(s+1)}\right)\left(\frac{10}{s-1}\right)=\frac{1}{s+1}
$$

- ROC of $H$ :

- overall system is BIBO stable


## Stabilization Example: Using Feedback (1)

- feedback system (with causal LTI compensator and sensor):


$$
P(s)=\frac{10}{s-1}, \quad C(s)=\beta, \quad Q(s)=1
$$

- system function $H$ of feedback system:

$$
H(s)=\frac{C(s) P(s)}{1+C(s) P(s) Q(s)}=\frac{10 \beta}{s-(1-10 \beta)}
$$

- ROC of $H$ :


■ feedback system is BIBO stable if and only if $1-10 \beta<0$ or equivalently $\beta>\frac{1}{10}$

## Stabilization Example: Using Feedback (2)



$$
\begin{gathered}
R(s)=X(s)-Q(s) Y(s) \\
Y(s)=C(s) P(s) R(s) \\
Y(s)=C(s) P(s) R(s) \\
=C(s) P(s)[X(s)-Q(s) Y(s)] \\
=C(s) P(s) X(s)-C(s) P(s) Q(s) Y(s) \\
{[1+C(s) P(s) Q(s)] Y(s)=C(s) P(s) X(s)} \\
H(s)=\frac{Y(s)}{X(s)}=\frac{C(s) P(s)}{1+C(s) P(s) Q(s)}
\end{gathered}
$$

$$
P(s)=\frac{10}{s-1}, \quad C(s)=\beta, \quad Q(s)=1
$$

$$
\begin{aligned}
H(s) & =\frac{C(s) P(s)}{1+C(s) P(s) Q(s)} \\
& =\frac{\beta\left(\frac{10}{s-1}\right)}{1+\beta\left(\frac{10}{s-1}\right)(1)} \\
& =\frac{10 \beta}{s-1+10 \beta} \\
& =\frac{10 \beta}{s-(1-10 \beta)}
\end{aligned}
$$

- Pole-zero cancellation is not achievable in practice, and therefore it cannot be used to stabilize real-world systems.
■ The theoretical models used to represent real-world systems are only approximations due to many factors, including the following:
- Determining the system function of a system involves measurement, which always has some error.
- A system cannot be built with such precision that it will have exactly some prescribed system function.
- The system function of most systems will vary at least slightly with changes in the physical environment.
- Although a LTI model is used to represent a system, the likely reality is that the system is not exactly LTI, which introduces error.
- Due to approximation error, the effective poles and zeros of the system function will only be approximately where they are expected to be.
- Since pole-zero cancellation requires that a pole and zero be placed at exactly the same location, any error will prevent this cancellation from being achieved.


## Section 7.8

## Unilateral Laplace Transform

## Unilateral Laplace Transform

- The unilateral Laplace transform of the function $x$, denoted $\mathcal{L}_{u} x$ or $X$, is defined as

$$
\mathcal{L}_{u} x(s)=X(s)=\int_{0^{-}}^{\infty} x(t) e^{-s t} d t .
$$

- The unilateral Laplace transform is related to the bilateral Laplace transform as follows:

$$
\mathcal{L}_{u} x(s)=\int_{0^{-}}^{\infty} x(t) e^{-s t} d t=\int_{-\infty}^{\infty} x(t) u(t) e^{-s t} d t=\mathcal{L}\{x u\}(s) .
$$

- In other words, the unilateral Laplace transform of the function $x$ is simply the bilateral Laplace transform of the function $x u$.
- Since $\mathcal{L}_{u} x=\mathcal{L}\{x u\}$ and $x u$ is always a right-sided function, the ROC associated with $\mathcal{L}_{u} x$ is always either a RHP or the entire complex plane.
- For this reason, we often do not explicitly indicate the ROC when working with the unilateral Laplace transform.
- With the unilateral Laplace transform, the same inverse transform equation is used as in the bilateral case.
- The unilateral Laplace transform is only invertible for causal functions.
- In particular, we have

$$
\begin{aligned}
\mathcal{L}_{u}^{-1}\left\{\mathcal{L}_{u} x\right\}(t) & =\mathcal{L}_{u}^{-1}\{\mathcal{L}\{x u\}\}(t) \\
& =\mathcal{L}^{-1}\{\mathcal{L}\{x u\}\}(t) \\
& =x(t) u(t) \\
& = \begin{cases}x(t) & t \geq 0 \\
0 & t<0 .\end{cases}
\end{aligned}
$$

- For a noncausal function $x$, we can only recover $x(t)$ for $t \geq 0$.
- Due to the close relationship between the unilateral and bilateral Laplace transforms, these two transforms have some similarities in their properties.
- Since these two transforms are not identical, however, their properties differ in some cases, often in subtle ways.
- In the unilateral case, we have that:

1 the time-domain convolution property has the additional requirement that the functions being convolved must be causal;
2 the time/Laplace-domain scaling property has the additional constraint that the scaling factor must be positive;
3 the time-domain differentiation property has an extra term in the expression for $\mathcal{L}_{\mathrm{u}}\{\mathcal{D} x\}(t)$, where $\mathcal{D}$ denotes the derivative operator (namely, $-x\left(0^{-}\right)$);
4 the time-domain integration property has a different lower limit in the time-domain integral (namely, $0^{-}$instead of $-\infty$ ); and
5 the time-domain shifting property does not hold (except in special circumstances).

| Property | Time Domain | Laplace Domain |
| :--- | :--- | :--- |
| Linearity | $a_{1} x_{1}(t)+a_{2} x_{2}(t)$ | $a_{1} X_{1}(s)+a_{2} X_{2}(s)$ |
| Laplace-Domain Shifting | $e^{s_{0} t} x(t)$ | $X\left(s-s_{0}\right)$ |
| Time/Laplace-Domain Scaling | $x(a t), a>0$ | $\frac{1}{a} X\left(\frac{s}{a}\right)$ |
| Conjugation | $x^{*}(t)$ | $X^{*}\left(s^{*}\right)$ |
| Time-Domain Convolution | $x_{1} * x_{2}(t), x_{1}$ and $x_{2}$ are causal | $X_{1}(s) X_{2}(s)$ |
| Time-Domain Differentiation | $\frac{d}{d t} x(t)$ | $s X(s)-x\left(0^{-}\right)$ |
| Laplace-Domain Differentiation | $-t x(t)$ | $\frac{d}{d s} X(s)$ |
| Time-Domain Integration | $\int_{0^{-}}^{t} x(\tau) d \tau$ | $\frac{1}{s} X(s)$ |


| Property |  |
| :--- | :--- |
| Initial Value Theorem | $x\left(0^{+}\right)=\lim _{s \rightarrow \infty} s X(s)$ |
| Final Value Theorem | $\lim _{t \rightarrow \infty} x(t)=\lim _{s \rightarrow 0} s X(s)$ |

## Unilateral Laplace Transform Pairs

| Pair | $x(t), t \geq 0$ | $X(s)$ |
| :--- | :--- | :--- |
| 1 | $\delta(t)$ | 1 |
| 2 | 1 | $\frac{1}{s}$ |
| 3 | $t^{n}$ | $\frac{n!}{s^{n+1}}$ |
| 4 | $e^{-a t}$ | $\frac{1}{s+a}$ |
| 5 | $t^{n} e^{-a t}$ | $\frac{n!}{(s+a)^{n+1}}$ |
| 6 | $\cos \left(\omega_{0} t\right)$ | $\frac{s}{s^{2}+\omega_{0}^{2}}$ |
| 7 | $\sin \left(\omega_{0} t\right)$ | $\frac{\omega_{0}}{s^{2}+\omega_{0}^{2}}$ |
| 8 | $e^{-a t} \cos \left(\omega_{0} t\right)$ | $\frac{s+a}{(s+a)^{2}+\omega_{0}^{2}}$ |
| 9 | $e^{-a t} \sin \left(\omega_{0} t\right)$ | $\frac{\left(+\omega_{0}\right.}{(s+a)^{2}+\omega_{0}^{2}}$ |

- Many systems of interest in engineering applications can be characterized by constant-coefficient linear differential equations.
- One common use of the unilateral Laplace transform is in solving constant-coefficient linear differential equations with nonzero initial conditions.


## Part 8

## Discrete-Time (DT) Signals and Systems

## Section 8.1

## Independent- and Dependent-Variable Transformations

- Time shifting (also called translation) maps the input sequence $x$ to the output sequence $y$ as given by

$$
y(n)=x(n-b),
$$

where $b$ is an integer.

- Such a transformation shifts the sequence (to the left or right) along the time axis.
- If $b>0, y$ is shifted to the right by $|b|$, relative to $x$ (i.e., delayed in time).
- If $b<0, y$ is shifted to the left by $|b|$, relative to $x$ (i.e., advanced in time).

Time Shifting (Translation): Example




- Time reversal (also known as reflection) maps the input sequence $x$ to the output sequence $y$ as given by

$$
y(n)=x(-n) .
$$

- Geometrically, the output sequence $y$ is a reflection of the input sequence $x$ about the (vertical) line $n=0$.




## Downsampling

- Downsampling maps the input sequence $x$ to the output sequence $y$ as given by

$$
y(n)=(\downarrow a) x(n)=x(a n),
$$

where $a$ is a strictly positive integer.

- The output sequence $y$ is produced from the input sequence $x$ by keeping only every $a$ th sample of $x$.


- Upsampling maps the input sequence $x$ to the output sequence $y$ as given by

$$
y(n)=(\uparrow a) x(n)= \begin{cases}x(n / a) & n / a \text { is an integer } \\ 0 & \text { otherwise },\end{cases}
$$

where $a$ is a strictly positive integer.

- The output sequence $y$ is produced from the input sequence $x$ by inserting $a-1$ zeros between all of the samples of $x$.




## Combined Independent-Variable Transformations

- Consider a transformation that maps the input sequence $x$ to the output sequence $y$ as given by

$$
y(n)=x(a n-b),
$$

where $a$ and $b$ are integers and $a \neq 0$.

- Such a transformation is a combination of time shifting, downsampling, and time reversal operations.
- Time reversal commutes with downsampling.
- Time shifting does not commute with time reversal or downsampling.
- The above transformation is equivalent to:

11 first, time shifting $x$ by $b$;
[ then, downsampling the result by $|a|$ and, if $a<0$, time reversing as well.

- If $\frac{b}{a}$ is an integer, the above transformation is also equivalent to:

11 first, downsampling $x$ by $|a|$ and, if $a<0$, time reversing;
2 then, time shifting the result by $\frac{b}{a}$.

- Note that the time shift is not by the same amount in both cases.


## Section 8.2

## Properties of Sequences

- Sums involving even and odd sequences have the following properties:
- The sum of two even sequences is even.
- The sum of two odd sequences is odd.
- The sum of an even sequence and odd sequence is neither even nor odd, provided that neither of the sequences is identically zero.
- That is, the sum of sequences with the same type of symmetry also has the same type of symmetry.
- Products involving even and odd sequences have the following properties:
- The product of two even sequences is even.
- The product of two odd sequences is even.
- The product of an even sequence and an odd sequence is odd.
- That is, the product of sequences with the same type of symmetry is even, while the product of sequences with opposite types of symmetry is odd.


## Decomposition of a Sequence into Even and Odd Parts

- Every sequence $x$ has a unique representation of the form

$$
x(n)=x_{\mathrm{e}}(n)+x_{\mathrm{o}}(n)
$$

where the sequences $x_{\mathrm{e}}$ and $x_{\mathrm{o}}$ are even and odd, respectively.

- In particular, the sequences $x_{\mathrm{e}}$ and $x_{\mathrm{o}}$ are given by

$$
x_{\mathrm{e}}(n)=\frac{1}{2}[x(n)+x(-n)] \quad \text { and } \quad x_{\circ}(n)=\frac{1}{2}[x(n)-x(-n)] .
$$

- The sequences $x_{\mathrm{e}}$ and $x_{\mathrm{o}}$ are called the even part and odd part of $x$, respectively.
- For convenience, the even and odd parts of $x$ are often denoted as Even $\{x\}$ and $\operatorname{Odd}\{x\}$, respectively.
- The least common multiple of two (strictly positive) integers $a$ and $b$, denoted $\operatorname{lcm}(a, b)$, is the smallest positive integer that is divisible by both $a$ and $b$.
- The quantity $\operatorname{lcm}(a, b)$ can be easily determined from a prime factorization of the integers $a$ and $b$ by taking the product of the highest power for each prime factor appearing in these factorizations. Example:

$$
\begin{gathered}
\operatorname{lcm}(20,6)=\operatorname{lcm}\left(2^{2} \cdot 5^{1}, 2^{1} \cdot 3^{1}\right)=2^{2} \cdot 3^{1} \cdot 5^{1}=60 \\
\operatorname{lcm}(54,24)=\operatorname{lcm}\left(2^{1} \cdot 3^{3}, 2^{3} \cdot 3^{1}\right)=2^{3} \cdot 3^{3}=216 ; \quad \text { and } \\
\operatorname{lcm}(24,90)=\operatorname{lcm}\left(2^{3} \cdot 3^{1}, 2^{1} \cdot 3^{2} \cdot 5^{1}\right)=2^{3} \cdot 3^{2} \cdot 5^{1}=360
\end{gathered}
$$

- Sum of periodic sequences. For any two periodic sequences $x_{1}$ and $x_{2}$ with fundamental periods $N_{1}$ and $N_{2}$, respectively, the sum $x_{1}+x_{2}$ is periodic with period $\operatorname{lcm}\left(N_{1}, N_{2}\right)$.


## Right-Sided Sequences

- A sequence $x$ is said to be right sided if, for some (finite) integer constant $n_{0}$, the following condition holds:

$$
x(n)=0 \quad \text { for all } n<n_{0}
$$

(i.e., $x$ is only potentially nonzero to the right of $n_{0}$ ).

- An example of a right-sided sequence is shown below.

- A sequence $x$ is said to be causal if

$$
x(n)=0 \quad \text { for all } n<0 .
$$

- A causal sequence is a special case of a right-sided sequence.
- A causal sequence is not to be confused with a causal system. In these two contexts, the word "causal" has very different meanings.


## Left-Sided Sequences

- A sequence $x$ is said to be left sided if, for some (finite) integer constant $n_{0}$, the following condition holds:

$$
x(n)=0 \quad \text { for all } n>n_{0}
$$

(i.e., $x$ is only potentially nonzero to the left of $n_{0}$ ).

- An example of a left-sided sequence is shown below.

- A sequence $x$ is said to be anticausal if

$$
x(n)=0 \quad \text { for all } n>0 .
$$

- An anticausal sequence is a special case of a left-sided sequence.
- An anticausal sequence is not to be confused with an anticausal system. In these two contexts, the word "anticausal" has very different meanings.


## Finite-Duration and Two-Sided Sequences

- A sequence that is both left sided and right sided is said to be finite duration (or time limited).
- An example of a finite-duration sequence is shown below.

- A sequence that is neither left sided nor right sided is said to be two sided.
- An example of a two-sided sequence is shown below.

- A sequence $x$ is said to be bounded if there exists some (finite) positive real constant $A$ such that

$$
|x(n)| \leq A \quad \text { for all } n
$$

(i.e., $x(n)$ is finite for all $n$ ).

- Examples of bounded sequences include any constant sequence.
- Examples of unbounded sequences include any nonconstant polynomial sequence.


## Energy of a Sequence

- The energy $E$ contained in the sequence $x$ is given by

$$
E=\sum_{k=-\infty}^{\infty}|x(k)|^{2}
$$

- A signal with finite energy is said to be an energy signal.

Section 8.3

## Elementary Sequences

## Real Sinusoidal Sequences

- A real sinusoidal sequence is a sequence of the form

$$
x(n)=A \cos (\Omega n+\theta)
$$

where $A, \Omega$, and $\theta$ are real constants.

- A real sinusoid is periodic if and only if $\frac{\Omega}{2 \pi}$ is a rational number, in which case the fundamental period is the smallest integer of the form $\frac{2 \pi k}{|\Omega|}$ where $k$ is a (strictly) positive integer.
■ For all integer $k, x_{k}(n)=A \cos ([\Omega+2 \pi k] n+\theta)$ is the same sequence.
- An example of a periodic real sinusoid with fundamental period 12 is shown plotted below.



## Oscillation Rate of Real Sinusoidal Sequences

■ Unlike their continuous-time counterparts, real sinusoidal sequences have an upper bound on the rate at which they can oscillate.

- Since $x_{k}(n)=A \cos ([\Omega+2 \pi k] n+\theta)$ is the same sequence for all integer $k$, we consider only $0 \leq \Omega<2 \pi$ without loss of generality.
- Consider the set of real sinusoidal sequences of the form

$$
x(n)=A \cos (\Omega n+\theta),
$$

where $0 \leq \Omega<2 \pi$.

- The rate of oscillation of $x$ is least (i.e., $x$ is constant) when $\Omega=0$.
- The rate of oscillation of $x$ is greatest when $\Omega=\pi$.
- As $\Omega$ increases from 0 to $\pi$, the rate of oscillation of $x$ increases.
- As $\Omega$ increases from $\pi$ to $2 \pi$, the rate of oscillation of $x$ decreases.


## Effect of Increasing Frequency on Oscillation Rate








## Complex Exponential Sequences

- A complex exponential sequence is a sequence of the form

$$
x(n)=c a^{n}
$$

where $c$ and $a$ are complex constants.

- Such a sequence can also be equivalently expressed in the form

$$
x(n)=c e^{b n}
$$

where $b$ is a complex constant chosen as $b=\ln a$. (This this form is more similar to that presented for CT complex exponentials).

- A complex exponential can exhibit one of a number of distinct modes of behavior, depending on the values of the parameters $c$ and $a$.
- For example, as special cases, complex exponentials include real exponentials and complex sinusoids.
- A real exponential sequence is a special case of a complex exponential

$$
x(n)=c a^{n}
$$

where $c$ and $a$ are restricted to be real numbers.

- A real exponential can exhibit one of several distinct modes of behavior, depending on the magnitude and sign of $a$.
- If $|a|>1$, the magnitude of $x(n)$ increases exponentially as $n$ increases (i.e., a growing exponential).
- If $|a|<1$, the magnitude of $x(n)$ decreases exponentially as $n$ increases (i.e., a decaying exponential).
- If $|a|=1$, the magnitude of $x(n)$ is a constant, independent of $n$.
- If $a>0, x(n)$ has the same sign for all $n$.
- If $a<0, x(n)$ alternates in sign as $n$ increases/decreases.


$$
|a|>1, a>0 \quad\left[a=\frac{5}{4} ; c=1\right]
$$


$|a|<1, a>0 \quad\left[a=\frac{4}{5} ; c=1\right]$


$$
|a|=1, a>0 \quad[a=1 ; c=1]
$$

## Real Exponential Sequences (Continued 2)


$|a|>1, a<0 \quad\left[a=-\frac{5}{4} ; c=1\right]$

$|a|<1, a<0 \quad\left[a=-\frac{4}{5} ; c=1\right]$


$$
|a|=1, a<0 \quad[a=-1 ; c=1]
$$

## Complex Sinusoidal Sequences

■ A complex sinusoidal sequence is a special case of a complex exponential $x(n)=c a^{n}$, where $c$ and $a$ are complex and $|a|=1$ (i.e., $a$ is of the form $e^{j \Omega}$ where $\Omega$ is real).

- That is, a complex sinusoidal sequence is a sequence of the form

$$
x(n)=c e^{j \Omega n}
$$

where $c$ is complex and $\Omega$ is real.

- Using Euler's relation, we can rewrite $x(n)$ as

$$
x(n)=\underbrace{|c| \cos (\Omega n+\arg c)}_{\operatorname{Re}\{x(n)\}}+j \underbrace{|c| \sin (\Omega n+\arg c)}_{\operatorname{Im}\{x(n)\}} .
$$

- Thus, $\operatorname{Re}\{x\}$ and $\operatorname{Im}\{x\}$ are real sinusoids.
- A complex sinusoid is periodic if and only if $\frac{\Omega}{2 \pi}$ is a rational number, in which case the fundamental period is the smallest integer of the form $\frac{2 \pi k}{|\Omega|}$ where $k$ is a (strictly) positive integer.


## Complex Sinusoidal Sequences (Continued)

- For $x(n)=e^{j(2 \pi / 7) n}$, the graphs of $\operatorname{Re}\{x\}$ and $\operatorname{Im}\{x\}$ are shown below.




## Oscillation Rate of Complex Sinusoidal Sequences

- Unlike their continuous-time counterparts, complex sinusoidal sequences have an upper bound on the rate at which they can oscillate.
- Since $x_{k}(n)=c e^{j(\Omega+2 \pi k) n}$ is the same sequence for all integer $k$, we consider only $0 \leq \Omega<2 \pi$ without loss of generality.
■ Consider the set of complex sinusoidal sequences of the form

$$
x(n)=c e^{j \Omega n}
$$

where $0 \leq \Omega<2 \pi$.

- The rate of oscillation of $x$ is least (i.e., $x$ is constant) when $\Omega=0$.
- The rate of oscillation of $x$ is greatest when $\Omega=\pi$.

■ As $\Omega$ increases from 0 to $\pi$, the rate of oscillation of $x$ increases.

- As $\Omega$ increases from $\pi$ to $2 \pi$, the rate of oscillation of $x$ decreases.


## General Complex Exponential Sequences

- In the most general case of a complex exponential sequence $x(n)=c a^{n}$, $c$ and $a$ are both complex.
- Letting $c=|c| e^{j \theta}$ and $a=|a| e^{j \Omega}$ where $\theta$ and $\Omega$ are real, and using Euler's relation, we can rewrite $x(n)$ as

$$
x(n)=\underbrace{|c||a|^{n} \cos (\Omega n+\theta)}_{\operatorname{Re}\{x(n)\}}+j \underbrace{|c||a|^{n} \sin (\Omega n+\theta)}_{\operatorname{Im}\{x(n)\}} .
$$

- Thus, $\operatorname{Re}\{x\}$ and $\operatorname{Im}\{x\}$ are each the product of a real exponential and real sinusoid.
- One of several distinct modes of behavior is exhibited by $x$, depending on the value of $a$.
- If $|a|=1, \operatorname{Re}\{x\}$ and $\operatorname{Im}\{x\}$ are real sinusoids.
- If $|a|>1, \operatorname{Re}\{x\}$ and $\operatorname{Im}\{x\}$ are each the product of a real sinusoid and a growing real exponential.
- If $|a|<1, \operatorname{Re}\{x\}$ and $\operatorname{Im}\{x\}$ are each the product of a real sinusoid and a decaying real exponential.


## General Complex Exponential Sequences (Continued)

- The various modes of behavior for $\operatorname{Re}\{x\}$ and $\operatorname{Im}\{x\}$ are illustrated below.

$|a|>1$

$|a|<1$


$$
|a|=1
$$

## Relationship Between Complex Exponentials and Real Sinusoids

- From Euler's relation, a complex sinusoid can be expressed as the sum of two real sinusoids as

$$
c e^{j \Omega n}=c \cos (\Omega n)+j c \sin (\Omega n) .
$$

- Moreover, a real sinusoid can be expressed as the sum of two complex sinusoids using the identities

$$
\begin{gathered}
c \cos (\Omega n+\theta)=\frac{c}{2}\left[e^{j(\Omega n+\theta)}+e^{-j(\Omega n+\theta)}\right] \text { and } \\
c \sin (\Omega n+\theta)=\frac{c}{2 j}\left[e^{j(\Omega n+\theta)}-e^{-j(\Omega n+\theta)}\right] .
\end{gathered}
$$

- Note that, above, we are simply restating results from the (appendix) material on complex analysis.


## Unit-Step Sequence

- The unit-step sequence, denoted $u$, is defined as

$$
u(n)= \begin{cases}1 & n \geq 0 \\ 0 & \text { otherwise }\end{cases}
$$

- A plot of this sequence is shown below.



## Unit Rectangular Pulses

- A unit rectangular pulse is a sequence of the form

$$
p(n)= \begin{cases}1 & a \leq n<b \\ 0 & \text { otherwise }\end{cases}
$$

where $a$ and $b$ are integer constants satisfying $a<b$.

- Such a sequence can be expressed in terms of the unit-step sequence as

$$
p(n)=u(n-a)-u(n-b) .
$$

- The graph of a unit rectangular pulse has the general form shown below.



## Unit-Impulse Sequence

- The unit-impulse sequence (also known as the delta sequence), denoted $\delta$, is defined as

$$
\delta(n)= \begin{cases}1 & n=0 \\ 0 & \text { otherwise }\end{cases}
$$

- The first-order difference of $u$ is $\delta$. That is,

$$
\delta(n)=u(n)-u(n-1) .
$$

- The running sum of $\delta$ is $u$. That is,

$$
u(n)=\sum_{k=-\infty}^{n} \delta(k)
$$

- A plot of $\delta$ is shown below.


■ For any sequence $x$ and any integer constant $n_{0}$, the following identity holds:

$$
x(n) \boldsymbol{\delta}\left(n-n_{0}\right)=x\left(n_{0}\right) \boldsymbol{\delta}\left(n-n_{0}\right) .
$$

■ For any sequence $x$ and any integer constant $n_{0}$, the following identity holds:

$$
\sum_{n=-\infty}^{\infty} x(n) \delta\left(n-n_{0}\right)=x\left(n_{0}\right)
$$

- Trivially, the sequence $\delta$ is also even.
- For integer constants $a$ and $b$ where $a<b$, consider a sequence $x$ of the form

$$
x(n)= \begin{cases}1 & a \leq n<b \\ 0 & \text { otherwise }\end{cases}
$$

(i.e., $x$ is a rectangular pulse of height one that is nonzero from $a$ to $b-1$ inclusive).

- The sequence $x$ can be equivalently written as

$$
x(n)=u(n-a)-u(n-b)
$$

(i.e., the difference of two time-shifted unit-step sequences).

- Unlike the original expression for $x$, this latter expression for $x$ does not involve multiple cases.
- In effect, by using unit-step sequences, we have collapsed a formula involving multiple cases into a single expression.
- The idea from the previous slide can be extended to handle any sequence that is defined in a piecewise manner (i.e., via an expression involving multiple cases).
- That is, by using unit-step sequences, we can always collapse a formula involving multiple cases into a single expression.
■ Often, simplifying a formula in this way can be quite beneficial.


## Section 8.4

## Discrete-Time (DT) Systems

## DT Systems

- A system with input $x$ and output $y$ can be described by the equation

$$
y=\mathcal{H} x
$$

where $\mathcal{H}$ denotes an operator (i.e., transformation).

- Note that the operator $\mathcal{H}$ maps a sequence to a sequence (not a number to a number).
- Alternatively, we can express the above relationship using the notation

$$
x \xrightarrow{\mathcal{H}} y .
$$

- If clear from the context, the operator $\mathcal{H}$ is often omitted, yielding the abbreviated notation

$$
x \rightarrow y .
$$

- Note that the symbols " $\rightarrow$ " and " $=$ " have very different meanings.

■ The symbol " $\rightarrow$ " should be read as "produces" (not as "equals").

## Block Diagram Representations

- Often, a system defined by the operator $\mathcal{H}$ and having the input $x$ and output $y$ is represented in the form of a block diagram as shown below.



## Interconnection of Systems

- Two basic ways in which systems can be interconnected are shown below.

- A series (or cascade) connection ties the output of one system to the input of the other.
- The overall series-connected system is described by the equation

$$
y=\mathcal{H}_{2} \mathcal{H}_{1} x
$$

- A parallel connection ties the inputs of both systems together and sums their outputs.
- The overall parallel-connected system is described by the equation

$$
y=\mathcal{H}_{1} x+\mathcal{H}_{2} x .
$$

## Section 8.5

## Properties of (DT) Systems

- A system $\mathcal{H}$ is said to be memoryless if, for every integer constant $n_{0}$, $\mathcal{H} x\left(n_{0}\right)$ does not depend on $x(n)$ for some $n \neq n_{0}$.
- In other words, a memoryless system is such that the value of its output at any given point in time can depend on the value of its input at only the same point in time.
- A system that is not memoryless is said to have memory.
- Although simple, a memoryless system is not very flexible, since its current output value cannot rely on past or future values of the input.


## Memory (Continued)

If the system $\mathcal{H}$ is memoryless, the output $\mathcal{H} x$ at $n_{0}$ can depend on the input $x$ only at $n_{0}$.


Consider the calculation of the output $\mathcal{H} x$ at $n_{0}$.

## Causality

■ A system $\mathcal{H}$ is said to be causal if, for every integer constant $n_{0}, \mathcal{H} x\left(n_{0}\right)$ does not depend on $x(n)$ for some $n>n_{0}$.

- In other words, a causal system is such that the value of its output at any given point in time can depend on the value of its input at only the same or earlier points in time (i.e., not later points in time).
- If the independent variable $n$ represents time, a system must be causal in order to be physically realizable.
- Noncausal systems can sometimes be useful in practice, however, since the independent variable need not always represent time (e.g., the independent variable might represent position).
- A memoryless system is always causal, although the converse is not necessarily true.


## Causality (Continued)

If the system $\mathcal{H}$ is causal, the output $\mathcal{H} x$ at $n_{0}$ can depend on the input $x$ only at points $n \leq n_{0}$.


Consider the calculation of the output $\mathcal{H} x$ at $n_{0}$.

## Invertibility

- The inverse of a system $\mathcal{H}$ is another system $\mathcal{H}^{-1}$ such that, for every sequence $x$,

$$
\mathcal{H}^{-1} \mathcal{H} x=x
$$

(i.e., the system formed by the cascade interconnection of $\mathcal{H}$ followed by $\mathcal{H}^{-1}$ is a system whose input and output are equal).

- A system is said to be invertible if it has a corresponding inverse system (i.e., its inverse exists).
- Equivalently, a system is invertible if its input $x$ can always be uniquely determined from its output $y$.
- An invertible system will always produce distinct outputs from any two distinct inputs.
- To show that a system is invertible, we simply find the inverse system.
- To show that a system is not invertible, we find two distinct inputs that result in identical outputs.
- In practical terms, invertible systems are "nice" in the sense that their effects can be undone.


## Invertibility (Continued)

■ A system $\mathcal{H}^{-1}$ being the inverse of $\mathcal{H}$ means that the following two systems are equivalent (i.e., $\mathcal{H}^{-1} \mathcal{H}$ is an identity):


System 1: $y=\mathcal{H}^{-1} \mathcal{H} x$


System 2: $y=x$

- A system $\mathcal{H}$ is BIBO stable if, for every bounded sequence $x, \mathcal{H} x$ is bounded (i.e., $|x(n)|<\infty$ for all $n$ implies that $|\mathcal{H} x(n)|<\infty$ for all $n$ ).
- In other words, a BIBO stable system is such that it guarantees to always produce a bounded output as long as its input is bounded.
- To show that a system is BIBO stable, we must show that every bounded input leads to a bounded output.
- To show that a system is not BIBO stable, we need only find a single bounded input that leads to an unbounded output.
- In practical terms, a BIBO stable system is well behaved in the sense that, as long as the system input remains finite for all time, the output will also remain finite for all time.
■ Usually, a system that is not BIBO stable will have serious safety issues.
- For example, a portable music player with a battery input of 3.7 volts and headset output of $\infty$ volts would result in one vaporized human (and likely a big lawsuit as well).
- A system $\mathcal{H}$ is said to be time invariant (TII) (or shift invariant (SII)) if, for every sequence $x$ and every integer $n_{0}$, the following condition holds:

$$
\mathcal{H} x\left(n-n_{0}\right)=\mathcal{H} x^{\prime}(n) \text { for all } n, \quad \text { where } \quad x^{\prime}(n)=x\left(n-n_{0}\right)
$$

(i.e., $\mathcal{H}$ commutes with time shifts).

- In other words, a system is time invariant if a time shift (i.e., advance or delay) in the input always results only in an identical time shift in the output.
- A system that is not time invariant is said to be time varying.
- In simple terms, a time invariant system is a system whose behavior does not change with respect to time.
- Practically speaking, compared to time-varying systems, time-invariant systems are much easier to design and analyze, since their behavior does not change with respect to time.
- Let $S_{n_{0}}$ denote an operator that applies a time shift of $n_{0}$ to a sequence (i.e., $\mathcal{S}_{n_{0}} x(n)=x\left(n-n_{0}\right)$ ).
- A system $\mathcal{H}$ is time invariant if and only if the following two systems are equivalent (i.e., $\mathcal{H}$ commutes with $\mathcal{S}_{n_{0}}$ ):


$$
\begin{gathered}
\text { System 1: } y=\mathcal{H}_{n_{0}} x \\
{\left[\begin{array}{c}
y(n)=\mathcal{H} x^{\prime}(n) \\
x^{\prime}(n)=\mathcal{S}_{n_{0}} x(n)=x\left(n-n_{0}\right)
\end{array}\right]}
\end{gathered}
$$



System 2: $y=\mathcal{S}_{n_{0}} \mathcal{F} x$
$\left[y(n)=\mathcal{H}\left(x\left(n-n_{0}\right)\right]\right.$

$$
\left[y(n)=\mathcal{H} x\left(n-n_{0}\right)\right]
$$

## Additivity, Homogeneity, and Linearity

- A system $\mathcal{H}$ is said to be addlitive if, for all sequences $x_{1}$ and $x_{2}$, the following condition holds:

$$
\mathcal{H}\left(x_{1}+x_{2}\right)=\mathcal{H} x_{1}+\mathcal{H} x_{2}
$$

(i.e., $\mathcal{H}$ commutes with sums).

■ A system $\mathcal{H}$ is said to be homogeneous if, for every sequence $x$ and every complex constant $a$, the following condition holds:

$$
\mathcal{H}(a x)=a \mathcal{H} x
$$

(i.e., $\mathcal{H}$ commutes with multiplication by a constant).

- A system that is both additive and homogeneous is said to be linear.
- In other words, a system $\mathcal{H}$ is linear, if for all sequences $x_{1}$ and $x_{2}$ and all complex constants $a_{1}$ and $a_{2}$, the following condition holds:

$$
\mathcal{H}\left(a_{1} x_{1}+a_{2} x_{2}\right)=a_{1} \mathcal{H} x_{1}+a_{2} \mathcal{H} x_{2}
$$

(i.e., $\mathcal{H}$ commutes with linear combinations).

- The linearity property is also referred to as the superposition property.
- Practically speaking, linear systems are much easier to design and analyze than nonlinear systems.


## Additivity, Homogeneity, and Linearity (Continued 1)

- The system $\mathcal{H}$ is additive if and only if the following two systems are equivalent (i.e., $\mathcal{H}$ commutes with addition):



System 2: $y=\mathcal{H} x_{1}+\mathcal{H} x_{2}$

- The system $\mathcal{H}$ is homogeneous if and only if the following two systems are equivalent (i.e., $\mathcal{H}$ commutes with scalar multiplication):


System 1: $y=\mathcal{H}(a x)$


System 2: $y=a \mathcal{H}(x$

## Additivity, Homogeneity, and Linearity (Continued 2)

- The system $\mathcal{H}$ is linear if and only if the following two systems are equivalent (i.e., $\mathcal{H}$ commutes with linear combinations):


System 1: $y=\mathcal{H}\left(a_{1} x_{1}+a_{2} x_{2}\right)$


System 2: $y=a_{1} \mathcal{H} x_{1}+a_{2} \mathcal{H} x_{2}$

- A sequence $x$ is said to be an eigensequence of the system $\mathcal{H}$ with the eigenvalue $\lambda$ if

$$
\mathcal{H} x=\lambda x
$$

where $\lambda$ is a complex constant.

- In other words, the system $\mathcal{H}$ acts as an ideal amplifier for each of its eigensequences $x$, where the amplifier gain is given by the corresponding eigenvalue $\lambda$.
- Different systems have different eigensequences.
- Many of the mathematical tools developed for the study of DT systems have eigensequences as their basis.


## Part 9

## Discrete-Time Linear Time-Invariant (LTI) Systems

- In engineering, linear time-invariant (LTI) systems play a very important role.
- Very powerful mathematical tools have been developed for analyzing LTI systems.
- LTI systems are much easier to analyze than systems that are not LTI.
- In practice, systems that are not LTI can be well approximated using LTI models.
■ So, even when dealing with systems that are not LTI, LTI systems still play an important role.


## Section 9.1

## Convolution

- The (DT) convolution of the sequences $x$ and $h$, denoted $x * h$, is defined as the sequence

$$
x * h(n)=\sum_{k=-\infty}^{\infty} x(k) h(n-k)
$$

- The convolution $x * h$ evaluated at the point $n$ is simply a weighted sum of elements of $x$, where the weighting is given by $h$ time reversed and shifted by $n$.
■ Herein, the asterisk symbol (i.e., "*") will always be used to denote convolution, not multiplication.
- As we shall see, convolution is used extensively in the theory of (DT) systems.
- In particular, convolution has a special significance in the context of (DT) LTI systems.
- To compute the convolution

$$
x * h(n)=\sum_{k=-\infty}^{\infty} x(k) h(n-k),
$$

we proceed as follows:
1 Plot $x(k)$ and $h(n-k)$ as a function of $k$.
2 Initially, consider an arbitrarily large negative value for $n$. This will result in $h(n-k)$ being shifted very far to the left on the time axis.
3 Write the mathematical expression for $x * h(n)$.
4 Increase $n$ gradually until the expression for $x * h(n)$ changes form. Record the interval over which the expression for $x * h(n)$ was valid.
5 Repeat steps 3 and 4 until $n$ is an arbitrarily large positive value. This corresponds to $h(n-k)$ being shifted very far to the right on the time axis.
6 The results for the various intervals can be combined in order to obtain an expression for $x * h(n)$ for all $n$.

- The convolution operation is commutative. That is, for any two sequences $x$ and $h$,

$$
x * h=h * x
$$

- The convolution operation is associative. That is, for any sequences $x, h_{1}$, and $h_{2}$,

$$
\left(x * h_{1}\right) * h_{2}=x *\left(h_{1} * h_{2}\right) .
$$

- The convolution operation is distributive with respect to addition. That is, for any sequences $x, h_{1}$, and $h_{2}$,

$$
x *\left(h_{1}+h_{2}\right)=x * h_{1}+x * h_{2} .
$$

- For any sequence $x$,

$$
x(n)=\sum_{k=-\infty}^{\infty} x(k) \delta(n-k)=x * \delta(n)
$$

- Thus, any sequence $x$ can be written in terms of an expression involving $\delta$.
- Moreover, $\delta$ is the convolutional identity. That is, for any sequence $x$,

$$
x * \delta=x
$$

## Circular Convolution

- The convolution of two periodic sequences is usually not well defined.
- This motivates an alternative notion of convolution for periodic sequences known as circular convolution.
- The circular convolution (also known as the DT periodic convolution) of the $N$-periodic sequences $x$ and $h$, denoted $x \circledast h$, is defined as

$$
x \circledast h(n)=\sum_{k=\langle N\rangle} x(k) h(n-k)=\sum_{k=0}^{N-1} x(k) h(\bmod (n-k, N)),
$$

where $\bmod (a, b)$ is the remainder after division when $a$ is divided by $b$.

- The circular convolution and (linear) convolution of the $N$-periodic sequences $x$ and $h$ are related as follows:

$$
x \circledast h(n)=x_{0} * h(n) \quad \text { where } \quad x(n)=\sum_{k=-\infty}^{\infty} x_{0}(n-k N)
$$

(i.e., $x_{0}(n)$ equals $x(n)$ over a single period of $x$ and is zero elsewhere).

## Section 9.2

## Convolution and LTI Systems

## Impulse Response

- The response $h$ of a system $\mathcal{H}$ to the input $\delta$ is called the impulse response of the system (i.e., $h=\mathcal{H} \delta$ ).
■ For any LTI system with input $x$, output $y$, and impulse response $h$, the following relationship holds:

$$
y=x * h .
$$

■ In other words, a LTI system simply computes a convolution.

- Furthermore, a LTI system is completely characterized by its impulse response.
- That is, if the impulse response of a LTI system is known, we can determine the response of the system to any input.


## Step Response

- The response $s$ of a system $\mathcal{H}$ to the input $u$ is called the step response of the system (i.e., $s=\mathcal{H} u$ ).
- The impulse response $h$ and step response $s$ of a system are related as

$$
h(n)=s(n)-s(n-1)
$$

- Therefore, the impulse response of a system can be determined from its step response by (first-order) differencing.


## Block Diagram of LTI Systems

- Often, it is convenient to represent a (DT) LTI system in block diagram form.
- Since such systems are completely characterized by their impulse response, we often label a system with its impulse response.
- That is, we represent a system with input $x$, output $y$, and impulse response $h$, as shown below.



## Interconnection of LTI Systems

- The series interconnection of the LTI systems with impulse responses $h_{1}$ and $h_{2}$ is the LTI system with impulse response $h=h_{1} * h_{2}$. That is, we have the equivalences shown below.

- The parallel interconnection of the LTI systems with impulse responses $h_{1}$ and $h_{2}$ is a LTI system with the impulse response $h=h_{1}+h_{2}$. That is, we have the equivalence shown below.



## Section 9.3

## Properties of LTI Systems

- A LTI system with impulse response $h$ is memoryless if and only if

$$
h(n)=0 \quad \text { for all } n \neq 0 .
$$

- That is, a LTI system is memoryless if and only if its impulse response $h$ is of the form

$$
h(n)=K \delta(n)
$$

where $K$ is a complex constant.

- Consequently, every memoryless LTI system with input $x$ and output $y$ is characterized by an equation of the form

$$
y=x *(K \delta)=K x
$$

(i.e., the system is an ideal amplifier).

- For a LTI system, the memoryless constraint is extremely restrictive (as every memoryless LTI system is an ideal amplifier).


## Causality

- A LTI system with impulse response $h$ is causal if and only if

$$
h(n)=0 \quad \text { for all } n<0
$$

(i.e., $h$ is a causal sequence).

- It is due to the above relationship that we call a sequence $x$, satisfying

$$
x(n)=0 \quad \text { for all } n<0,
$$

a causal sequence.

## Invertibility

- The inverse of a LTI system, if such a system exists, is a LTI system.

■ Let $h$ and $h_{\text {inv }}$ denote the impulse responses of a LTI system and its (LTI) inverse, respectively. Then,

$$
h * h_{\mathrm{inv}}=\delta
$$

- Consequently, a LTI system with impulse response $h$ is invertible if and only if there exists a sequence $h_{\text {inv }}$ such that

$$
h * h_{\mathrm{inv}}=\delta
$$

- Except in simple cases, the above condition is often quite difficult to test.
- A LTI system with impulse response $h$ is BIBO stable if and only if

$$
\sum_{n=-\infty}^{\infty}|h(n)|<\infty
$$

(i.e., $h$ is absolutely summable).

## Eigensequences of LTI Systems

■ As it turns out, every complex exponential is an eigensequence of all LTI systems.

- For a LTI system $\mathcal{H}$ with impulse response $h$,

$$
\mathcal{H}\left\{z^{n}\right\}(n)=H(z) z^{n},
$$

where $z$ is a complex constant and

$$
H(z)=\sum_{n=-\infty}^{\infty} h(n) z^{-n} .
$$

- That is, $z^{n}$ is an eigensequence of a LTI system and $H(z)$ is the corresponding eigenvalue.
- We refer to $H$ as the system function (or transfer function) of the system $\mathcal{H}$.
- From above, we can see that the response of a LTI system to a complex exponential is the same complex exponential multiplied by the complex factor $H(z)$.

■ Consider a LTI system with input $x$, output $y$, and system function $H$.

- Suppose that the input $x$ can be expressed as the linear combination of complex exponentials

$$
x(n)=\sum_{k} a_{k} z_{k}^{n},
$$

where the $a_{k}$ and $z_{k}$ are complex constants.

- Using the fact that complex exponentials are eigenfunctions of LTI systems, we can conclude

$$
y(n)=\sum_{k} a_{k} H\left(z_{k}\right) z_{k}^{n} .
$$

- Thus, if an input to a LTI system can be expressed as a linear combination of complex exponentials, the output can also be expressed as linear combination of the same complex exponentials.
- The above formula can be used to determine the output of a LTI system from its input in a way that does not require convolution.


## Part 10

## Discrete-Time Fourier Series (DTFS)

- The Fourier series is a representation for periodic sequences.
- With a Fourier series, a sequence is represented as a linear combination of complex sinusoids.
- The use of complex sinusoids is desirable due to their numerous attractive properties.
- Perhaps, most importantly, complex sinusoids are eigensequences of (DT) LTI systems.


## Section 10.1

## Fourier Series

- A set of periodic complex sinusoids is said to be harmonically related if there exists some constant $\frac{2 \pi}{N}$ such that the fundamental frequency of each complex sinusoid is an integer multiple of $\frac{2 \pi}{N}$.
■ Consider the set of harmonically-related complex sinusoids given by

$$
\phi_{k}(n)=e^{j(2 \pi / N) k n} \quad \text { for all integer } k
$$

- In the above set $\left\{\phi_{k}\right\}$, only $N$ elements are distinct, since

$$
\phi_{k}=\phi_{k+N} \quad \text { for all integer } k .
$$

- Since the fundamental frequency of each of the harmonically-related complex sinusoids is an integer multiple of $\frac{2 \pi}{N}$, a linear combination of these complex sinusoids must be $N$-periodic.


## DT Fourier Series (DTFS)

- An $N$-periodic complex-valued sequence $x$ can be represented as a linear combination of harmonically-related complex sinusoids as

$$
x(n)=\sum_{k=\langle N\rangle} a_{k} e^{j(2 \pi / N) k n}
$$

where $\sum_{k=\langle N\rangle}$ denotes summation over any $N$ consecutive integers (e.g., $[0 . . N-1]$ ). (The summation can be taken over any $N$ consecutive integers, due to the $N$-periodic nature of $x$ and $e^{j(2 \pi / N) k n}$.)

- The above representation of $x$ is known as the (DT) Fourier series and the $a_{k}$ are called Fourier series coefficients.
- The above formula for $x$ is often called the Fourier series synthesis equation.
- To denote that the sequence $x$ has the Fourier series coefficient sequence $a$, we write

$$
x(n) \stackrel{\text { DTFS }}{\longleftrightarrow} a_{k} .
$$

- A periodic sequence $x$ with fundamental period $N$ has the Fourier series coefficient sequence $a$ given by

$$
a_{k}=\frac{1}{N} \sum_{n=\langle N\rangle} x(n) e^{-j(2 \pi / N) k n} .
$$

(The summation can be taken over any $N$ consecutive integers due to the $N$-periodic nature of $x$ and $e^{-j(2 \pi / N) k n}$.)

- The above equation for $a_{k}$ is often referred to as the Fourier series analysis equation.
- Due to the $N$-periodic nature of $x$ and $e^{-j(2 \pi / N) k n}$, the sequence $a$ is also $N$-periodic.


## Convergence of Fourier Series

- Since the analysis and synthesis equations for (DT) Fourier series involve only finite sums (as opposed to infinite series), convergence is not a significant issue of concern.
- If an $N$-periodic sequence is bounded (i.e., is finite in value), its Fourier series coefficient sequence will exist and be bounded and the Fourier series analysis and synthesis equations must converge.


## Section 10.2

## Properties of Fourier Series

$$
x(n) \stackrel{\text { DTFS }}{\longleftrightarrow} a_{k} \quad \text { and } \quad y(n) \stackrel{\text { DTFS }}{\longleftrightarrow} b_{k}
$$

| Property | Time Domain | Fourier Domain |
| :--- | :--- | :--- |
| Linearity | $\alpha x(n)+\beta y(n)$ | $\alpha a_{k}+\beta b_{k}$ |
| Translation | $x\left(n-n_{0}\right)$ | $e^{-j k(2 \pi / N) n_{0}} a_{k}$ |
| Modulation | $e^{j(2 \pi / N) k_{0} n} x(n)$ | $a_{k-k_{0}}$ |
| Reflection | $x(-n)$ | $a_{-k}$ |
| Conjugation | $x^{*}(n)$ | $a_{-k}^{*}$ |
| Duality | $a_{n}$ | $\frac{1}{N} x(-k)$ |
| Periodic Convolution | $x \circledast y(n)$ | $N a_{k} b_{k}$ |
| Multiplication | $x(n) y(n)$ | $a \circledast b_{k}$ |

## Property

Parseval's Relation
Even Symmetry
Odd Symmetry
Real / Conjugate Symmetry $x$ is real $\Leftrightarrow a$ is conjugate symmetric

■ Let $x$ and $y$ be $N$-periodic sequences. If $x(n) \stackrel{\text { DTFS }}{\longleftrightarrow} a_{k}$ and $y(n) \stackrel{\text { DTFS }}{\longleftrightarrow} b_{k}$, then

$$
\alpha x(n)+\beta y(n) \stackrel{\text { DTFS }}{\longleftrightarrow} \alpha a_{k}+\beta b_{k},
$$

where $\alpha$ and $\beta$ are complex constants.

- That is, a linear combination of sequences produces the same linear combination of their Fourier series coefficients.

■ Let $x$ denote a periodic sequence with period $N$. If $x(n) \stackrel{\text { DTFS }}{\longleftrightarrow} c_{k}$, then

$$
x\left(n-n_{0}\right) \stackrel{\text { DTFS }}{\longleftrightarrow} e^{-j k(2 \pi / N) n_{0}} c_{k},
$$

where $n_{0}$ is an integer constant.

- In other words, time shifting a periodic sequence changes the argument (but not magnitude) of its Fourier series coefficients.

■ Let $x$ denote a periodic sequence with period $N$. If $x(n) \stackrel{\text { DTFS }}{\longleftrightarrow} c_{k}$, then

$$
e^{j(2 \pi / N) k_{0} n} x(n) \stackrel{\text { DTFS }}{\longleftrightarrow} c_{k-k_{0}},
$$

where $k_{0}$ is an integer constant.

- That is, multiplying a sequence by a complex sinusoid whose frequency is an integer multiple of $\frac{2 \pi}{N}$ results in a translation of the corresponding Fourier series coefficient sequence.

■ Let $x$ denote a periodic sequence with period $N$. If $x(n) \stackrel{\text { DTFS }}{\longleftrightarrow} c_{k}$, then

$$
x(-n) \stackrel{\mathrm{DTFS}}{\longleftrightarrow} c_{-k} .
$$

- That is, time reversing a sequence results in a time reversal of the corresponding Fourier series coefficient sequence.


## Conjugation

■ Let $x$ denote a periodic sequence with period $N$. If $x(n) \stackrel{\text { DTFS }}{\longleftrightarrow} c_{k}$, then

$$
x^{*}(n) \stackrel{\mathrm{DTFS}}{\longleftrightarrow} c_{-k}^{*} .
$$

- In other words, conjugating a sequence has the effect of time reversing and conjugating the corresponding Fourier series coefficient sequence.

■ Let $x$ denote a periodic sequence with period $N$. If $x(n) \stackrel{\text { DTFS }}{\longleftrightarrow} a(k)$, then

$$
a(n) \stackrel{\text { DTFS }}{\longleftrightarrow} \frac{1}{N} x(-k) .
$$

- This is known as the duality property of Fourier series.
- This property follows from the high degree of symmetry in the analysis and synthesis Fourier-series equations, which are respectively given by

$$
x(m)=\sum_{\ell=\langle N\rangle} a(\ell) e^{j(2 \pi / N) \ell m} \quad \text { and } \quad a(m)=\frac{1}{N} \sum_{\ell=\langle N\rangle} x(\ell) e^{-j(2 \pi / N) m \ell}
$$

- That is, the analysis and synthesis equations are identical except for a factor of $N$ and different sign in the parameter for the exponential function.

■ Let $x$ and $y$ be $N$-periodic sequences. If $x(n) \stackrel{\text { DTFS }}{\longleftrightarrow} a_{k}$ and $y(n) \stackrel{\text { DTFS }}{\longleftrightarrow} b_{k}$, then

$$
x \circledast y(n) \stackrel{\text { DTFS }}{\longleftrightarrow} N a_{k} b_{k} .
$$

- That is, periodic convolution of two sequences multiplies their corresponding Fourier series coefficient sequences (up to a scale factor).

■ Let $x$ and $y$ be $N$-periodic sequences. If $x(n) \stackrel{\text { DTFS }}{\longleftrightarrow} a_{k}$ and $y(n) \stackrel{\text { DTFS }}{\longleftrightarrow} b_{k}$, then

$$
x(n) y(n) \stackrel{\text { DTFS }}{\longrightarrow} a \circledast b(k) .
$$

- That is, multiplying two sequences results in a circular convolution of their corresponding Fourier series coefficient sequences.
- A sequence $x$ and its Fourier series coefficient sequence $a$ satisfy the following relationship:

$$
\frac{1}{N} \sum_{n=\langle N\rangle}|x(n)|^{2}=\sum_{k=\langle N\rangle}\left|a_{k}\right|^{2} .
$$

- The above relationship is simply stating that the amount of energy in a single period of $x$ and the amount of energy in a single period of $a$ are equal up to a scale factor.
- In other words, the transformation between a sequence and its Fourier series coefficient sequence preserves energy (up to a scale factor).
- For an $N$-periodic sequence $x$ with Fourier-series coefficient sequence $a$, the following properties hold:

$$
\begin{gathered}
x \text { is even } \Leftrightarrow a \text { is even; and } \\
x \text { is odd } \Leftrightarrow a \text { is odd. }
\end{gathered}
$$

- In other words, the even/odd symmetry properties of $x$ and $a$ always match.
- A sequence $x$ is real if and only if its Fourier series coefficient sequence $a$ satisfies

$$
a_{k}=a_{-k}^{*} \text { for all } k
$$

(i.e., $a$ is conjugate symmetric).

- From properties of complex numbers, one can show that $a_{k}=a_{-k}^{*}$ is equivalent to

$$
\left|a_{k}\right|=\left|a_{-k}\right| \quad \text { and } \quad \arg a_{k}=-\arg a_{-k}
$$

(i.e., $\left|a_{k}\right|$ is even and $\arg a_{k}$ is $o d d$ ).

- Note that $x$ being real does not necessarily imply that $a$ is real.


## Trigonometric Form of a Fourier Series

- Consider the $N$-periodic sequence $x$ with Fourier series coefficient sequence $a$.
- If $x$ is real, then its Fourier series can be rewritten in trigonometric form as shown below.
- The trigonometric form of a Fourier series has the appearance

$$
x(n)= \begin{cases}\alpha_{0}+\sum_{k=1}^{N / 2-1}\left[\alpha_{k} \cos \left(\frac{2 \pi k n}{N}\right)+\beta_{k} \sin \left(\frac{2 \pi k n}{N}\right)\right]+ & \\ \alpha_{N / 2} \cos (\pi n) & N \text { even } \\ \alpha_{0}+\sum_{k=1}^{(N-1) / 2}\left[\alpha_{k} \cos \left(\frac{2 \pi k n}{N}\right)+\beta_{k} \sin \left(\frac{2 \pi k n}{N}\right)\right] & N \text { odd },\end{cases}
$$

where $\alpha_{0}=a_{0}, \alpha_{N / 2}=a_{N / 2}, \alpha_{k}=2 \operatorname{Re} a_{k}$, and $\beta_{k}=-2 \operatorname{Im} a_{k}$.

- Note that the above trigonometric form contains only real quantities.


## Other Properties of Fourier Series

- For an $N$-periodic sequence $x$ with Fourier-series coefficient sequence $a$, the following properties hold:
$1 a_{0}$ is the average value of $x$ over a single period;
$2 x$ is real and even $\Leftrightarrow a$ is real and even; and
$3 x$ is real and odd $\Leftrightarrow a$ is purely imaginary and odd.


## Section 10.3

## Discrete Fourier Transform (DFT)

■ Letting $a_{k}^{\prime}=N a_{k}$, we can rewrite the Fourier series synthesis and analysis equations, respectively, as

$$
x(n)=\frac{1}{N} \sum_{k=0}^{N-1} a_{k}^{\prime} e^{j(2 \pi / N) k n} \quad \text { and } \quad a_{k}^{\prime}=\sum_{n=0}^{N-1} x(n) e^{-j(2 \pi / N) k n}
$$

- Since $x$ and $a^{\prime}$ are both $N$-periodic, each of these sequences is completely characterized by its $N$ samples over a single period.
- If we only consider the behavior of $x$ and $a^{\prime}$ over a single period, this leads to the equations

$$
\begin{array}{cc}
x(n)=\frac{1}{N} \sum_{k=0}^{N-1} a_{k}^{\prime} e^{j(2 \pi / N) k n} \quad \text { for } n \in[0 \ldots N-1] \quad \text { and } \\
a_{k}^{\prime}=\sum_{n=0}^{N-1} x(n) e^{-j(2 \pi / N) k n} \quad \text { for } k \in[0 \ldots N-1] .
\end{array}
$$

- As it turns out, the above two equations define what is known as the discrete Fourier transform (DFT).


## Discrete Fourier Transform (DFT)

- The discrete Fourier transform (DFT) $X$ of the $N$-element sequence $x$ is defined as

$$
X(k)=\sum_{n=0}^{N-1} x(n) e^{-j(2 \pi / N) k n} \quad \text { for } k \in[0 \ldots N-1] .
$$

- The preceding equation is known as the DFT analysis equation.
- The inverse DFT $x$ of the $N$-element sequence $X$ is given by

$$
x(n)=\frac{1}{N} \sum_{k=0}^{N-1} X(k) e^{j(2 \pi / N) k n} \quad \text { for } n \in[0 . . N-1]
$$

- The preceding equation is known as the DFT synthesis equation.
- The DFT maps a finite-length sequence of $N$ samples to another finite-length sequence of $N$ samples.
- The DFT will be considered in more detail later.

| Property | Time Domain | Fourier Domain |
| :--- | :--- | :--- |
| Linearity | $a_{1} x_{1}(n)+a_{2} x_{2}(n)$ | $a_{1} X_{1}(k)+a_{2} X_{2}(k)$ |
| Translation | $x\left(n-n_{0}\right)$ | $e^{-j k(2 \pi / N) n_{0}} X(k)$ |
| Modulation | $e^{j(2 \pi / N) k_{0} n} x(n)$ | $X\left(k-k_{0}\right)$ |
| Reflection | $x(-n)$ | $X(-k)$ |
| Conjugation | $x^{*}(n)$ | $X^{*}(-k)$ |
| Duality | $X(n)$ | $N x(-k)$ |
| Periodic Convolution | $x_{1} \circledast x_{2}(n)$ | $X_{1}(k) X_{2}(k)$ |
| Multiplication | $x_{1}(n) x_{2}(n)$ | $\frac{1}{N} X_{1} \circledast X_{2}(k)$ |

Property

Parseval's Relation
Even Symmetry
Odd Symmetry
Real / Conjugate Symmetry $x$ is real $\Leftrightarrow X$ is conjugate symmetric

## Section 10.4

## Fourier Series and Frequency Spectra

- The Fourier series provides us with an entirely new way to view sequences.
- Instead of viewing a sequence as having information distributed with respect to time (i.e., a function whose domain is time), we view a sequence as having information distributed with respect to frequency (i.e., a function whose domain is frequency).
- This so called frequency-domain perspective is of fundamental importance in engineering.
- Many engineering problems can be solved much more easily using the frequency domain than the time domain.
- The Fourier series coefficients of a sequence $x$ provide a means to quantify how much information $x$ has at different frequencies.
- The distribution of information in a sequence over different frequencies is referred to as the frequency spectrum of the sequence.
- To gain further insight into the role played by the Fourier series coefficients $a_{k}$ in the context of the frequency spectrum of the $N$-periodic sequence $x$, it is helpful to write the Fourier series with the $a_{k}$ expressed in polar form as

$$
x(n)=\sum_{k=0}^{N-1} a_{k} e^{j(2 \pi / N) k n}=\sum_{k=0}^{N-1}\left|a_{k}\right| e^{j\left([2 \pi / N] k n+\arg a_{k}\right)}
$$

- Clearly, the $k$ th term in the summation corresponds to a complex sinusoid with fundamental frequency $\frac{2 \pi}{N} k$ that has been amplitude scaled by a factor of $\left|a_{k}\right|$ and time-shifted by an amount that depends on $\arg a_{k}$.
- For a given $k$, the larger $\left|a_{k}\right|$ is, the larger is the amplitude of its corresponding complex sinusoid $e^{j(2 \pi / N) k n}$, and therefore the larger the contribution the $k$ th term (which is associated with frequency $\frac{2 \pi}{N} k$ ) will make to the overall summation.
- In this way, we can use $\left|a_{k}\right|$ as a measure of how much information a sequence $x$ has at the frequency $\frac{2 \pi}{N} k$.
- The Fourier series coefficients $a_{k}$ of the sequence $x$ are referred to as the frequency spectrum of $x$.
- The magnitudes $\left|a_{k}\right|$ of the Fourier series coefficients $a_{k}$ are referred to as the magnitude spectrum of $x$.
- The arguments $\arg a_{k}$ of the Fourier series coefficients $a_{k}$ are referred to as the phase spectrum of $x$.
- The frequency spectrum $a_{k}$ of an $N$-periodic sequence is $N$-periodic in the coefficient index $k$ and $2 \pi$-periodic in the frequency $\Omega=\frac{2 \pi}{N} k$.
- The range of frequencies between $-\pi$ and $\pi$ are referred to as the baseband.
- Often, the spectrum of a sequence is plotted against frequency $\Omega=\frac{2 \pi}{N} k$ (over the single $2 \pi$ period of the baseband) instead of the Fourier series coefficient index $k$.
- Since the Fourier series only has frequency components at integer multiples of the fundamental frequency, the frequency spectrum is discrete in the independent variable (i.e., frequency).
- Due to the general appearance of frequency-spectrum plot (i.e., a number of vertical lines at various frequencies), we refer to such spectra as line spectra.
- Let $x$ denote an $N$-periodic sequence with the corresponding Fourier-series coefficient sequence $c$.
- As we saw earlier:
$x$ is real $\Leftrightarrow c$ is conjugate symmetric.
- Furthermore, if $x$ is real, the following assertions hold for $c_{k}$ for $k \in[0 . . N-1]$ :
$11 c_{k}=c_{N-k}^{*}$ for $k \in[1 . . N-1]$;
■ of the $N$ coefficients $c_{k}$ for $k \in[0 \ldots N-1]$, only $\left\lfloor\frac{N}{2}\right\rfloor+1$ coefficients are independent; for example, $c_{k}$ for $k \in\left[0 . .\left\lfloor\frac{N}{2}\right\rfloor\right]$ completely determines $c_{k}$ for all $k \in[0 . . N-1]$;
${ }^{3} c_{0}$ is real; and
4 if $N$ is even, $c_{N / 2}$ is real.
■ Note that approximately half of the coefficients in a single period of $c$ are redundant if $x$ is real.


## Section 10.5

## Fourier Series and LTI Systems

## Frequency Response

■ Recall that a LTI system $\mathcal{H}$ with impulse response $h$ is such that $\mathcal{H}\left\{z^{n}\right\}(n)=H_{\mathrm{Z}}(z) z^{n}$, where $H_{\mathrm{Z}}(z)=\sum_{n=-\infty}^{\infty} h(n) z^{-n}$. (That is, complex exponentials are eigensequences of LTI systems.)

- Since a complex sinusoid is a special case of a complex exponential, we can reuse the above result for the special case of complex sinusoids.
- For a LTI system $\mathcal{H}$ with impulse response $h$,

$$
\mathcal{H}\left\{e^{j \Omega n}\right\}(n)=H(\Omega) e^{j \Omega n}
$$

where $\Omega$ is real and

$$
H(\Omega)=\sum_{n=-\infty}^{\infty} h(n) e^{-j \Omega n}
$$

■ That is, $e^{j \Omega n}$ is an eigensequence of a LTI system and $H(\Omega)$ is the corresponding eigenvalue.

- The function $H$ is $2 \pi$-periodic, since $e^{j \Omega}$ is $2 \pi$-periodic.
- We refer to $H$ as the frequency response of the system $\mathcal{H}$.

■ Consider a LTI system with input $x$, output $y$, and frequency response $H$.

- Suppose that the $N$-periodic input $x$ is expressed as the Fourier series

$$
x(n)=\sum_{k=0}^{N-1} a_{k} e^{j k \Omega_{0} n}, \quad \text { where } \Omega_{0}=\frac{2 \pi}{N}
$$

- Using our knowledge about the eigensequences of LTI systems, we can conclude

$$
y(n)=\sum_{k=0}^{N-1} a_{k} H\left(k \Omega_{0}\right) e^{j k \Omega_{0} n}
$$

- Thus, if the input $x$ to a LTI system is a Fourier series, the output $y$ is also a Fourier series. More specifically, if $x(n) \stackrel{\text { DTFS }}{\longleftrightarrow} a_{k}$ then $y(n) \stackrel{\text { DTFS }}{\longleftrightarrow} H\left(k \Omega_{0}\right) a_{k}$.
- The above formula can be used to determine the output of a LTI system from its input in a way that does not require convolution.

■ In many applications, we want to modify the spectrum of a sequence by either amplifying or attenuating certain frequency components.

- This process of modifying the frequency spectrum of a sequence is called filtering.
- A system that performs a filtering operation is called a filter.
- Many types of filters exist.
- Frequency selective filters pass some frequencies with little or no distortion, while significantly attenuating other frequencies.
- Several basic types of frequency-selective filters include: lowpass, highpass, and bandpass.


## Ideal Lowpass Filter

- An ideal lowpass filter eliminates all baseband frequency components with a frequency whose magnitude is greater than some cutoff frequency, while leaving the remaining baseband frequency components unaffected.
- Such a filter has a frequency response of the form

$$
H(\Omega)= \begin{cases}1 & |\Omega| \leq \Omega_{c} \\ 0 & \Omega_{c}<|\Omega| \leq \pi\end{cases}
$$

where $\Omega_{c}$ is the cutoff frequency.

- A plot of this frequency response is given below.



## Ideal Highpass Filter

- An ideal highpass filter eliminates all baseband frequency components with a frequency whose magnitude is less than some cutoff frequency, while leaving the remaining baseband frequency components unaffected.
- Such a filter has a frequency response of the form

$$
H(\Omega)= \begin{cases}1 & \Omega_{c}<|\Omega| \leq \pi \\ 0 & |\Omega| \leq \Omega_{c}\end{cases}
$$

where $\Omega_{c}$ is the cutoff frequency.

- A plot of this frequency response is given below.



## Ideal Bandpass Filter

- An ideal lbandpass filter eliminates all baseband frequency components with a frequency whose magnitude does not lie in a particular range, while leaving the remaining baseband frequency components unaffected.
- Such a filter has a frequency response of the form

$$
H(\Omega)= \begin{cases}1 & \Omega_{c 1} \leq|\Omega| \leq \Omega_{c 2} \\ 0 & |\Omega|<\Omega_{c 1} \text { or } \Omega_{c 2}<|\Omega|<\pi\end{cases}
$$

where the limits of the passband are $\Omega_{c 1}$ and $\Omega_{c 2}$.

- A plot of this frequency response is given below.



## Part 11

## Discrete-Time Fourier Transform (DTFT)

- The (DT) Fourier series provide an extremely useful representation for periodic sequences.
- Often, however, we need to deal with sequences that are not periodic.
- A more general tool than the Fourier series is needed in this case.
- The (DT) Fourier transform can be used to represent both periodic and aperiodic sequences.
- Since the (DT) Fourier transform is essentially derived from (DT) Fourier series through a limiting process, the Fourier transform has many similarities with Fourier series.

Section 11.1

## Fourier Transform

## Development of the Fourier Transform [Aperiodic Case]

- The (DT) Fourier series is an extremely useful signal representation.

■ Unfortunately, this signal representation can only be used for periodic sequences, since a Fourier series is inherently periodic.

- Many sequences are not periodic, however.
- Rather than abandoning Fourier series, one might wonder if we can somehow use Fourier series to develop a representation that can also be applied to aperiodic sequences.
- By viewing an aperiodic sequence as the limiting case of an $N$-periodic sequence where $N \rightarrow \infty$, we can use the Fourier series to develop a signal representation that can be used for aperiodic sequences, known as the Fourier transform.


## Development of the Fourier Transform [Aperiodic Case] (coninueo)

- Recall that the Fourier series representation of an $N$-periodic sequence $x$ is given by

$$
x(n)=\sum_{k=\langle N\rangle} \underbrace{\left(\frac{1}{N} \sum_{\ell=\langle N\rangle} x(\ell) e^{-j(2 \pi / N) k \ell}\right)}_{c_{k}} e^{j(2 \pi / N) k n} .
$$

- In the above representation, if we take the limit as $N \rightarrow \infty$, we obtain

$$
x(n)=\frac{1}{2 \pi} \int_{2 \pi} \underbrace{\left(\sum_{\ell=-\infty}^{\infty} x(\ell) e^{-j \Omega \ell}\right)}_{X(\Omega)} e^{j \Omega n} d \Omega
$$

(i.e., as $N \rightarrow \infty$, the two finite summations become an integral and infinite summation, $\frac{1}{N}$ becomes $\frac{1}{2 \pi} d \Omega$, and $\left(\frac{2 \pi}{N}\right) k$ becomes $\left.\Omega\right)$.

- This representation for aperiodic sequences is known as the Fourier transform representation.
- The classical Fourier transform for aperiodic sequences does not exist (i.e., $\sum_{n=-\infty}^{\infty} x(n) e^{-j \Omega n}$ fails to converge) for some sequences of great practical interest, such as:
- a nonzero constant sequence;
- a periodic sequence (e.g., a real or complex sinusoid); and
- the unit-step sequence (i.e., $u$ ).
- Fortunately, the Fourier transform can be extended to handle such sequences, resulting in what is known as the generalized Fourier transform.
- For our purposes, we can think of the classical and generalized Fourier transforms as being defined by the same formulas.
- Therefore, in what follows, we will not typically make a distinction between the classical and generalized Fourier transforms.


## DT Fourier Transform (DTFT)

- The Fourier transform of the sequence $x$, denoted $\mathcal{F} x$ or $X$, is given by

$$
\mathcal{F} x(\Omega)=X(\Omega)=\sum_{n=-\infty}^{\infty} x(n) e^{-j \Omega n}
$$

- The preceding equation is sometimes referred to as Fourier transform analysis equation (or forward Fourier transform equation).
- The inverse Fourier transform of $X$, denoted $\mathcal{F}^{-1} X$ or $x$, is given by

$$
\mathcal{F}^{-1} X(n)=x(n)=\frac{1}{2 \pi} \int_{2 \pi} X(\Omega) e^{j \Omega n} d \Omega
$$

- The preceding equation is sometimes referred to as the Fourier transform synthesis equation (or inverse Fourier transform equation).
- As a matter of notation, to denote that a sequence $x$ has the Fourier transform $X$, we write $x(n) \stackrel{\text { DTFTT }}{\longleftrightarrow} X(\Omega)$.
- A sequence $x$ and its Fourier transform $X$ constitute what is called a Fourier transform pair.


## Section 11.2

## Convergence Properties of the Fourier Transform

## Convergence of the Fourier Transform

- For a sequence $x$, the Fourier transform analysis equation (i.e., $\left.X(\Omega)=\sum_{-\infty}^{\infty} x(n) e^{-j \Omega n}\right)$ converges uniformly if

$$
\sum_{k=-\infty}^{\infty}|x(k)|<\infty
$$

(i.e., $x$ is absolutely summable).

- For a sequence $x$, the Fourier transform analysis equation (i.e., $\left.X(\Omega)=\sum_{-\infty}^{\infty} x(n) e^{-j \Omega n}\right)$ converges in the MSE sense if

$$
\sum_{k=-\infty}^{\infty}|x(k)|^{2}<\infty
$$

(i.e., $x$ is square summable).

- For a bounded Fourier transform $X$, the Fourier transform synthesis equation (i.e., $x(n)=\frac{1}{2 \pi} \int_{2 \pi} X(\Omega) e^{j \Omega n} d \Omega$ ) will always converge, since the integration interval is finite.


## Section 11.3

## Properties of the Fourier Transform

| Property | Time Domain | Frequency Domain |
| :--- | :--- | :--- |
| Linearity | $a_{1} x_{1}(n)+a_{2} x_{2}(n)$ | $a_{1} X_{1}(\Omega)+a_{2} X_{2}(\Omega)$ |
| Translation | $x\left(n-n_{0}\right)$ | $e^{-j \Omega n_{0}} X(\Omega)$ |
| Modulation | $e^{j \Omega_{0} n} x(n)$ | $X\left(\Omega-\Omega_{0}\right)$ |
| Conjugation | $x^{*}(n)$ | $X^{*}(-\Omega)$ |
| Time Reversal | $x(-n)$ | $X(-\Omega)$ |
| Upsampling | $(\uparrow M) x(n)$ | $X(M \Omega)$ |
| Downsampling | $(\downarrow M) x(n)$ | $\frac{1}{M} \sum_{k=0}^{M-1} X\left(\frac{\Omega-2 \pi k}{M}\right)$ |
| Convolution | $x_{1} * x_{2}(n)$ | $X_{1}(\Omega) X_{2}(\Omega)$ |
| Multiplication | $x_{1}(n) x_{2}(n)$ | $\frac{1}{2 \pi} \int_{2 \pi} X_{1}(\theta) X_{2}(\Omega-\theta) d \theta$ |
| Freq.-Domain Diff. | $n x(n)$ | $j \frac{d}{d \Omega} X(\Omega)$ |
| Differencing | $x(n)-x(n-1)$ | $\left(1-e^{-j \Omega}\right) X(\Omega)$ |
| Accumulation | $\sum_{k=-\infty}^{n} x(k)$ | $\frac{e^{j \Omega}}{e^{j \Omega}-1} X(\Omega)+\pi X(0) \sum_{k=-\infty}^{\infty} \delta(\Omega-2 \pi k)$ |

Property
Periodicity $\quad X(\Omega)=X(\Omega+2 \pi)$
Parseval's Relation
$\sum_{n=-\infty}^{\infty}|x(n)|^{2}=\frac{1}{2 \pi} \int_{2 \pi}|X(\Omega)|^{2} d \Omega$
Even Symmetry
$x$ is even $\Leftrightarrow X$ is even
Odd Symmetry
$x$ is odd $\Leftrightarrow X$ is odd
Real / Conjugate Symmetry $\quad x$ is real $\Leftrightarrow X$ is conjugate symmetric

## (DT) Fourier Transform Pairs

| Pair | $x(n)$ | $X(\Omega)$ |
| :---: | :---: | :---: |
| 1 | $\delta(n)$ | 1 |
| 2 | 1 | $2 \pi \sum_{k=-\infty}^{\infty} \delta(\Omega-2 \pi k)$ |
| 3 | $u(n)$ | $\frac{e^{j \Omega}}{e^{j \Omega}-1}+\sum_{k=-\infty}^{\infty} \pi \delta(\Omega-2 \pi k)$ |
| 4 | $a^{n} u(n),\|a\|<1$ | $\frac{e^{j \Omega}}{e^{j \Omega}-a}$ |
| 5 | $-a^{n} u(-n-1),\|a\|>1$ | $\frac{e^{j \Omega}}{e^{\text {S }}-a}$ |
| 6 | $a^{\|n\|},\|a\|<1$ | $\frac{1-a^{2}}{1-2 a \cos \Omega+a^{2}}$ |
| 7 | $\cos \left(\Omega_{0} n\right)$ | $\pi \sum_{k=-\infty}^{\infty}\left[\delta\left(\Omega-\Omega_{0}-2 \pi k\right)+\delta\left(\Omega+\Omega_{0}-2 \pi k\right)\right]$ |
| 8 | $\sin \left(\Omega_{0} n\right)$ | $j \pi \sum_{k=-\infty}^{\infty}\left[\delta\left(\Omega+\Omega_{0}-2 \pi k\right)-\delta\left(\Omega-\Omega_{0}-2 \pi k\right)\right]$ |
| 9 | $\cos \left(\Omega_{0} n\right) u(n)$ | $\frac{e^{\frac{j}{22 \Omega}-e^{j 2} \cos \Omega_{0}}}{e^{j 2 \Omega}-2 e^{j 2} \cos \Omega_{0}+1}+\frac{\pi}{2} \sum_{k=-\infty}^{\infty}\left[\delta\left(\Omega-2 \pi k-\Omega_{0}\right)+\delta\left(\Omega-2 \pi k+\Omega_{0}\right)\right]$ |
| 10 | $\sin \left(\Omega_{0} n\right) u(n)$ | $\frac{e^{j \Omega} \sin \Omega_{0}}{e^{j 2 \Omega}-2 e^{j 2} \cos \Omega_{0}+1}+\frac{\pi}{2 j} \sum_{k=-\infty}^{\infty}\left[\delta\left(\Omega-2 \pi k-\Omega_{0}\right)-\delta\left(\Omega-2 \pi k+\Omega_{0}\right)\right]$ |
| 11 | $\frac{B}{\pi} \operatorname{sinc}(B n), 0<B<\pi$ | $\sum_{k=-\infty}^{\infty} \operatorname{rect}\left(\frac{\Omega-2 \pi k}{2 B}\right)$ |
| 12 | $u(n)-u(n-M)$ | $e^{-j \Omega(M-1) / 2}\left(\frac{\sin (M \Omega / 2)}{\sin (\Omega / 2)}\right)$ |
| 13 | $n a^{n} u(n),\|a\|<1$ | $\frac{a e^{j \Omega}}{\left(e^{j \Omega}-a\right)^{2}}$ |

- Recall the definition of the Fourier transform $X$ of the sequence $x$ :

$$
X(\Omega)=\sum_{n=-\infty}^{\infty} x(n) e^{-j \Omega n}
$$

- For all integer $k$, we have that

$$
\begin{aligned}
X(\Omega+2 \pi k) & =\sum_{n=-\infty}^{\infty} x(n) e^{-j(\Omega+2 \pi k) n} \\
& =\sum_{n=-\infty}^{\infty} x(n) e^{-j(\Omega n+2 \pi k n)} \\
& =\sum_{n=-\infty}^{\infty} x(n) e^{-j \Omega n} \\
& =X(\Omega)
\end{aligned}
$$

- Thus, the Fourier transform $X$ of the sequence $x$ is always $2 \pi$-periodic.

■ If $x_{1}(n) \stackrel{\text { DTFT }}{\longleftrightarrow} X_{1}(\Omega)$ and $x_{2}(n) \stackrel{\text { DTFT }}{\longleftrightarrow} X_{2}(\Omega)$, then

$$
a_{1} x_{1}(n)+a_{2} x_{2}(n) \stackrel{\text { DTFT }}{\longleftrightarrow} a_{1} X_{1}(\Omega)+a_{2} X_{2}(\Omega),
$$

where $a_{1}$ and $a_{2}$ are arbitrary complex constants.

- This is known as the linearity property of the Fourier transform.


## Translation

■ If $x(n) \stackrel{\text { DTFT }}{\longleftrightarrow} X(\Omega)$, then

$$
x\left(n-n_{0}\right) \stackrel{\text { DTFT }}{\longleftrightarrow} e^{-j \Omega n_{0}} X(\Omega),
$$

where $n_{0}$ is an arbitrary integer.

- This is known as the translation (or time-domain shifting) property of the Fourier transform.


## Modulation

■ If $x(n) \stackrel{\text { DTFT }}{\longleftrightarrow} X(\Omega)$, then

$$
e^{j \Omega_{0} n} x(n) \stackrel{\text { DTFT }}{\longleftrightarrow} X\left(\Omega-\Omega_{0}\right),
$$

where $\Omega_{0}$ is an arbitrary real constant.

- This is known as the modulation (or frequency-domain shifting) property of the Fourier transform.


## Conjugation

■ If $x(n) \stackrel{\text { DTFT }}{\longleftrightarrow} X(\Omega)$, then

$$
x^{*}(n) \stackrel{\text { DTFT }}{\longleftrightarrow} X^{*}(-\Omega) .
$$

- This is known as the conjugation property of the Fourier transform.

■ If $x(n) \stackrel{\text { DTFT }}{\longleftrightarrow} X(\Omega)$, then

$$
x(-n) \stackrel{\text { DTFT }}{\longleftrightarrow} X(-\Omega) .
$$

- This is known as the time-reversal property of the Fourier transform.


## Upsampling

■ If $x(n) \stackrel{\text { DTFT }}{\longleftrightarrow} X(\Omega)$, then

$$
(\uparrow M) x(n) \stackrel{\text { DTFT }}{\longleftrightarrow} X(M \Omega) .
$$

- This is known as the upsampling property of the Fourier transform.


## Downsampling

■ If $x(n) \stackrel{\text { DTFT }}{\longleftrightarrow} X(\Omega)$, then

$$
(\downarrow M) x(n) \stackrel{\text { DTFT }}{\longleftrightarrow} \frac{1}{M} \sum_{k=0}^{M-1} X\left(\frac{\Omega-2 \pi k}{M}\right) .
$$

■ This is known as the downsampling property of the Fourier transform.

## Convolution

■ If $x_{1}(n) \stackrel{\text { DTFT }}{\longleftrightarrow} X_{1}(\Omega)$ and $x_{2}(n) \stackrel{\text { DTFT }}{\longleftrightarrow} X_{2}(\Omega)$, then

$$
x_{1} * x_{2}(n) \stackrel{\text { DTFT }}{\longleftrightarrow} X_{1}(\Omega) X_{2}(\Omega) .
$$

- This is known as the convolution (or time-domain convolution) property of the Fourier transform.
- In other words, a convolution in the time domain becomes a multiplication in the frequency domain.
- This suggests that the Fourier transform can be used to avoid having to deal with convolution operations.

■ If $x_{1}(n) \stackrel{\text { DTFT }}{\longleftrightarrow} X_{1}(\Omega)$ and $x_{2}(n) \stackrel{\text { DTFT }}{\longleftrightarrow} X_{2}(\Omega)$, then

$$
x_{1}(n) x_{2}(n) \stackrel{\text { DTFT }}{\longleftrightarrow} \frac{1}{2 \pi} \int_{2 \pi} X_{1}(\theta) X_{2}(\Omega-\theta) d \theta .
$$

- This is known as the multiplication (or time-domain multiplication) property of the Fourier transform.
- Do not forget the factor of $\frac{1}{2 \pi}$ in the above formula!
- This property of the Fourier transform is often tedious to apply (in the forward direction) as it turns a multiplication into a convolution.

■ If $x(n) \stackrel{\text { DTFT }}{\longleftrightarrow} X(\Omega)$, then

$$
n x(n) \stackrel{\text { DTFT }}{\longleftrightarrow} j \frac{d}{d \Omega} X(\Omega) .
$$

- This is known as the frequency-domain differentiation property of the Fourier transform.


## Differencing

- If $x(n) \stackrel{\text { DTFT }}{\longleftrightarrow} X(\Omega)$, then

$$
x(n)-x(n-1) \stackrel{\text { DTFT }}{\longleftrightarrow}\left(1-e^{-j \Omega}\right) X(\Omega) .
$$

- This is known as the differencing property of the Fourier transform.
- Note that this property follows quite trivially from the linearity and translation properties of the Fourier transform.


## Accumulation

- If $x(n) \xrightarrow{\text { Dret }} X(\Omega)$, then

$$
\sum_{k=-\infty}^{n} x(k) \stackrel{\text {,OTFT }}{\longleftrightarrow} \frac{e^{j \Omega}}{e^{j \Omega}-1} X(\Omega)+\pi X(0) \sum_{k=-\infty}^{\infty} \delta(\Omega-2 \pi k) .
$$

- This is known as the accumulation property of the Fourier transform.

■ If $x(n) \stackrel{\text { DTFT }}{\longleftrightarrow} X(\Omega)$, then

$$
\sum_{n=-\infty}^{\infty}|x(n)|^{2}=\frac{1}{2 \pi} \int_{2 \pi}|X(\Omega)|^{2} d \Omega
$$

(i.e., the energy of $x$ and energy of $X$ are equal up to a factor of $2 \pi$ ).

■ This is known as Parseval's relation.

- Since energy is often a quantity of great significance in engineering applications, it is extremely helpful to know that the Fourier transform preserves energy (up to a scale factor).
- For a sequence $x$ with Fourier transform $X$, the following assertions hold:
$11 x$ is even $\Leftrightarrow X$ is even; and
[2 $x$ is odd $\Leftrightarrow X$ is odd.
- In other words, the forward and inverse Fourier transforms preserve even/odd symmetry.
- A sequence $x$ is real if and only if its Fourier transform $X$ satisfies

$$
X(\Omega)=X^{*}(-\Omega) \text { for all } \Omega
$$

(i.e., $X$ is conjugate symmetric).

- Thus, for a real-valued sequence, the portion of the graph of a Fourier transform for negative values of frequency $\Omega$ is redundant, as it is completely determined by symmetry.
- From properties of complex numbers, one can show that $X(\Omega)=X^{*}(-\Omega)$ is equivalent to

$$
|X(\Omega)|=|X(-\Omega)| \quad \text { and } \quad \arg X(\Omega)=-\arg X(-\Omega)
$$

(i.e., $|X(\Omega)|$ is even and $\arg X(\Omega)$ is odd).

- Note that $x$ being real does not necessarily imply that $X$ is real.


## Section 11.4

## Fourier Transform of Periodic Sequences

## Fourier Transform of Periodic Sequences

- The Fourier transform can be generalized to also handle periodic sequences.
- Consider an $N$-periodic sequence $x$.
- Define the sequence $x_{N}$ as

$$
x_{N}(n)= \begin{cases}x(n) & 0 \leq n<N \\ 0 & \text { otherwise }\end{cases}
$$

(i.e., $x_{N}(n)$ is equal to $x(n)$ over a single period and zero elsewhere).

- Let $a$ denote the Fourier series coefficient sequence of $x$.
- Let $X$ and $X_{N}$ denote the Fourier transforms of $x$ and $x_{N}$, respectively.
- The following relationships can be shown to hold:

$$
\begin{gathered}
X(\Omega)=\frac{2 \pi}{N} \sum_{k=-\infty}^{\infty} X_{N}\left(\frac{2 \pi k}{N}\right) \delta\left(\Omega-\frac{2 \pi k}{N}\right), \\
a_{k}=\frac{1}{N} X_{N}\left(\frac{2 \pi k}{N}\right), \quad \text { and } \quad X(\Omega)=2 \pi \sum_{k=-\infty}^{\infty} a_{k} \delta\left(\Omega-\frac{2 \pi k}{N}\right) .
\end{gathered}
$$

- The Fourier series coefficient sequence $a$ is produced by sampling $X_{N}$ at integer multiples of the fundamental frequency $\frac{2 \pi}{N}$ and scaling the resulting sequence by $\frac{1}{N}$.
- The Fourier transform of a periodic sequence can only be nonzero at integer multiples of the fundamental frequency.


## Section 11.5

## Fourier Transform and Frequency Spectra of Sequences

- Like Fourier series, the Fourier transform also provides us with a frequency-domain perspective on sequences.
- That is, instead of viewing a sequence as having information distributed with respect to time (i.e., a function whose domain is time), we view a sequence as having information distributed with respect to frequency (i.e., a function whose domain is frequency).
- The Fourier transform $X$ of a sequence $x$ provides a means to quantify how much information $x$ has at different frequencies.
- The distribution of information in a sequence over different frequencies is referred to as the frequency spectrum of the sequence.
- To gain further insight into the role played by the Fourier transform $X$ in the context of the frequency spectrum of $x$, it is helpful to write the Fourier transform representation of $x$ with $X(\Omega)$ expressed in polar form as follows:

$$
x(n)=\frac{1}{2 \pi} \int_{2 \pi} X(\Omega) e^{j \Omega n} d \Omega=\frac{1}{2 \pi} \int_{2 \pi}|X(\Omega)| e^{j[\Omega n+\arg X(\Omega)]} d \Omega
$$

- In effect, the quantity $|X(\Omega)|$ is a weight that determines how much the complex sinusoid at frequency $\Omega$ contributes to the integration result $x(n)$.
- Perhaps, this can be more easily seen if we express the above integral as the limit of a sum, derived from an approximation of the integral using the area of rectangles, as shown on the next slide. [Recall that $\int_{a}^{b} f(x) d x=\lim _{n \rightarrow \infty} \sum_{k=1}^{n} f\left(x_{k}\right) \Delta x$ where $\Delta x=\frac{b-a}{n}$ and $x_{k}=a+k \Delta x$.]


## Fourier Transform and Frequency Spectra (Continued 1)

- Expressing the integral (from the previous slide) as the limit of a sum, we obtain

$$
x(n)=\lim _{\ell \rightarrow \infty} \frac{1}{2 \pi} \sum_{k=1}^{\ell} \Delta \Omega\left|X\left(\Omega^{\prime}\right)\right| e^{j\left[\Omega^{\prime} n+\arg X\left(\Omega^{\prime}\right)\right]}
$$

where $\Delta \Omega=\frac{2 \pi}{\ell}$ and $\Omega^{\prime}=k \Delta \Omega$.

- In the above equation, the $k$ th term in the summation corresponds to a complex sinusoid with fundamental frequency $\Omega^{\prime}=k \Delta \Omega$ that has had its amplitude scaled by a factor of $\left|X\left(\Omega^{\prime}\right)\right|$ and has been time shifted by an amount that depends on $\arg X\left(\Omega^{\prime}\right)$.
- For a given $\Omega^{\prime}=k \Delta \Omega$ (which is associated with the $k$ th term in the summation), the larger $\left|X\left(\Omega^{\prime}\right)\right|$ is, the larger the amplitude of its corresponding complex sinusoid $e^{j \Omega^{\prime} n}$ will be, and therefore the larger the contribution the $k$ th term will make to the overall summation.
- In this way, we can use $\left|X\left(\Omega^{\prime}\right)\right|$ as a measure of how much information a sequence $x$ has at the frequency $\Omega^{\prime}$.
- The Fourier transform $X$ of the sequence $x$ is referred to as the frequency spectrum of $x$.
- The magnitude $|X(\Omega)|$ of the Fourier transform $X$ is referred to as the magnitude spectrum of $x$.
- The argument $\arg X(\Omega)$ of the Fourier transform $X$ is referred to as the phase spectrum of $x$.
- Since the Fourier transform is a function of a real variable, a sequence can potentially have information at any real frequency.
- Earlier, we saw that for periodic sequences, the Fourier transform can only be nonzero at integer multiples of the fundamental frequency.
- So, the Fourier transform and Fourier series give a consistent picture in terms of frequency spectra.
- Since the frequency spectrum is complex (in the general case), it is usually represented using two plots, one showing the magnitude spectrum and one showing the phase spectrum.
- Recall that, for a real sequence $x$, the Fourier transform $X$ of $x$ satisfies

$$
X(\Omega)=X^{*}(-\Omega)
$$

(i.e., $X$ is conjugate symmetric), which is equivalent to

$$
|X(\Omega)|=|X(-\Omega)| \quad \text { and } \quad \arg X(\Omega)=-\arg X(-\Omega) .
$$

■ Since $|X(\Omega)|=|X(-\Omega)|$, the magnitude spectrum of a real sequence is always even.

- Similarly, since $\arg X(\Omega)=-\arg X(-\Omega)$, the phase spectrum of a real sequence is always odd.
- Due to the symmetry in the frequency spectra of real sequences, we typically ignore negative frequencies when dealing with such sequences.
- In the case of sequences that are complex but not real, frequency spectra do not possess the above symmetry, and negative frequencies become important.


## Bandwidth

- A sequence $x$ with Fourier transform $X$ satisfying $X(\Omega)=0$ for all $\Omega$ in $(-\pi, \pi]$ except for some interval $I$ is said to be bandllimited to frequencies in $I$.
- The bandwidth of a sequence $x$ with Fourier transform $X$ is the length of the interval in $(-\pi, \pi]$ over which $X$ is nonzero.
- For example, the sequence $x$ with the Fourier transform $X$ shown below is bandlimited to frequencies in $[-B, B]$ and has bandwidth $B-(-B)=2 B$.

- Since $x$ is real in the above example (as $X$ is conjugate symmetric), we might choose to ignore negative frequencies, in which case $x$ would be deemed to be bandlimited to frequencies in $[0, B]$ and have bandwidth $B-0=B$.


## Energy-Density Spectra

- By Parseval's relation, the energy $E$ in a sequence $x$ with Fourier transform $X$ is given by

$$
E=\frac{1}{2 \pi} \int_{2 \pi} E_{x}(\Omega) d \Omega
$$

where

$$
E_{x}(\Omega)=|X(\Omega)|^{2}
$$

- We refer to $E_{x}$ as the energy-density spectrum of the sequence $x$.
- The function $E_{x}$ indicates how the energy in $x$ is distributed with respect to frequency.
- For example, the energy contributed by frequencies in the range $\left[\Omega_{1}, \Omega_{2}\right]$ is given by

$$
\frac{1}{2 \pi} \int_{\Omega_{1}}^{\Omega_{2}} E_{x}(\Omega) d \Omega
$$

## Section 11.6

## Fourier Transform and LTI Systems

- Consider a LTI system with input $x$, output $y$, and impulse response $h$, and let $X, Y$, and $H$ denote the Fourier transforms of $x, y$, and $h$, respectively.
- Since $y(n)=x * h(n)$, we have that

$$
Y(\Omega)=X(\Omega) H(\Omega)
$$

- The function $H$ is called the frequency response of the system.
- A LTI system is completely characterized by its frequency response $H$.
- The above equation provides an alternative way of viewing the behavior of a LTI system. That is, we can view the system as operating in the frequency domain on the Fourier transforms of the input and output signals.
- The frequency spectrum of the output is the product of the frequency spectrum of the input and the frequency response of the system.
- In the general case, the frequency response $H$ is a complex-valued function.
- Often, we represent $H(\Omega)$ in terms of its magnitude $|H(\Omega)|$ and argument $\arg H(\Omega)$.
- The quantity $|H(\Omega)|$ is called the magnitude response of the system.
- The quantity $\arg H(\Omega)$ is called the phase response of the system.
- Since $Y(\Omega)=X(\Omega) H(\Omega)$, we trivially have that

$$
|Y(\Omega)|=|X(\Omega)||H(\Omega)| \quad \text { and } \quad \arg Y(\Omega)=\arg X(\Omega)+\arg H(\Omega)
$$

- The magnitude spectrum of the output equals the magnitude spectrum of the input times the magnitude response of the system.
- The phase spectrum of the output equals the phase spectrum of the input plus the phase response of the system.
- Since the frequency response $H$ is simply the frequency spectrum of the impulse response $h$, if $h$ is real, then

$$
|H(\Omega)|=|H(-\Omega)| \quad \text { and } \quad \arg H(\Omega)=-\arg H(-\Omega)
$$

(i.e., the magnitude response $|H(\Omega)|$ is even and the phase response $\arg H(\Omega)$ is $o d d)$.

## Unwrapped Phase

- For many types of analysis, restricting the range of a phase function to an interval of length $2 \pi$ (such as $(-\pi, \pi]$ ), often unnecessarily introduces discontinuities into the function.
- This motivates the notion of unwrapped phase.
- The unwrapped phase is simply the phase defined in such a way so as not to restrict the phase to an interval of length $2 \pi$ and to keep the phase function continuous to the greatest extent possible.
- For example, the function $H(\Omega)=e^{j 3 \Omega}$ has the unwrapped phase $\Theta(\Omega)=3 \Omega$.


Phase


Unwrapped Phase

## Interpretation of Magnitude and Phase Response

■ Recall that a LTI system $\mathcal{H}$ with frequency response $H$ is such that

$$
\mathcal{H}\left\{e^{j \Omega n}\right\}(n)=H(\Omega) e^{j \Omega n} .
$$

■ Expressing $H(\Omega)$ in polar form, we have

$$
\begin{aligned}
\mathcal{H}\left\{e^{j \Omega n}\right\}(n) & =|H(\Omega)| e^{j \arg H(\Omega)} e^{j \Omega n} \\
& =|H(\Omega)| e^{j[\Omega n+\arg H(\Omega)]} \\
& =|H(\Omega)| e^{j \Omega(n+\arg [H(\Omega)] / \Omega)} .
\end{aligned}
$$

- Thus, the response of the system to the sequence $e^{j \Omega n}$ is produced by applying two transformations to this sequence:
- (amplitude) scaling by $|H(\Omega)|$; and
- translating by $-\frac{\arg H(\Omega)}{\Omega}$ (using bandlimited interpolation if $-\frac{\arg H(\Omega)}{\Omega} \notin \mathbb{Z}$ ).
- Therefore, the magnitude response determines how different complex sinusoids are scaled (in amplitude) by the system.
- Similarly, the phase response determines how different complex sinusoids are translated (i.e., delayed/advanced) by the system.
- Recall that a LTI system $\mathcal{H}$ with frequency response $H$ is such that

$$
\mathcal{H}\left\{e^{j \Omega n}\right\}(n)=|H(\Omega)| e^{j \Omega(n+\arg [H(\Omega)] / \Omega)} .
$$

- If $|H(\Omega)|$ is a constant (for all $\Omega$ ), every complex sinusoid is scaled by the same amount when passing through the system.
- A system for which $|H(\Omega)|=1$ (for all $\Omega$ ) is said to be allpass.
- In the case of an allpass system, the magnitude spectra of the system's input and output are identical.
- If $|H(\Omega)|$ is not a constant, different complex sinusoids are scaled by different amounts, resulting in what is known as magnitude distortion.

■ Recall that a LTI system $\mathcal{H}$ with frequency response $H$ is such that

$$
\mathcal{H}\left\{e^{j \Omega n}\right\}(n)=|H(\Omega)| e^{j \Omega(n+\arg [H(\Omega)] / \Omega)}
$$

- The preceding equation can be rewritten as

$$
\mathcal{H}\left\{e^{j \Omega n}\right\}(n)=|H(\Omega)| e^{j \Omega\left[n-\tau_{\mathrm{p}}(\Omega)\right]} \quad \text { where } \quad \tau_{\mathrm{p}}(\Omega)=-\frac{\arg H(\Omega)}{\Omega} .
$$

- The function $\tau_{\mathrm{p}}$ is known as the phase delay of the system.
- If $\tau_{\mathrm{p}}(\Omega)=n_{\mathrm{d}}$ (where $n_{\mathrm{d}}$ is a constant), the system shifts all complex sinusoids by the same amount $n_{\mathrm{d}}$.
- Since $\tau_{\mathrm{p}}(\Omega)=n_{\mathrm{d}}$ is equivalent to the (unwrapped) phase response being of the form $\arg H(\Omega)=-n_{\mathrm{d}} \Omega$ (which is a linear function with a zero constant term), a system with a constant phase delay is said to have linear phase.
- In the case that $\tau_{\mathrm{p}}(\Omega)=0$, the system is said to have zero phase.
- If $\tau_{\mathrm{p}}(\Omega)$ is not a constant, different complex sinusoids are shifted by different amounts, resulting in what is known as phase distortion.


## Distortionless Transmission

- Consider a LTI system $\mathcal{H}$ with input $x$ and output $y$ given by

$$
y(n)=x\left(n-n_{0}\right),
$$

where $n_{0}$ is an integer constant.

- That is, the output of the system is simply the input delayed by $n_{0}$.
- This type of behavior is the ideal for which we strive in real-world communication systems (i.e., the received signal $y$ equals a delayed version of the transmitted signal $x$ ).
- Taking the Fourier transform of the preceding equation, we have

$$
Y(\Omega)=e^{-j \Omega n_{0}} X(\Omega)
$$

- Thus, the system has the frequency response $H$ given by

$$
H(\Omega)=e^{-j \Omega n_{0}}
$$

- Since the phase delay of the system is $\tau_{\mathrm{p}}(\Omega)=-\left(\frac{-\Omega n_{0}}{\Omega}\right)=n_{0}$, the phase delay is constant and the system has linear phase.
- Consider a LTI system with input $x$, output $y$, and impulse response $h$, and let $X, Y$, and $H$ denote the Fourier transforms of $x, y$, and $h$, respectively.
- Often, it is convenient to represent such a system in block diagram form in the frequency domain as shown below.

- Since a LTI system is completely characterized by its frequency response, we typically label the system with this quantity.


## Interconnection of LTI Systems

- The series interconnection of the LTI systems with frequency responses $H_{1}$ and $H_{2}$ is the LTI system with frequency response $H_{1} H_{2}$. That is, we have the equivalences shown below.

- The parallel interconnection of the LTI systems with frequency responses $H_{1}$ and $H_{2}$ is the LTI system with the frequency response $H_{1}+H_{2}$. That is, we have the equivalence shown below.



## LTI Systems and Difference Equations

■ Many LTI systems of practical interest can be represented using an Nth-order linear difference equation with constant coefficients.

- Consider a system with input $x$ and output $y$ that is characterized by an equation of the form

$$
\sum_{k=0}^{N} b_{k} y(n-k)=\sum_{k=0}^{M} a_{k} x(n-k)
$$

- Let $h$ denote the impulse response of the system, and let $X, Y$, and $H$ denote the Fourier transforms of $x, y$, and $h$, respectively.
- One can show that $H(\Omega)$ is given by

$$
H(\Omega)=\frac{Y(\Omega)}{X(\Omega)}=\frac{\sum_{k=0}^{M} a_{k}\left(e^{j \Omega}\right)^{-k}}{\sum_{k=0}^{N} b_{k}\left(e^{j \Omega}\right)^{-k}}=\frac{\sum_{k=0}^{M} a_{k} e^{-j k \Omega}}{\sum_{k=0}^{N} b_{k} e^{-j k \Omega}}
$$

- Each of the numerator and denominator of $H$ is a polynomial in $e^{-j \Omega}$.
- Thus, $H$ is a rational function in the variable $e^{-j \Omega}$.


## Section 11.7

## Fourier Transform Relationships

## Duality Between DTFT and CTFS

- The DTFT analysis and synthesis equations are, respectively, given by

$$
X(\Omega)=\sum_{k=-\infty}^{\infty} x(k) e^{-j k \Omega} \quad \text { and } \quad x(n)=\frac{1}{2 \pi} \int_{2 \pi} X(\Omega) e^{j n \Omega} d \Omega
$$

- The CTFS synthesis and analysis equations are, respectively, given by

$$
x_{\mathrm{c}}(t)=\sum_{k=-\infty}^{\infty} a(k) e^{j k(2 \pi / T) t} \quad \text { and } \quad a(n)=\frac{1}{T} \int_{T} x_{\mathrm{c}}(t) e^{-j n(2 \pi / T) t} d t
$$

which can be rewritten, respectively, as

$$
x_{\mathrm{c}}(t)=\sum_{k=-\infty}^{\infty} a(-k) e^{-j k(2 \pi / T) t} \quad \text { and } \quad a(-n)=\frac{1}{T} \int_{T} x_{\mathrm{c}}(t) e^{j n(2 \pi / T) t} d t
$$

- The CTFS synthesis equation with $T=2 \pi$ corresponds to the DTFT analysis equation with $X=x_{\mathrm{c}}, \Omega=t$, and $x(n)=a(-n)$.
- The CTFS analysis equation with $T=2 \pi$ corresponds to the DTFT synthesis equation with $X=x_{\mathrm{c}}$ and $x(n)=a(-n)$.
- Consequently, the DTFT $X$ of the sequence $x$ can be viewed as a CTFS representation of the $2 \pi$-periodic spectrum $X$.
- Let $x$ be a bandlimited function and let $T$ denote a sampling period for $x$ that satisfies the Nyquist condition.
- Let $\tilde{y}$ be the function obtained by impulse sampling $x$ with sampling period $T$. That is,

$$
\tilde{y}(t)=\sum_{n=-\infty}^{\infty} x(T n) \delta(t-T n) .
$$

■ Let $y$ denote the sequence obtaining by sampling $x$ with sampling period $T$. That is,

$$
y(n)=x(T n) .
$$

- Let $\tilde{Y}$ denote the (CT) Fourier transform of $\tilde{y}$ and let $Y$ denote the (DT) Fourier transform of $y$.
- Then, the following relationship holds:

$$
Y(\Omega)=\tilde{Y}\left(\frac{\Omega}{T}\right) \quad \text { for all } \Omega \in \mathbb{R}
$$

- Let $x$ be a sequence with (DT) Fourier transform $X$ such that

$$
x(n)=0 \quad \text { for all } n \notin[0 \ldots M-1] .
$$

- Let $\tilde{X}$ denote the $N$-point DFT of $X$. That is,

$$
\tilde{X}(k)=\sum_{n=0}^{N-1} x(n) e^{-j(2 \pi / N) k n} \quad \text { for } k \in[0 \ldots N-1] .
$$

- Suppose now that $N \geq M$.
- Then, the following relationship holds:

$$
X\left(\frac{2 \pi}{N} k\right)=\tilde{X}(k) \quad \text { for } k \in[0 \ldots N-1] .
$$

■ In other words, the elements of the sequence $\tilde{X}$ correspond to uniformly-spaced samples of the function $X$.

- Consider the sequence

$$
x(n)=u(n)-u(n-4)
$$

■ The Fourier transform $X$ of $x$ can be shown to be

$$
X(\Omega)=e^{-j(3 / 2) \Omega}\left[\frac{\sin (2 \Omega)}{\sin \left(\frac{1}{2} \Omega\right)}\right]
$$

- Clearly, $x(n)=0$ for all $n \notin[0 \ldots 3]$.
- Therefore, uniformly-spaced samples of $X$ can be obtained from an $N$-point DFT $\tilde{X}$ of $x$, where $N \geq 4$.
- The subsequent slides show the sampled spectrum obtained by the DFT for several values of $N$.


## Spectral Sampling Example: $N=4$




Phase Spectrum

## Spectral Sampling Example: $N=8$



Phase Spectrum

## Spectral Sampling Example: $N=16$



Phase Spectrum

## Spectral Sampling Example: $N=64$




## Section 11.8

## Application: Filtering

- In many applications, we want to modify the spectrum of a signal by either amplifying or attenuating certain frequency components.
- This process of modifying the frequency spectrum of a signal is called filltering.
- A system that performs a filtering operation is called a filter.
- Many types of filters exist.
- Frequency selective filters pass some frequencies with little or no distortion, while significantly attenuating other frequencies.
- Several basic types of frequency-selective filters include: lowpass, highpass, and bandpass.


## Ideal Lowpass Filter

- An ideal lowpass filter eliminates all baseband frequency components with a frequency whose magnitude is greater than some cutoff frequency, while leaving the remaining baseband frequency components unaffected.
- Such a filter has a frequency response $H$ of the form

$$
H(\Omega)= \begin{cases}1 & |\Omega| \leq \Omega_{c} \\ 0 & \Omega_{c}<|\Omega| \leq \pi\end{cases}
$$

where $\Omega_{c}$ is the cutoff frequency.

- A plot of this frequency response is given below.



## Ideal Highpass Filter

- An ideal highpass filter eliminates all baseband frequency components with a frequency whose magnitude is less than some cutoff frequency, while leaving the remaining baseband frequency components unaffected.
- Such a filter has a frequency response $H$ of the form

$$
H(\Omega)= \begin{cases}1 & \Omega_{c}<|\Omega| \leq \pi \\ 0 & |\Omega| \leq \Omega_{c}\end{cases}
$$

where $\Omega_{c}$ is the cutoff frequency.

- A plot of this frequency response is given below.



## Ideal Bandpass Filter

- An ideal lbandpass filter eliminates all baseband frequency components with a frequency whose magnitude does not lie in a particular range, while leaving the remaining baseband frequency components unaffected.
- Such a filter has a frequency response $H$ of the form

$$
H(\Omega)= \begin{cases}1 & \Omega_{c 1} \leq|\Omega| \leq \Omega_{c 2} \\ 0 & |\Omega|<\Omega_{c 1} \text { or } \Omega_{c 2}<|\Omega|<\pi\end{cases}
$$

where the limits of the passband are $\Omega_{c 1}$ and $\Omega_{c 2}$.

- A plot of this frequency response is given below.



## Part 12

## z Transform (ZT)

- Another important mathematical tool in the study of signals and systems is known as the $z$ transform.
- The $z$ transform can be viewed as a generalization of the (classical) Fourier transform.
- Due to its more general nature, the $z$ transform has a number of advantages over the (classical) Fourier transform.
- First, the z transform representation exists for some sequences that do not have a Fourier transform representation. So, we can handle some sequences with the $z$ transform that cannot be handled with the Fourier transform.
- Second, since the $z$ transform is a more general tool, it can provide additional insights beyond those facilitated by the Fourier transform.
- Earlier, we saw that complex exponentials are eigensequences of LTI systems.
- In particular, for a LTI system $\mathcal{H}$ with impulse response $h$, we have that

$$
\mathcal{H}\left\{z^{n}\right\}(n)=H(z) z^{n} \quad \text { where } \quad H(z)=\sum_{n=-\infty}^{\infty} h(n) z^{-n} .
$$

- Previously, we referred to $H$ as the system function.
- As it turns out, $H$ is the $z$ transform of $h$.
- Since the $z$ transform has already appeared earlier in the context of LTI systems, it is clearly a useful tool.
- Furthermore, as we will see, the z transform has many additional uses.


## Section 12.1

## z Transform

## (Bilateral) z Transform

- The (bilateral) z transform of the sequence $x$, denoted $Z x$ or $X$, is defined as

$$
z x(z)=X(z)=\sum_{n=-\infty}^{\infty} x(n) z^{-n}
$$

- The inverse z transform of $X$, denoted $z^{-1} X$ or $x$, is then given by

$$
z^{-1} X(n)=x(n)=\frac{1}{2 \pi j} \oint_{\Gamma} X(z) z^{n-1} d z
$$

where $\Gamma$ is a counterclockwise closed circular contour centered at the origin and with radius $r$ such that $\Gamma$ is in the ROC of $X$.

- We refer to $x$ and $X$ as a z transform pair and denote this relationship as

$$
x(n) \stackrel{\mathrm{ZT}}{\longleftrightarrow} X(z) .
$$

- In practice, we do not usually compute the inverse $z$ transform by directly using the formula from above. Instead, we resort to other means (to be discussed later).

■ Two different versions of the $z$ transform are commonly used:
1 the bilateral (or two-sided) z transform; and
2 the unilateral (or one-sided) z transform.

- The unilateral $z$ transform is most frequently used to solve systems of linear difference equations with nonzero initial conditions.
- As it turns out, the only difference between the definitions of the bilateral and unilateral $z$ transforms is in the lower limit of summation.
- In the bilateral case, the lower limit is $-\infty$, whereas in the unilateral case, the lower limit is 0 .
- For the most part, we will focus our attention primarily on the bilateral $z$ transform.
- We will, however, briefly introduce the unilateral z transform as a tool for solving difference equations.
- Unless otherwise noted, all subsequent references to the $z$ transform should be understood to mean bilateral z transform.


## Relationship Between Z and Fourier Transforms

- Let $X$ and $X_{\mathrm{F}}$ denote the $z$ and (DT) Fourier transforms of $x$, respectively.
- The function $X(z)$ evaluated at $z=e^{j \Omega}$ (where $\Omega$ is real) yields $X_{\mathrm{F}}(\Omega)$. That is,

$$
X\left(e^{j \Omega}\right)=X_{\mathrm{F}}(\Omega)
$$

- Due to the preceding relationship, the Fourier transform of $x$ is sometimes written as $X\left(e^{j \Omega}\right)$.
- The function $X(z)$ evaluated at an arbitrary complex value $z=r e^{j \Omega}$ (where $r=|z|$ and $\Omega=\arg z$ ) can also be expressed in terms of a Fourier transform involving $x$. In particular, we have

$$
X\left(r e^{j \Omega}\right)=X_{\mathrm{F}}^{\prime}(\Omega),
$$

where $X_{\mathrm{F}}^{\prime}$ is the (DT) Fourier transform of $x^{\prime}(n)=r^{-n} x(n)$.

- So, in general, the $z$ transform of $x$ is the Fourier transform of an exponentially-weighted version of $x$.
- Due to this weighting, the $z$ transform of a sequence may exist when the Fourier transform of the same sequence does not.


## z Transform Examples

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## Section 12.2

## Region of Convergence (ROC)

- A disk with center 0 and radius $r$ is the set of all complex numbers $z$ satisfying

$$
|z|<r
$$

where $r$ is a real constant and $r>0$.


## Annulus

- An annullus with center 0 , inner radius $r_{0}$, and outer radius $r_{1}$ is the set of all complex numbers $z$ satisfying

$$
r_{0}<|z|<r_{1}
$$

where $r_{0}$ and $r_{1}$ are real constants and $0<r_{0}<r_{1}$.


## Circle Exterior

- The exterior of a circle with center 0 and radius $r$ is the set of all complex numbers $z$ satisfying

$$
|z|>r,
$$

where $r$ is a real constant and $r>0$.


## Example: Set Intersection


$R_{1}$

$R_{2}$

$R_{1} \cap R_{2}$

## Example: Scalar Multiple of a Set




## Example: Reciprocal of a Set




- As we saw earlier, for a sequence $x$, the complete specification of its $z$ transform $X$ requires not only an algebraic expression for $X$, but also the ROC associated with $X$.
- Two very different sequences can have the same algebraic expressions for $X$.
- Now, we examine some of the constraints on the ROC (of the $z$ transform) for various classes of sequences.

■ The ROC of a z transform consists of concentric circles centered at 0 in the complex plane.

- That is, if a point $z_{0}$ is in the ROC, then the circle centered at 0 passing through $z_{0}$ (i.e., $|z|=\left|z_{0}\right|$ ) is also in the ROC.
■ Some examples of sets that would be either valid or invalid as ROCs are shown below.


Valid


Valid


Invalid

- If a z transform $X$ is a rational function, then the ROC of $X$ does not contain any poles and is bounded by poles or extends to infinity.
- Some examples of sets that would be either valid or invalid as ROCs of rational $z$ transforms are shown below.


Valid


Valid


Invalid

- If a sequence $x$ is finite duration and its $z$ transform $X$ converges for at least one point, then $X$ converges for all points the complex plane, except possibly 0 and/or $\infty$.
- Some examples of sets that would be either valid or invalid as ROCs for $X$, if $x$ is finite duration, are shown below.


Valid


Invalid


Invalid

- If a sequence $x$ is right sided and the circle $|z|=r_{0}$ is in the ROC of $X=z x$, then all (finite) values of $z$ for which $|z|>r_{0}$ will also be in the ROC of $X$ (i.e., the ROC contains the exterior of a circle centered at 0 , possibly including $\infty$ ).
- Thus, if $x$ is right sided but not left sided, the ROC of $X$ is the exterior of a circle centered at 0 , possibly including $\infty$.
- Examples of sets that would be either valid or invalid as ROCs for $X$, if $x$ is right sided but not left sided, are shown below.


Valid


Invalid


Invalid

- If a sequence $x$ is left sided and the circle $|z|=r_{0}$ is in the ROC of $X=z x$, then all values of $z$ for which $0<|z|<r_{0}$ will also be in the ROC of $X$ (i.e., the ROC contains a disk centered at 0 , possibly excluding 0 ).
- Thus, if $x$ is left sided but not right sided, the ROC of $X$ is a disk centered at 0 , possibly excluding 0 .
- Examples of sets that would be either valid or invalid as ROCs for $X$, if $x$ is left sided but not right sided, are shown below.


Valid


Invalid


Invalid

- If a sequence $x$ is two sided and the circle $|z|=r_{0}$ is in the ROC of $X=z x$, then the ROC of $X$ will consist of a ring that contains this circle (i.e., the ROC is an annulus centered at 0 ).
- Examples of sets that would be either valid or invalid as ROCs for $X$, if $x$ is two sided, are shown below.


Valid


Invalid


Invalid

- If a sequence $x$ has a rational $z$ transform $X$ (with at least one pole), then:

11 If $x$ is right sided, then the ROC of $X$ is the region outside the circle of radius equal to the largest magnitude of the poles of $X$ (i.e., outside the outermost pole), possibly including $\infty$.
2 If $x$ is left sided, then the ROC of $X$ is the region inside the circle of radius equal to the smallest magnitude of the nonzero poles of $X$ and extending inward to, and possibly including, 0 (i.e., inside the innermost nonzero pole).

- This property is implied by properties $1,2,4$, and 5 .
- Some examples of sets that would be either valid or invalid as ROCs for $X$, if $X$ is rational and $x$ is left/right sided, are given below.


Valid


Invalid


Valid


Invalid

## General Form of the ROC

- To summarize the results of properties $3,4,5$, and 6 , if the $z$ transform $X$ of the sequence $x$ exists, the ROC of $X$ depends on the left- and right-sidedness of $x$ as follows:

| $x$ |  |  |
| :--- | :--- | :--- |
| left sided | right sided |  |
| yes | yes | ROC of $X$ |
| no | yes | exterior of circle centered at 0, possibly including $\infty$ |
| yes | no | disk centered at 0, possibly excluding 0 |
| no | no | annulus centered at 0 |

- Thus, we can infer that, if $X$ exists, the ROC can only be of one of the forms listed above.
- For example, the sets shown below would not be valid as ROCs.


Invalid


Invalid

## Section 12.3

## Properties of the z Transform

Properties of the z Transform

| Property | Time Domain | Z Domain | ROC |
| :--- | :--- | :--- | :--- |
| Linearity | $a_{1} x_{1}(n)+a_{2} x_{2}(n)$ | $a_{1} X_{1}(z)+a_{2} X_{2}(z)$ | At least $R_{1} \cap R_{2}$ |
| Translation | $x\left(n-n_{0}\right)$ | $z^{-n_{0}} X(z)$ | $R$ except possible addition/deletion of o |
| Modulation | $a^{n} x(n)$ | $X\left(a^{-1} z\right)$ | $\|a\| R$ |
| Conjugation | $x^{*}(n)$ | $X^{*}\left(z^{*}\right)$ | $R$ |
| Time Reversal | $x(-n)$ | $X(1 / z)$ | $R^{-1}$ |
| Upsampling | $(\uparrow M) x(n)$ | $X\left(z^{M}\right)$ | $R^{1 / M}$ |
| Downsampling | $(\downarrow M) x(n)$ | $\frac{1}{M} \sum_{k=0}^{M-1} X\left(e^{-j 2 \pi k / M} z^{1 / M}\right)$ | $R^{M}$ |
| Convolution | $x_{1} * x_{2}(n)$ | $X_{1}(z) X_{2}(z)$ | At least $R_{1} \cap R_{2}$ |
| Z-Domain Diff. | $n x(n)$ | $-z \frac{d}{d z} X(z)$ | $R$ |
| Differencing | $x(n)-x(n-1)$ | $\left(1-z^{-1}\right) X(z)$ | At least $R \cap\|z\|>0$ |
| Accumulation | $\sum_{k=-\infty}^{n} x(k)$ | $\frac{z}{z-1} X(z)$ | At least $R \cap\|z\|>1$ |


| Property |  |
| :--- | :--- |
| Initial Value Theorem | $x(0)=\lim _{z \rightarrow \infty} X(z)$ |
| Final Value Theorem | $\lim _{n \rightarrow \infty} x(n)=\lim _{z \rightarrow 1}[(z-1) X(z)]$ |


| Pair | $x(n)$ | $X(z)$ | ROC |
| :--- | :--- | :--- | :--- |
| 1 | $\delta(n)$ | 1 | All $z$ |
| 2 | $u(n)$ | $\frac{z}{z-1}=\frac{1}{1-z^{-1}}$ | $\|z\|>1$ |
| 3 | $-u(-n-1)$ | $\frac{z}{z-1}=\frac{1}{1-z^{-1}}$ | $\|z\|<1$ |
| 4 | $n u(n)$ | $\frac{z}{(z-1)^{2}}=\frac{z^{-1}}{\left(1-z^{-1}\right)^{2}}$ | $\|z\|>1$ |
| 5 | $-n u(-n-1)$ | $\frac{z}{(z-1)^{2}}=\frac{z^{-1}}{\left(1-z^{-1}\right)^{2}}$ | $\|z\|<1$ |
| 6 | $a^{n} u(n)$ | $\frac{z}{z-a}=\frac{1}{1-a z^{-1}}$ | $\|z\|>\|a\|$ |
| 7 | $-a^{n} u(-n-1)$ | $\frac{z}{z-a}=\frac{1}{1-a z^{-1}}$ | $\|z\|<\|a\|$ |
| 8 | $n a^{n} u(n)$ | $\frac{a z}{(z-a)^{2}}=\frac{a z^{-1}}{\left(1-a z^{-1}\right)^{2}}$ | $\|z\|>\|a\|$ |
| 9 | $-n a^{n} u(-n-1)$ | $\frac{a z}{(z-a)^{2}}=\frac{a z^{-1}}{\left(1-a z^{-1}\right)^{2}}$ | $\|z\|<\|a\|$ |
| 10 | $\frac{(n+1)(n+2) \cdots(n+m-1)}{(m-1)!} a^{n} u(n)$ | $\frac{z^{n}}{(z-a)^{m}}=\frac{1}{\left(1-a z^{-1}\right)^{m}}$ | $\|z\|>\|a\|$ |
| 11 | $-\frac{(n+1)(n+2) \cdots(n+m-1)}{(m-1)!} a^{n} u(-n-1)$ | $\frac{z^{n}}{(z-a)^{m}}=\frac{1}{\left(1-a z^{-1}\right)^{m}}$ | $\|z\|<\|a\|$ |

## z Transform Pairs (Continued)

| Pair | $x(n)$ | $X(z)$ | ROC |
| :--- | :--- | :--- | :--- |
| 12 | $\cos \left(\Omega_{0} n\right) u(n)$ | $\frac{z\left(z-\cos \Omega_{0}\right)}{z^{2}-2 z \cos \Omega_{0}+1}=\frac{1-\left(\cos \Omega_{0}\right) z^{-1}}{1-\left(2 \cos \Omega_{0} z^{-1}+z^{-2}\right.}$ | $\|z\|>1$ |
| 13 | $-\cos \left(\Omega_{0} n\right) u(-n-1)$ | $\frac{z\left(z-\cos \Omega_{0}\right)}{z^{2}-2 z \cos \Omega_{0}+1}=\frac{1-\left(\cos \Omega_{0}\right) z^{-1}}{1-\left(2 \cos \Omega_{0}\right) z^{-1}+z^{-2}}$ | $\|z\|<1$ |
| 14 | $\sin \left(\Omega_{0} n\right) u(n)$ | $\frac{z \sin \Omega_{0}}{z^{2}-2 z \cos \Omega_{0}+1}=\frac{\left(\sin \Omega_{0}\right) z^{-1}}{1-\left(2 \cos \Omega_{0} z^{-1}+z^{-2}\right.}$ | $\|z\|>1$ |
| 15 | $-\sin \left(\Omega_{0} n\right) u(-n-1)$ | $\frac{z \sin \Omega_{0}}{z^{2}-2 z \cos \Omega_{0}+1}=\frac{\left(\sin \Omega_{0}\right) z^{-1}}{1-\left(2 \cos \Omega_{0}\right) z^{-1}+z^{-2}}$ | $\|z\|<1$ |
| 16 | $a^{n} \cos \left(\Omega_{0} n\right) u(n)$ | $\frac{z\left(z-a \cos \Omega_{0}\right)}{z^{2}-2 a z \cos \Omega_{0}+a^{2}}=\frac{1-\left(a \cos \Omega_{0} z^{-1}\right.}{1-\left(2 a \cos \Omega_{0}\right) z^{-1}+a^{2} z^{-2}}$ | $\|z\|>\|a\|$ |
| 17 | $a^{n} \sin \left(\Omega_{0} n\right) u(n)$ | $\frac{a z \sin \Omega_{0}}{z^{2}-2 a \cos \Omega_{0}+a^{2}}=\frac{\left(a \sin \Omega_{0}\right) z^{-1}}{1-\left(2 a \cos \Omega_{0}\right) z^{-1}+a^{2} z^{-2}}$ | $\|z\|>\|a\|$ |
| 18 | $u(n)-u(n-M), M>0$ | $\frac{z\left(1-z^{-M}\right)}{z-1}=\frac{1-z^{-M}}{1-z^{-1}}$ | $\|z\|>0$ |
| 19 | $a^{\|n\|},\|a\|<1$ | $\frac{\left(a-a^{-1}\right) z}{(z-a)\left(z-a^{-1}\right)}$ | $\|a\|<\|z\|<\left\|a^{-1}\right\|$ |

■ If $x_{1}(n) \stackrel{\text { ZT }}{\longleftrightarrow} X_{1}(z)$ with ROC $R_{1}$ and $x_{2}(n) \stackrel{\text { zT }}{\longleftrightarrow} X_{2}(z)$ with ROC $R_{2}$, then $a_{1} x_{1}(n)+a_{2} x_{2}(n) \stackrel{\mathrm{zT}}{\longleftrightarrow} a_{1} X_{1}(z)+a_{2} X_{2}(z) \quad$ with ROC $R$ containing $R_{1} \cap R_{2}$, where $a_{1}$ and $a_{2}$ are arbitrary complex constants.

- This is known as the linearity property of the $z$ transform.
- The ROC always contains the intersection but could be larger (in the case that pole-zero cancellation occurs).
- If $x(n) \stackrel{\text { T }}{\longleftrightarrow} X(z)$ with ROC $R$, then

$$
x\left(n-n_{0}\right) \stackrel{Z T}{\longleftrightarrow} z^{-n_{0}} X(z) \text { with ROC } R^{\prime},
$$

where $n_{0}$ is an integer constant and $R^{\prime}$ is the same as $R$ except for the possible addition or deletion of zero or infinity.

- This is known as the translation (or time-shifting) property of the $z$ transform.


## Z-Domain Scaling

■ If $x(n) \stackrel{\text { zT }}{\longleftrightarrow} X(z)$ with ROC $R$, then

$$
a^{n} x(n) \stackrel{\mathrm{zT}}{\longleftrightarrow} X(z / a) \quad \text { with ROC }|a| R,
$$

where $a$ is a nonzero constant.

- This is known as the z -domain scaling property of the z transform.
- As illustrated below, the ROC $R$ is scaled by $|a|$.


R

$|a| R$

- If $x(n) \stackrel{\mathrm{zT}}{\longleftrightarrow} X(z)$ with ROC $R$, then

$$
x(-n) \stackrel{\mathrm{zT}}{\longleftrightarrow} X(1 / z) \quad \text { with ROC } 1 / R .
$$

■ This is known as the time-reversal property of the $z$ transform.

- As illustrated below, the ROC $R$ is reciprocated.


- Define $(\uparrow M) x(n)$ as

$$
(\uparrow M) x(n)= \begin{cases}x(n / M) & n / M \text { is an integer } \\ 0 & \text { otherwise }\end{cases}
$$

- If $x(n) \stackrel{z T}{\longleftrightarrow} X(z)$ with ROC $R$, then

$$
(\uparrow M) x(n) \stackrel{\mathrm{ZT}}{\longleftrightarrow} X\left(z^{M}\right) \quad \text { with ROC } R^{1 / M}
$$

- This is known as the upsampling (or time-expansion) property of the $z$ transform.


## Downsampling

- If $x(n) \stackrel{\text { zT }}{\longleftrightarrow} X(z)$ with ROC $R$, then

$$
(\downarrow M) x(n) \stackrel{\mathrm{ZT}}{\longleftrightarrow} \frac{1}{M} \sum_{k=0}^{M-1} X\left(e^{-j 2 \pi k / M} z^{1 / M}\right) \quad \text { with ROC } R^{M} .
$$

- This is known as the downsampling property of the $z$ transform.


## Conjugation

■ If $x(n) \stackrel{\text { zT }}{\longleftrightarrow} X(z)$ with ROC $R$, then

$$
x^{*}(n) \stackrel{\mathrm{ZT}}{\longleftrightarrow} X^{*}\left(z^{*}\right) \quad \text { with ROC } R .
$$

- This is known as the conjugation property of the $z$ transform.


## Convolution

- If $x_{1}(n) \stackrel{\text { zT }}{\longleftrightarrow} X_{1}(z)$ with ROC $R_{1}$ and $x_{2}(n) \stackrel{\text { zT }}{\longleftrightarrow} X_{2}(z)$ with ROC $R_{2}$, then

$$
x_{1} * x_{2}(n) \stackrel{z T}{\longleftrightarrow} X_{1}(z) X_{2}(z) \quad \text { with ROC containing } R_{1} \cap R_{2} .
$$

- This is known that the convolution (or time-domain convolution) property of the $z$ transform.
- The ROC always contains the intersection but can be larger than the intersection (if pole-zero cancellation occurs).
- Convolution in the time domain becomes multiplication in the $z$ domain.
- This can make dealing with LTI systems much easier in the $z$ domain than in the time domain.


## Z-Domain Differentiation

■ If $x(n) \stackrel{2 T}{\longleftrightarrow} X(z)$ with ROC $R$, then

$$
n x(n) \stackrel{Z T}{\longleftrightarrow}-z \frac{d}{d z} X(z) \quad \text { with ROC } R .
$$

- This is known as the z-domain differentiation property of the $z$ transform.


## Differencing

- If $x(n) \stackrel{\text { ZT }}{\longleftrightarrow} X(z)$ with ROC $R$, then

$$
x(n)-x(n-1) \stackrel{\text { ZT }}{\longleftrightarrow}\left(1-z^{-1}\right) X(z) \text { for ROC containing } R \cap|z|>0 .
$$

- This is known as the differencing property of the $z$ transform.
- Differencing in the time domain becomes multiplication by $1-z^{-1}$ in the $z$ domain.
- This can make dealing with difference equations much easier in the $z$ domain than in the time domain.


## Accumulation

- If $x(n) \stackrel{\text { ZT }}{\longleftrightarrow} X(z)$ with ROC $R$, then

$$
\sum_{k=-\infty}^{n} x(k) \stackrel{z T}{\longleftrightarrow} \frac{z}{z-1} X(z) \text { for ROC containing } R \cap|z|>1
$$

- This is known as the accumulation property of the $z$ transform.


## Initial Value Theorem

- For a sequence $x$ with z transform $X$, if $x$ is causal, then

$$
x(0)=\lim _{z \rightarrow \infty} X(z) .
$$

- This result is known as the initial-value theorem.
- For a sequence $x$ with $z$ transform $X$, if $x$ is causal and $\lim _{n \rightarrow \infty} x(n)$ exists, then

$$
\lim _{n \rightarrow \infty} x(n)=\lim _{z \rightarrow 1}[(z-1) X(z)] .
$$

- This result is known as the final-value theorem.


## More z Transform Examples

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Section 12.4

## Determination of Inverse z Transform

- Recall that the inverse $z$ transform $x$ of $X$ is given by

$$
x(n)=\frac{1}{2 \pi j} \oint_{\Gamma} X(z) z^{n-1} d z
$$

where $\Gamma$ is a counterclockwise closed circular contour centered at the origin and with radius $r$ such that $\Gamma$ is in the ROC of $X$.

- Unfortunately, the above contour integration can often be quite tedious to compute.
- Consequently, we do not usually compute the inverse z transform directly using the above equation.
- For rational functions, the inverse z transform can be more easily computed using partial fraction expansions.
- Using a partial fraction expansion, we can express a rational function as a sum of lower-order rational functions whose inverse $z$ transforms can typically be found in tables.


## Section 12.5

## z Transform and LTI Systems

- Consider a LTI system with input $x$, output $y$, and impulse response $h$, and let $X, Y$, and $H$ denote the $z$ transforms of $x, y$, and $h$, respectively.
- Since $y(n)=x * h(n)$, the system is characterized in the z domain by

$$
Y(z)=X(z) H(z)
$$

- As a matter of terminology, we refer to $H$ as the system function (or transfer function) of the system (i.e., the system function is the $z$ transform of the impulse response).
- When viewed in the $z$ domain, a LTI system forms its output by multiplying its input with its system function.
- A LTI system is completely characterized by its system function $H$.
- If the ROC of $H$ includes the unit circle $|z|=1$, then $H\left(e^{j \Omega}\right)$ is the frequency response of the LTI system.

■ Consider a LTI system with input $x$, output $y$, and impulse response $h$, and let $X, Y$, and $H$ denote the $z$ transforms of $x, y$, and $h$, respectively.

- Often, it is convenient to represent such a system in block diagram form in the $z$ domain as shown below.

- Since a LTI system is completely characterized by its system function, we typically label the system with this quantity.


## Interconnection of LTI Systems

- The series interconnection of the LTI systems with system functions $H_{1}$ and $H_{2}$ is the LTI system with system function $H=H_{1} H_{2}$. That is, we have the equivalences shown below.

- The parallel interconnection of the LTI systems with impulse responses $H_{1}$ and $H_{2}$ is a LTI system with the system function $H=H_{1}+H_{2}$. That is, we have the equivalence shown below.



## Causality

- If a LTI system is causal, its impulse response is causal, and therefore right sided. From this, we have the result below.
- Theorem. A LTI system is causal if and only if the ROC of the system function is:

1 the exterior of a circle, including $\infty$; or
2 the entire complex plane, including $\infty$ and possibly excluding 0 .

- Theorem. A LTI system with a rational system function $H$ is causal if and only if:

1 the ROC of $H$ is the exterior of a (possibly degenerate) circle outside the outermost pole of $H$ or, if $H$ has no poles, the entire complex plane; and
$2 H$ is proper (i.e., when $H(z)$ is expressed as a ratio of polynomials in $z$, the order of the numerator polynomial does not exceed the order of the denominator polynomial).

- Whether or not a system is BIBO stable depends on the ROC of its system function.
- Theorem. A LTI system is BIBO stable if and only if the ROC of its system function contains the unit circle (i.e., $|z|=1$ ).
- Theorem. A causal LTI system with a rational system function $H$ is BIBO stable if and only if all of the poles of $H$ lie inside the unit circle (i.e., each of the poles has a magnitude less than one).


## Invertibility

- A LTI system $\mathcal{H}$ with system function $H$ is invertible if and only if there exists another LTI system with system function $H_{\text {inv }}$ such that

$$
H(z) H_{\mathrm{inv}}(z)=1
$$

in which case $H_{\text {inv }}$ is the system function of $\mathcal{H}^{-1}$ and

$$
H_{\mathrm{inv}}(z)=\frac{1}{H(z)}
$$

- Since distinct systems can have identical system functions (but with differing ROCs), the inverse of a LTI system is not necessarily unique.
- In practice, however, we often desire a stable and/or causal system. So, although multiple inverse systems may exist, we are frequently only interested in one specific choice of inverse system (due to these additional constraints of stability and/or causality).


## LTI Systems and Difference Equations

■ Many LTI systems of practical interest can be represented using an Nth-order linear difference equation with constant coefficients.

- Consider a system with input $x$ and output $y$ that is characterized by an equation of the form

$$
\sum_{k=0}^{N} b_{k} y(n-k)=\sum_{k=0}^{M} a_{k} x(n-k) \quad \text { where } \quad M \leq N
$$

■ Let $h$ denote the impulse response of the system, and let $X, Y$, and $H$ denote the $z$ transforms of $x, y$, and $h$, respectively.

- One can show that $H(z)$ is given by

$$
H(z)=\frac{Y(z)}{X(z)}=\frac{\sum_{k=0}^{M} a_{k} z^{-k}}{\sum_{k=0}^{N} b_{k} z^{-k}}
$$

- Observe that, for a system of the form considered above, the system function is always rational.


## Section 12.6

## Application: Analysis of Control Systems

Reference


Signal

- input: desired value of the quantity to be controlled
- output: actual value of the quantity to be controlled
- error: difference between the desired and actual values
- plant: system to be controlled
- sensor: device used to measure the actual output
- controller: device that monitors the error and changes the input of the plant with the goal of forcing the error to zero
- Often, we want to ensure that a system is BIBO stable.
- The BIBO stability property is more easily characterized in the $z$ domain than in the time domain.
- Therefore, the $z$ domain is extremely useful for the stability analysis of systems.


## Section 12.7

## Unilateral z Transform

- The unilateral z transform of the sequence $x$, denoted $\mathcal{Z}_{u} x$ or $X$, is defined as

$$
Z_{\mathrm{u}} x(z)=X(z)=\sum_{n=0}^{\infty} x(n) z^{-n}
$$

- The unilateral $z$ transform is related to the bilateral $z$ transform as follows:

$$
z_{u} x(z)=\sum_{n=0}^{\infty} x(n) z^{-n}=\sum_{n=-\infty}^{\infty} x(n) u(n) z^{-n}=z\{x u\}(z)
$$

- In other words, the unilateral $z$ transform of the sequence $x$ is simply the bilateral $z$ transform of the sequence $x u$.
- Since $z_{u} x=Z\{x u\}$ and $x u$ is always a right-sided sequence, the ROC associated with $z_{u} x$ is always the exterior of a circle.
- For this reason, we often do not explicitly indicate the ROC when working with the unilateral $z$ transform.
- With the unilateral $z$ transform, the same inverse transform equation is used as in the bilateral case.
- The unilateral $z$ transform is only invertible for causal sequences. In particular, we have

$$
\begin{aligned}
Z_{\mathrm{u}}^{-1}\left\{Z_{\mathrm{u}}\{x\}\right\}(n) & =Z_{\mathrm{u}}^{-1}\{Z\{x u\}\}(n) \\
& =Z^{-1}\{Z\{x u\}\}(n) \\
& =x(n) u(n) \\
& = \begin{cases}x(n) & n \geq 0 \\
0 & \text { otherwise }\end{cases}
\end{aligned}
$$

■ For a noncausal sequence $x$, we can only recover $x(n)$ for $n \geq 0$.

- Due to the close relationship between the unilateral and bilateral $z$ transforms, these two transforms have some similarities in their properties.
- Since these two transforms are not identical, however, their properties differ in some cases, often in subtle ways.


## Properties of the Unilateral z Transform

| Property | Time Domain | Z Domain |
| :--- | :--- | :--- |
| Linearity | $a_{1} x_{1}(n)+a_{2} x_{2}(n)$ | $a_{1} X_{1}(z)+a_{2} X_{2}(z)$ |
| Time Delay | $x(n-1)$ | $z^{-1} X(z)+x(-1)$ |
| Time Advance | $x(n+1)$ | $z X(z)-z x(0)$ |
| Modulation | $a^{n} x(n)$ | $X\left(a^{-1} z\right)$ |
|  | $e^{j \Omega_{0} n} x(n)$ | $X\left(e^{\left.-j \Omega_{0} z\right)}\right.$ |
| Conjugation | $x^{*}(n)$ | $X^{*}\left(z^{*}\right)$ |
| Upsampling | $(\uparrow M) x(n)$ | $X\left(z^{M}\right)$ |
| Downsampling | $(\downarrow M) x(n)$ | $\frac{1}{M} \sum_{k=0}^{M-1} X\left(e^{\left.-j 2 \pi k / M z^{1 / M}\right)}\right.$ |
| Convolution | $x_{1} * x_{2}(n), x_{1}$ and $x_{2}$ are causal | $X_{1}(z) X_{2}(z)$ |
| Z-Domain Diff. | $n x(n)$ | $-z \frac{d}{d z} X(z)$ |
| Differencing | $x(n)-x(n-1)$ | $\left(1-z^{-1}\right) X(z)-x(-1)$ |
| Accumulation | $\sum_{k=0}^{n} x(k)$ | $\frac{1}{1-z^{-1} X(z)}$ |


| Property |  |
| :--- | :--- |
| Initial Value Theorem | $x(0)=\lim _{z \rightarrow \infty} X(z)$ |
| Final Value Theorem | $\lim _{n \rightarrow \infty} x(n)=\lim _{z \rightarrow 1}[(z-1) X(z)]$ |

## Unilateral z Transform Pairs

| Pair | $x(n), n \geq 0$ | $X(z)$ |
| :--- | :--- | :--- |
| 1 | $\delta(n)$ | 1 |
| 2 | 1 | $\frac{z}{z-1}$ |
| 3 | $n$ | $\frac{z}{(z-1)^{2}}$ |
| 4 | $a^{n}$ | $\frac{z}{z-a}$ |
| 5 | $a^{n} n$ | $\frac{a z}{(z-a)^{2}}$ |
| 6 | $\cos \left(\Omega_{0} n\right)$ | $\frac{z\left(z-\cos \Omega_{0}\right)}{z^{2}-2\left(\cos \Omega_{0}\right) z+1}$ |
| 7 | $\sin \left(\Omega_{0} n\right)$ | $\frac{z \sin \Omega_{0}}{z^{2}-2\left(\cos \Omega_{0}\right) z+1}$ |
| 8 | $\|a\|^{n} \cos \left(\Omega_{0} n\right)$ | $\frac{z\left(z-\|a\| \cos \Omega_{0}\right)}{z^{2}-2\|a\|\left(\cos \Omega_{0}\right) z+\|a\|^{2}}$ |
| 9 | $\|a\|^{n} \sin \left(\Omega_{0} n\right)$ | $\frac{z\|a\| \sin \Omega_{0}}{z^{2}-2\|a\|\left(\cos \Omega_{0}\right) z+\|a\|^{2}}$ |

■ Many systems of interest in engineering applications can be characterized by constant-coefficient linear difference equations.
■ One common use of the unilateral $z$ transform is in solving constant-coefficient linear difference equations with nonzero initial conditions.

## Part 13

## Complex Analysis

## Complex Numbers

- A complex number is a number of the form $z=x+j y$ where $x$ and $y$ are real numbers and $j$ is the constant defined by $j^{2}=-1$ (i.e., $j=\sqrt{-1}$ ).
- The Cartesian form of the complex number $z$ expresses $z$ in the form

$$
z=x+j y
$$

where $x$ and $y$ are real numbers. The quantities $x$ and $y$ are called the real part and imaginary part of $z$, and are denoted as $\operatorname{Re} z$ and $\operatorname{Im} z$, respectively.

- The polar form of the complex number $z$ expresses $z$ in the form

$$
z=r(\cos \theta+j \sin \theta) \quad \text { or equivalently } \quad z=r e^{j \theta}
$$

where $r$ and $\theta$ are real numbers and $r \geq 0$. The quantities $r$ and $\theta$ are called the magnitude and argument of $z$, and are denoted as $|z|$ and $\arg z$, respectively. [Note: $e^{j \theta}=\cos \theta+j \sin \theta$.]

## Complex Numbers (Continued)

■ Since $e^{j \theta}=e^{j(\theta+2 \pi k)}$ for all real $\theta$ and all integer $k$, the argument of a complex number is only uniquely determined to within an additive multiple of $2 \pi$.

- The principal argument of a complex number $z$, denoted $\operatorname{Arg} z$, is the particular value $\theta$ of $\arg z$ that satisfies $-\pi<\theta \leq \pi$.
- The principal argument of a complex number (excluding zero) is unique.


## Geometric Interpretation of Cartesian and Polar Forms




Cartesian form:
$z=x+j y$
where $x=\operatorname{Re} z$ and $y=\operatorname{Im} z$

$$
\begin{aligned}
& \text { Polar form: } \\
& z=r(\cos \theta+j \sin \theta)=r e^{j \theta} \\
& \text { where } r=|z| \text { and } \theta=\arg z
\end{aligned}
$$

- The range of the arctan function is $-\pi / 2$ (exclusive) to $\pi / 2$ (exclusive).
- Consequently, the arctan function always yields an angle in either the first or fourth quadrant.


- The angle $\theta$ that a vector from the origin to the point $(x, y)$ makes with the positive $x$ axis is given by $\theta=\operatorname{atan} 2(y, x)$, where

$$
\operatorname{atan} 2(y, x) \triangleq \begin{cases}\arctan (y / x) & x>0 \\ \pi / 2 & x=0 \text { and } y>0 \\ -\pi / 2 & x=0 \text { and } y<0 \\ \arctan (y / x)+\pi & x<0 \text { and } y \geq 0 \\ \arctan (y / x)-\pi & x<0 \text { and } y<0\end{cases}
$$

- The range of the atan2 function is from $-\pi$ (exclusive) to $\pi$ (inclusive).
- For the complex number $z$ expressed in Cartesian form $x+j y$, $\operatorname{Arg} z=\operatorname{atan} 2(y, x)$.
- Although the atan2 function is quite useful for computing the principal argument (or argument) of a complex number, it is not advisable to memorize the definition of this function. It is better to simply understand what this function is doing (namely, intelligently applying the arctan function).


## Conversion Between Cartesian and Polar Form

- Let $z$ be a complex number with the Cartesian and polar form representations given respectively by

$$
z=x+j y \quad \text { and } \quad z=r e^{j \theta}
$$

- To convert from polar to Cartesian form, we use the following identities:

$$
x=r \cos \theta \quad \text { and } \quad y=r \sin \theta .
$$

- To convert from Cartesian to polar form, we use the following identities:

$$
r=\sqrt{x^{2}+y^{2}} \quad \text { and } \quad \theta=\operatorname{atan} 2(y, x)+2 \pi k
$$

where $k$ is an arbitrary integer.

- Since the atan 2 function simply amounts to the intelligent application of the arctan function, instead of memorizing the definition of the atan 2 function, one should simply understand how to use the arctan function to achieve the same result.
- For complex numbers, addition and multiplication are commutative. That is, for any two complex numbers $z_{1}$ and $z_{2}$,

$$
\begin{gathered}
z_{1}+z_{2}=z_{2}+z_{1} \quad \text { and } \\
z_{1} z_{2}=z_{2} z_{1}
\end{gathered}
$$

- For complex numbers, addition and multiplication are associative. That is, for any three complex numbers $z_{1}, z_{2}$, and $z_{3}$,

$$
\begin{gathered}
\left(z_{1}+z_{2}\right)+z_{3}=z_{1}+\left(z_{2}+z_{3}\right) \quad \text { and } \\
\left(z_{1} z_{2}\right) z_{3}=z_{1}\left(z_{2} z_{3}\right)
\end{gathered}
$$

- For complex numbers, the distributive property holds. That is, for any three complex numbers $z_{1}, z_{2}$, and $z_{3}$,

$$
z_{1}\left(z_{2}+z_{3}\right)=z_{1} z_{2}+z_{1} z_{3}
$$

## Conjugation

- The conjugate of the complex number $z=x+j y$ is denoted as $z^{*}$ and defined as

$$
z^{*}=x-j y .
$$

- Geometrically, the conjugation operation reflects a point in the complex plane about the real axis.
- The geometric interpretation of the conjugate is illustrated below.

- For every complex number $z$, the following identities hold:

$$
\begin{gathered}
\left|z^{*}\right|=|z|, \\
\arg z^{*}=-\arg z, \\
z z^{*}=|z|^{2}, \\
\operatorname{Re} z=\frac{1}{2}\left(z+z^{*}\right), \quad \text { and } \\
\operatorname{Im} z=\frac{1}{2 j}\left(z-z^{*}\right) .
\end{gathered}
$$

- For all complex numbers $z_{1}$ and $z_{2}$, the following identities hold:

$$
\begin{gathered}
\left(z_{1}+z_{2}\right)^{*}=z_{1}^{*}+z_{2}^{*} \\
\left(z_{1} z_{2}\right)^{*}=z_{1}^{*} z_{2}^{*}, \quad \text { and } \\
\left(z_{1} / z_{2}\right)^{*}=z_{1}^{*} / z_{2}^{*} .
\end{gathered}
$$

## Addition

- Cartesian form: Let $z_{1}=x_{1}+j y_{1}$ and $z_{2}=x_{2}+j y_{2}$. Then,

$$
\begin{aligned}
z_{1}+z_{2} & =\left(x_{1}+j y_{1}\right)+\left(x_{2}+j y_{2}\right) \\
& =\left(x_{1}+x_{2}\right)+j\left(y_{1}+y_{2}\right) .
\end{aligned}
$$

- That is, to add complex numbers expressed in Cartesian form, we simply add their real parts and add their imaginary parts.
- Polar form: Let $z_{1}=r_{1} e^{j \theta_{1}}$ and $z_{2}=r_{2} e^{j \theta_{2}}$. Then,

$$
\begin{aligned}
z_{1}+z_{2} & =r_{1} e^{j \theta_{1}}+r_{2} e^{j \theta_{2}} \\
& =\left(r_{1} \cos \theta_{1}+j r_{1} \sin \theta_{1}\right)+\left(r_{2} \cos \theta_{2}+j r_{2} \sin \theta_{2}\right) \\
& =\left(r_{1} \cos \theta_{1}+r_{2} \cos \theta_{2}\right)+j\left(r_{1} \sin \theta_{1}+r_{2} \sin \theta_{2}\right) .
\end{aligned}
$$

■ That is, to add complex numbers expressed in polar form, we first rewrite them in Cartesian form, and then add their real parts and add their imaginary parts.

- For the purposes of addition, it is easier to work with complex numbers expressed in Cartesian form.
- Cartesian form: Let $z_{1}=x_{1}+j y_{1}$ and $z_{2}=x_{2}+j y_{2}$. Then,

$$
\begin{aligned}
z_{1} z_{2} & =\left(x_{1}+j y_{1}\right)\left(x_{2}+j y_{2}\right) \\
& =x_{1} x_{2}+j x_{1} y_{2}+j x_{2} y_{1}-y_{1} y_{2} \\
& =\left(x_{1} x_{2}-y_{1} y_{2}\right)+j\left(x_{1} y_{2}+x_{2} y_{1}\right)
\end{aligned}
$$

- That is, to multiply two complex numbers expressed in Cartesian form, we use the distributive law along with the fact that $j^{2}=-1$.
$\square$ Polar form: Let $z_{1}=r_{1} e^{j \theta_{1}}$ and $z_{2}=r_{2} e^{j \theta_{2}}$. Then,

$$
z_{1} z_{2}=\left(r_{1} e^{j \theta_{1}}\right)\left(r_{2} e^{j \theta_{2}}\right)=r_{1} r_{2} e^{j\left(\theta_{1}+\theta_{2}\right)}
$$

- That is, to multiply two complex numbers expressed in polar form, we use exponent rules.
- For the purposes of multiplication, it is easier to work with complex numbers expressed in polar form.

■ Cartesian form: Let $z_{1}=x_{1}+j y_{1}$ and $z_{2}=x_{2}+j y_{2}$. Then,

$$
\begin{aligned}
\frac{z_{1}}{z_{2}} & =\frac{z_{1} z_{2}^{*}}{z_{2} z_{2}^{*}}=\frac{z_{1} z_{2}^{*}}{\left|z_{2}\right|^{2}}=\frac{\left(x_{1}+j y_{1}\right)\left(x_{2}-j y_{2}\right)}{x_{2}^{2}+y_{2}^{2}} \\
& =\frac{x_{1} x_{2}-j x_{1} y_{2}+j x_{2} y_{1}+y_{1} y_{2}}{x_{2}^{2}+y_{2}^{2}}=\frac{x_{1} x_{2}+y_{1} y_{2}+j\left(x_{2} y_{1}-x_{1} y_{2}\right)}{x_{2}^{2}+y_{2}^{2}}
\end{aligned}
$$

- That is, to compute the quotient of two complex numbers expressed in Cartesian form, we convert the problem into one of division by a real number.
■ Polar form: Let $z_{1}=r_{1} e^{j \theta_{1}}$ and $z_{2}=r_{2} e^{j \theta_{2}}$. Then,

$$
\frac{z_{1}}{z_{2}}=\frac{r_{1} e^{j \theta_{1}}}{r_{2} e^{j \theta_{2}}}=\frac{r_{1}}{r_{2}} e^{j\left(\theta_{1}-\theta_{2}\right)}
$$

- That is, to compute the quotient of two complex numbers expressed in polar form, we use exponent rules.
- For the purposes of division, it is easier to work with complex numbers expressed in polar form.
- For arbitrary complex numbers $z_{1}$ and $z_{2}$, the following identities hold:

$$
\begin{gathered}
\left|z_{1} z_{2}\right|=\left|z_{1}\right|\left|z_{2}\right| \\
\left|\frac{z_{1}}{z_{2}}\right|=\frac{\left|z_{1}\right|}{\left|z_{2}\right|} \quad \text { for } z_{2} \neq 0 ; \\
\arg \left(z_{1} z_{2}\right)=\arg z_{1}+\arg z_{2} ; \quad \text { and } \\
\arg \left(\frac{z_{1}}{z_{2}}\right)=\arg z_{1}-\arg z_{2} \quad \text { for } z_{2} \neq 0 .
\end{gathered}
$$

- The above four identities trivially follow from the polar representation of complex numbers.
- The first and third identities above imply two additional identities, as given below.
- For all complex numbers $z$ and all integers $n$, the following identities hold:

$$
\begin{aligned}
& \left|z^{n}\right|=|z|^{n} ; \quad \text { and } \\
& \arg \left(z^{n}\right)=n \arg z .
\end{aligned}
$$

- Euler's relation. For all real $\theta$,

$$
e^{j \theta}=\cos \theta+j \sin \theta
$$

- From Euler's relation, we can deduce the following useful identities:

$$
\begin{gathered}
\cos \theta=\frac{1}{2}\left(e^{j \theta}+e^{-j \theta}\right) \quad \text { and } \\
\sin \theta=\frac{1}{2 j}\left(e^{j \theta}-e^{-j \theta}\right) .
\end{gathered}
$$

- De Moivre's theorem. For all real $\theta$ and all integer $n$,

$$
e^{j n \theta}=\left(e^{j \theta}\right)^{n}
$$

[Note: This relationship does not necessarily hold for real n.]

- Every complex number $z=r e^{j \theta}$ (where $r=|z|$ and $\theta=\arg z$ ) has $n$ distinct $n$th roots given by

$$
\sqrt[n]{r} e^{j(\theta+2 \pi k) / n} \quad \text { for } k=0,1, \ldots, n-1
$$

- For example, 1 has the two distinct square roots 1 and -1 .


## Quadratic Formula

- Consider the equation

$$
a z^{2}+b z+c=0,
$$

where $a, b$, and $c$ are real, $z$ is complex, and $a \neq 0$.

- The roots of this equation are given by

$$
z=\frac{-b \pm \sqrt{b^{2}-4 a c}}{2 a}
$$

■ This formula is often useful in factoring quadratic polynomials.

- The quadratic $a z^{2}+b z+c$ can be factored as $a\left(z-z_{0}\right)\left(z-z_{1}\right)$, where

$$
z_{0}=\frac{-b-\sqrt{b^{2}-4 a c}}{2 a} \quad \text { and } \quad z_{1}=\frac{-b+\sqrt{b^{2}-4 a c}}{2 a} .
$$

## Complex Functions

- A complex function maps complex numbers to complex numbers. For example, the function $F(z)=z^{2}+2 z+1$, where $z$ is complex, is a complex function.
- A complex polynomial function is a mapping of the form

$$
F(z)=a_{0}+a_{1} z+a_{2} z^{2}+\cdots+a_{n} z^{n}
$$

where $z, a_{0}, a_{1}, \ldots, a_{n}$ are complex.

- A complex rational function is a mapping of the form

$$
F(z)=\frac{a_{0}+a_{1} z+a_{2} z^{2}+\ldots+a_{n} z^{n}}{b_{0}+b_{1} z+b_{2} z^{2}+\ldots+b_{m} z^{m}}
$$

where $a_{0}, a_{1}, \ldots, a_{n}, b_{0}, b_{1}, \ldots, b_{m}$ and $z$ are complex.

- Observe that a polynomial function is a special case of a rational function.
- Herein, we will mostly focus our attention on polynomial and rational functions.


## Continuity

- A function $F$ is said to be continuous at a point $z_{0}$ if $F\left(z_{0}\right)$ is defined and given by

$$
F\left(z_{0}\right)=\lim _{z \rightarrow z_{0}} F(z)
$$

- A function that is continuous at every point in its domain is said to be continuous.
- Polynomial functions are continuous everywhere.
- Rational functions are continuous everywhere except at points where the denominator polynomial becomes zero.


## Differentiability

- A function $F$ is said to be differentiable at a point $z=z_{0}$ if the limit

$$
F^{\prime}\left(z_{0}\right)=\lim _{z \rightarrow z_{0}} \frac{F(z)-F\left(z_{0}\right)}{z-z_{0}}
$$

exists. This limit is called the derivative of $F$ at the point $z=z_{0}$.

- A function is said to be differentiable if it is differentiable at every point in its domain.
- The rules for differentiating sums, products, and quotients are the same for complex functions as for real functions. If $F^{\prime}\left(z_{0}\right)$ and $G^{\prime}\left(z_{0}\right)$ exist, then
$1(a F)^{\prime}\left(z_{0}\right)=a F^{\prime}\left(z_{0}\right)$ for any complex constant $a$;
2 $(F+G)^{\prime}\left(z_{0}\right)=F^{\prime}\left(z_{0}\right)+G^{\prime}\left(z_{0}\right)$;
3 $(F G)^{\prime}\left(z_{0}\right)=F^{\prime}\left(z_{0}\right) G\left(z_{0}\right)+F\left(z_{0}\right) G^{\prime}\left(z_{0}\right)$;
$4(F / G)^{\prime}\left(z_{0}\right)=\frac{G\left(z_{0}\right) F^{\prime}\left(z_{0}\right)-F\left(z_{0}\right) G^{\prime}\left(z_{0}\right)}{G\left(z_{0}\right)^{2}}$; and
5 if $z_{0}=G\left(w_{0}\right)$ and $G^{\prime}\left(w_{0}\right)$ exists, then the derivative of $F(G(z))$ at $w_{0}$ is $F^{\prime}\left(z_{0}\right) G^{\prime}\left(w_{0}\right)$ (i.e., the chain rule).
- A polynomial function is differentiable everywhere.
- A rational function is differentiable everywhere except at the points where its denominator polynomial becomes zero.


## Open Disks

- An open disk in the complex plane with center $z_{0}$ and radius $r$ is the set of complex numbers $z$ satisfying

$$
\left|z-z_{0}\right|<r,
$$

where $r$ is a strictly positive real number.

- A plot of an open disk is shown below.



## Analyticity

- A function is said to be analytic at a point $z_{0}$ if it is differentiable at every point in an open disk about $z_{0}$.
- A function is said to be analytic if it is analytic at every point in its domain.
- A polynomial function is analytic everywhere.
- A rational function is analytic everywhere, except at the points where its denominator polynomial becomes zero.


## Zeros and Singularities

- If a function $F$ is zero at the point $z_{0}$ (i.e., $F\left(z_{0}\right)=0$ ), $F$ is said to have a zero at $z_{0}$.
- If a function $F$ is such that $F\left(z_{0}\right)=0, F^{(1)}\left(z_{0}\right)=0, \ldots, F^{(n-1)}\left(z_{0}\right)=0$ (where $F^{(k)}$ denotes the $k$ th order derivative of $F$ ), $F$ is said to have an $n$th order zero at $z_{0}$.
- A point at which a function fails to be analytic is called a singularity.
- Polynomials do not have singularities.
- Rational functions can have a type of singularity called a pole.
- If a function $F$ is such that $G(z)=1 / F(z)$ has an $n$th order zero at $z_{0}, F$ is said to have an $n$th order pole at $z_{0}$.
- A pole of first order is said to be simple, whereas a pole of order two or greater is said to be repeated. A similar terminology can also be applied to zeros (i.e., simple zero and repeated zero).
- Given a rational function $F$, we can always express $F$ in factored form as

$$
F(z)=\frac{K\left(z-a_{1}\right)^{\alpha_{1}}\left(z-a_{2}\right)^{\alpha_{2}} \cdots\left(z-a_{M}\right)^{\alpha_{M}}}{\left(z-b_{1}\right)^{\beta_{1}}\left(z-b_{2}\right)^{\beta_{2}} \cdots\left(z-b_{N}\right)^{\beta_{N}}}
$$

where $K$ is complex, $a_{1}, a_{2}, \ldots, a_{M}, b_{1}, b_{2}, \ldots, b_{N}$ are distinct complex numbers, and $\alpha_{1}, \alpha_{2}, \ldots, \alpha_{M}$ and $\beta_{1}, \beta_{2}, \ldots, \beta_{N}$ are strictly positive integers.
■ One can show that $F$ has poles at $b_{1}, b_{2}, \ldots, b_{N}$ and zeros at $a_{1}, a_{2}, \ldots, a_{M}$.
■ Furthermore, the $k$ th pole (i.e., $b_{k}$ ) is of $\operatorname{order} \beta_{k}$, and the $k$ th zero (i.e., $a_{k}$ ) is of order $\alpha_{k}$.

- When plotting zeros and poles in the complex plane, the symbols "o" and " $x$ " are used to denote zeros and poles, respectively.


## Part 14

## Partial Fraction Expansions (PFEs)

## Motivation for PFEs

- Sometimes it is beneficial to be able to express a rational function as a sum of lower-order rational functions.
- This can be accomplished using a type of decomposition known as a partial fraction expansion.
- Partial fraction expansions are often useful in the calculation of inverse Laplace transforms, inverse z transforms, and inverse CT/DT Fourier transforms.
- Consider a rational function

$$
F(v)=\frac{\alpha_{m} v^{m}+\alpha_{m-1} v^{m-1}+\ldots+\alpha_{1} v+\alpha_{0}}{\beta_{n} v^{n}+\beta_{n-1} v^{n-1}+\ldots+\beta_{1} v+\beta_{0}}
$$

- The function $F$ is said to be strictly proper if $m<n$ (i.e., the order of the numerator polynomial is strictly less than the order of the denominator polynomial).
- Through polynomial long division, any rational function can be written as the sum of a polynomial and a strictly-proper rational function.
- A strictly-proper rational function can be expressed as a sum of lower-order rational functions, with such an expression being called a partial fraction expansion.

Section 14.1

## PFEs for First Form of Rational Functions

- Any rational function $F$ can be expressed in the form of

$$
F(v)=\frac{a_{m} v^{m}+a_{m-1} v^{m-1}+\ldots+a_{0}}{v^{n}+b_{n-1} v^{n-1}+\ldots+b_{0}}
$$

■ Furthermore, the denominator polynomial $D(v)=v^{n}+b_{n-1} v^{n-1}+\ldots+b_{0}$ in the above expression for $F(v)$ can be factored to obtain

$$
D(v)=\left(v-p_{1}\right)^{q_{1}}\left(v-p_{2}\right)^{q_{2}} \cdots\left(v-p_{n}\right)^{q_{n}},
$$

where the $p_{k}$ are distinct and the $q_{k}$ are integers.

- If $F$ has only simple poles, $q_{1}=q_{2}=\cdots=q_{n}=1$.
- Suppose that $F$ is strictly proper (i.e., $m<n$ ).
- In the determination of a partial fraction expansion of $F$, there are two cases to consider:

I $F$ has only simple poles; and
[ $F$ has at least one repeated pole.

- Suppose that the (rational) function $F$ has only simple poles.
- Then, the denominator polynomial $D$ for $F$ is of the form

$$
D(v)=\left(v-p_{1}\right)\left(v-p_{2}\right) \cdots\left(v-p_{n}\right),
$$

where the $p_{k}$ are distinct.

- In this case, $F$ has a partial fraction expansion of the form

$$
F(v)=\frac{A_{1}}{v-p_{1}}+\frac{A_{2}}{v-p_{2}}+\ldots+\frac{A_{n-1}}{v-p_{n-1}}+\frac{A_{n}}{v-p_{n}}
$$

where

$$
A_{k}=\left.\left(v-p_{k}\right) F(v)\right|_{v=p_{k}}
$$

- Note that the (simple) pole $p_{k}$ contributes a single term to the partial fraction expansion.
- Suppose that the (rational) function $F$ has at least one repeated pole.
- In this case, $F$ has a partial fraction expansion of the form

$$
\begin{aligned}
& F(v)=\left[\frac{A_{1,1}}{v-p_{1}}+\frac{A_{1,2}}{\left(v-p_{1}\right)^{2}}+\ldots+\frac{A_{1, q_{1}}}{\left(v-p_{1}\right)^{q_{1}}}\right] \\
&+\left[\frac{A_{2,1}}{v-p_{2}}+\ldots+\frac{A_{2, q_{2}}}{\left(v-p_{2}\right)^{q_{2}}}\right] \\
&+\ldots+\left[\frac{A_{P, 1}}{v-p_{P}}+\ldots+\frac{A_{P, q_{P}}}{\left(v-p_{P}\right)^{q_{P}}}\right]
\end{aligned}
$$

where

$$
A_{k, \ell}=\left.\frac{1}{\left(q_{k}-\ell\right)!}\left[\left[\frac{d}{d v}\right]^{q_{k}-\ell}\left[\left(v-p_{k}\right)^{q_{k}} F(v)\right]\right]\right|_{v=p_{k}}
$$

- Note that the $q_{k}$ th-order pole $p_{k}$ contributes $q_{k}$ terms to the partial fraction expansion.
- Note that $n!=(n)(n-1)(n-2) \cdots(1)$ and $0!=1$.


## PFEs for Second Form of Rational Functions

- Any rational function $F$ can be expressed in the form of

$$
F(v)=\frac{a_{m} v^{m}+a_{m-1} v^{m-1}+\ldots+a_{1} v+a_{0}}{b_{n} v^{n}+b_{n-1} v^{n-1}+\ldots+b_{1} v+1} .
$$

- Furthermore, the denominator polynomial
$D(v)=b_{n} v^{n}+b_{n-1} v^{n-1}+\ldots+b_{1} v+1$ in the above expression for $F(v)$ can be factored to obtain

$$
D(v)=\left(1-p_{1}^{-1} v\right)^{q_{1}}\left(1-p_{2}^{-1} v\right)^{q_{2}} \cdots\left(1-p_{n}^{-1} v\right)^{q_{n}}
$$

where the $p_{k}$ are distinct and the $q_{k}$ are integers.

- If $F$ has only simple poles, $q_{1}=q_{2}=\cdots=q_{n}=1$.
- Suppose that $F$ is strictly proper (i.e., $m<n$ ).
- In the determination of a partial fraction expansion of $F$, there are two cases to consider:
$11 F$ has only simple poles; and
- $F$ has at least one repeated pole.
- Suppose that the (rational) function $F$ has only simple poles.
- Then, the denominator polynomial $D$ for $F$ is of the form

$$
D(v)=\left(1-p_{1}^{-1} v\right)\left(1-p_{2}^{-1} v\right) \cdots\left(1-p_{n}^{-1} v\right)
$$

where the $p_{k}$ are distinct.

- In this case, $F$ has a partial fraction expansion of the form

$$
F(v)=\frac{A_{1}}{1-p_{1}^{-1} v}+\frac{A_{2}}{1-p_{2}^{-1} v}+\ldots+\frac{A_{n-1}}{1-p_{n-1}^{-1} v}+\frac{A_{n}}{1-p_{n}^{-1} v},
$$

where

$$
A_{k}=\left.\left(1-p_{k}^{-1} v\right) F(v)\right|_{v=p_{k}}
$$

- Note that the (simple) pole $p_{k}$ contributes a single term to the partial fraction expansion.
- Suppose that the (rational) function $F$ has at least one repeated pole.
- In this case, $F$ has a partial fraction expansion of the form

$$
\begin{aligned}
& F(v)=[ \left.\frac{A_{1,1}}{1-p_{1}^{-1} v}+\frac{A_{1,2}}{\left(1-p_{1}^{-1} v\right)^{2}}+\ldots+\frac{A_{1, q_{1}}}{\left(1-p_{1}^{-1} v\right)^{q_{1}}}\right] \\
&+\left[\frac{\left.A_{2,1}^{1-p_{2}^{-1} v}+\ldots+\frac{A_{2, q_{2}}}{\left(1-p_{2}^{-1} v\right)^{q_{2}}}\right]}{}\right. \\
&+\ldots+\left[\frac{A_{P, 1}}{1-p_{P}^{-1} v}+\ldots+\frac{A_{P, q_{P}}}{\left(1-p_{P}^{-1} v\right)^{q_{P}}}\right]
\end{aligned}
$$

where

$$
A_{k, \ell}=\left.\frac{1}{\left(q_{k}-\ell\right)!}\left(-p_{k}\right)^{q_{k}-\ell}\left[\left[\frac{d}{d v}\right]^{q_{k}-\ell}\left[\left(1-p_{k}^{-1} v\right)^{q_{k}} F(v)\right]\right]\right|_{v=p_{k}} .
$$

■ Note that the $q_{k}$ th-order pole $p_{k}$ contributes $q_{k}$ terms to the partial fraction expansion.

- Note that $n!=(n)(n-1)(n-2) \cdots(1)$ and $0!=1$.


## Part 15

## MATLAB

## MATLAB

- MATLAB is a proprietary programming language and numeric computing environment developed by MathWorks.
- MATLAB allows such things as: matrix computation, data visualization (e.g., 2-D and 3-D plots), system simulation, implementation of algorithms, and interfacing with code written in other programming languages.
- The base functionality of MATLAB can be extended through the use of numerous toolkits.
- MATLAB is available on many platforms, including Windows, macOS, Linux, and numerous other Unix variants.
- Extensive information on MATLAB (including detailed guides and manuals) is available from MathWorks web site:
- https://www.mathworks.com.
- Relative to general-purpose programming languages such as C++ and C, MATLAB is much easier to learn, and therefore is particular well suited for teaching purposes.
- The earliest version of MATLAB was a simple interactive matrix calculator (without an associated programming language) developed by Cleve Moler (then a professor at the University of New Mexico and later a cofounder of MathWorks).
- MATLAB stands for "matrix laboratory", with this name reflecting the roots of the software as a simple interactive matrix calculator.


## GNU Octave

- GNU Octave (or simply Octave) is a software package that is very similar in functionality to MATLAB.
- In fact, Octave is largely (but not fully) compatible with MATLAB.
- Unlike MATLAB, Octave is free software and is available under the GNU General Public License.
■ The Octave software was created by John W. Eaton, and was named after one of his professors, Octave Levenspiel.
- The Octave software is available for many platforms, including Windows, macOS, Linux, and other Unix variants, and can be obtained from its official web site:
- https://www.octave.org
- For anyone who does not have a license to use MATLAB, the Octave software provides an excellent alternative.


## Section 15.1

## Running MATLAB

## MATLAB Command-Line Interface

- MATLAB can be invoked using a command of the following form, where options denotes a list of zero or more options:
matlab options
- Some of the supported options are listed in the table below.

| Option | Description |
| :--- | :--- |
| -help | Print help information and exit. |
| -desktop | Enable the starting of the MATLAB desktop. |
| -nodesktop | Disable the starting of the MATLAB desktop and use <br> the terminal window for commands. |
| -display display | Set the X11 display to display. |
| -nodisplay | Do not use any X11 display. |
| -nosplash | Do not show the splash screen during startup. |
| -noFigureWindows | Disable the display of figure windows. |
| -batch command | Execute MATLAB command command upon startup <br> without the MATLAB desktop. |

- For example, to obtain help on the command-line interface of MATLAB, use the command:
matlab -help
- On Unix-based systems, the graphics output of MATLAB can typically be displayed on a different machine from the one running the MATLAB program (via the X11 protocol).
- This is typically used in a scenario where the user logs into remote machine and runs the MATLAB program on this remote machine with the graphics output being directed to the user's local machine.
- In order to take advantage of this functionality, the following software is required:
- a Secure Shell (SSH) client (that supports X11 forwarding); and
- an X11 server.
- With the necessary configuration in place, one can simply login via SSH to another machine with MATLAB installed and then run MATLAB.
- Any graphics windows opened by the MATLAB program will appear on the local machine via its X11 server.


## Secure Shell (SSH) Client Software

- Most Unix variants, including Linux and macOS, typically provide SSH client software (through a program that is usually called ssh ).
- Since recent versions of ChromeOS can run Linux in a container, the preceding comment also applies to ChromeOS. For information on containerized Linux in ChromeOS, see:
- https://chromeos.dev/en/linux

■ Some popular SSH client software for non-Unix-based platforms is as follows:

- PuTTY
- SSH client for Windows (and numerous Unix platforms as well)
- open source implementation of SSH and telnet
- web site: https://www.chiark.greenend.org.uk/~sgtatham/putty
- MobaXterm
- for Windows platform
- provides both an X11 server and a SSH client
- see slide on X11 server software for more details


## X11 Server Software

■ Most Unix variants, including Linux but excluding macOS, typically provide X11 server software.

- Some popular X11 server software for various platforms is as follows:
- MobaXterm
- X11 server for Windows platform
- enhanced terminal with X11 server and SSH client
- commercial software but Home Edition is free
- web site: https://mobaxterm.mobatek.net
- Xming
- X11 server for Windows platform
- open source
- web site: http://www.straightrunning.com/XmingNotes
- XQuartz
- X11 server for macOS platform
- web site: https://www.xquartz.org

■ For ChromeOS, use containerized Linux for an X11 server. See:

- https://chromeos.dev/en/linux


## MATLAB With Desktop



## MATLAB Without Desktop



## Section 15.2

## MATLAB Examples

```
1 figure
2 t = linspace (-1, 1, 1024);
3 hold on
4 plot(t, \(\sin \left(2{ }^{*}\right.\) pi * t), 'b');
5 plot(t, cos(2 * pi * t), 'r:', 'LineWidth', 2);
6 hold off
7 grid
8 legend('sin(2\pi\{\}t)', 'cos(2\pi\{\}t)');
9 print('output.pdf', '-dpdf', '-bestfit', '-opengl');
```


## Plot of Bivariate Function With Shading

```
1 figure
2 [x y] = meshgrid(linspace(-8, 8, 129));
3 r = sqrt(x .^ 2 + y .^ 2);
4 z = sinc(2 / pi * r);
5 surf(z);
6 shading interp;
7 light;
8 lighting phong;
9 material([0.9 0.5 0.5 1.0]);
10 colormap(parula);
11 print('output.pdf', '-dpdf', '-bestfit');
```



```
1 figure
2 d = @(x,y) sqrt(x .^ 2 + y .^ 2);
3 f = @(x,y) (d(x,y) <= 4) .* (4 - d(x,y));
4 fmesh(f, [-5 5 -5 5]);
5 print('output.pdf', '-dpdf', '-bestfit');
```



## Plot of Sequence

```
1 figure
2 n = 0 : 50;
3 h = stem(n, exp(-0.04 * n) .* sin(0.08 * pi * n), 'filled', 'b');
4 set(h(1), 'MarkerSize', 3)
5 xlabel('n');
6 ylabel('f(n)');
7 title('f(n) = e^{-0.04n} sin(0.08{\pi}n)');
8 print('output.pdf', '-dpdf','-bestfit', '-opengl');
```



## Histogram

```
1 clip = @(x, low, high) max(min(x, high), low);
2 figure
3 randn('seed', 5);
4 x = clip(67 + 10 * randn(250, 1), 0, 100);
5 histogram(x, 10);
6 title('Distribution of Student Grades');
7 xlabel('Grade (%)');
8 ylabel('Number of Students');
9 print('output.pdf', '-dpdf', '-bestfit');
```



## Pie Chart

```
figure
values = [65 18.5 9.5 3.2];
values = [values (100 - sum(values))];
names = {'Oxygen', 'Carbon', 'Hydrogen', 'Nitrogen', 'Other'};
explode = [00 0 0 0 0}]\mp@code{;
labels = cellstr([num2str(values', '%.1f%%' )])';
pie(values, explode, labels);
legend(names, 'Location', 'bestoutside');
title('Composition of Human Body')
print('output.pdf', '-dpdf', '-bestfit');
```



## Design of Analog Lowpass Filter

```
1 figure
2 [a b] = butter(30, 2000, 's');
3 [f, w] = freqs(a, b, linspace(0, 6000, 512));
4 subplot(2, 1, 1);
5 plot(w, abs(f));
6 xlabel('Frequency (rad/s)'); ylabel('Magnitude');
7 subplot(2, 1, 2);
8 plot(w, 180 / pi * unwrap(angle(f)));
9 xlabel('Frequency (rad/s)'); ylabel('Phase (\circ)');
10 print('output.pdf', '-dpdf', '-bestfit');
```



## Design of Digital Bandpass Filter

```
1 [a b] = firpm(63, [0 0.40 0.45 0.55 . 60 1], [0 0 1 1 0 0]);
2 freqz(a, b);
3 print('output.pdf', '-dpdf', '-bestfit');
```



## Solve Differential Equation: $y^{\prime \prime}=9 y, y(0)=a, y^{\prime}(0)=b$

```
1 syms y(t) a b
2 eqn = diff(y, t, 2) == 9 * y;
3 Dy = diff(y, t);
4 ics = [y(0) == a, Dy(0) == b];
5 ySol(t) = dsolve(eqn, ics);
6 S = evalc('pretty(ySol(t))');
7 disp(sprintf('y(t) = \n%s', s))
```

$$
\begin{aligned}
& y(t)= \\
& \exp \left(\left.\begin{array}{ll}
-3 & t) \quad\left|\begin{array}{cc}
- & - \\
2 & 6
\end{array}\right|+\exp (3 \mathrm{t}) \quad\left|\begin{array}{l}
- \\
2
\end{array}\right| \\
\hline
\end{array} \right\rvert\,\right.
\end{aligned}
$$

## Solve Nonlinear System: $x^{2}+y^{2}=a, 2 x-y=3$

```
1 syms x y a
2 [solx soly] = solve(x ^ 2 + y ^ 2 == a, 2 * x - y == 3, [x y]);
3 pretty([solx, soly]);
```



## Image Processing

```
figure
color = imread('peppers.png');
mono = rgb2gray(color);
edges = 1 - edge(mono, 'canny');
a = 8; blocky = imresize(imresize(color, 1/a), a, 'box');
subplot(2, 2, 1); set(gca, 'Position', [0 0.5 0.5 0.5])
imshow(color);
subplot(2, 2, 2); set(gca, 'Position', [0.5 0.5 0.5 0.5])
imshow (mono);
subplot(2, 2, 3); set(gca, 'Position', [0 0 0.5 0.5])
imshow(edges);
subplot(2, 2, 4); set(gca, 'Position', [0.5 0 0.5 0.5])
imshow(blocky);
print('output.pdf', '-dpdf', '-bestfit');
```



## Some Basic MATLAB Commands

- The quit or exit command terminate the MATLAB program.
- The clc command clears the command window.
- The clear command clears all variables.


## Section 15.3

## Accessing MATLAB Documentation

- The lookfor command searches for a string in the first comment line of the help text in all MATLAB files found on the MATLAB search path and prints each first comment line for which a match occurs.
- For example, to search for information that contains "butterworth" (as in Butterworth filter), use the following command:

lookfor butterworth

## lookfor butterworth Example



## Displaying MATLAB Documentation

- Two very helpful commands for accessing documentation about various aspects of MATLAB are the help and doc commands.
- The help command prints information to the command window about a specified item, such as a function, command, operator symbol, or toolbox.
- For example, the following command prints information about the abs function:
help abs
- For example, the following command prints information about the . * operator:
help .*
- The doc command opens the help browser (if not already open) and uses the browser to display the documentation for the specified item, such as a function, command, operator symbol, method, class, or toolbox.
- For example, the following command displays documentation for the abs function in the help browser:
doc abs


## Command-Window Help Information for abs Function



## Help Browser Information for abs Function



Section 15.4

## Comments and Variables

## Comments

- Two styles of comments are provided:

1 block comments
2 short comments

- A block comment:
- starts with a line that contains a percent character (i.e., "o") immediately followed by a left-brace character (i.e., "\{") and has no other characters except spaces and tabs;
- ends with a line that contains a percent character immediately followed by a right-brace character (i.e., "\}") with no other characters except spaces and tabs.
- A short comment starts with a percent character (i.e., "응) that does not start a block comment and continues until the end of line.
- Some examples of comments are given in the following code fragment:

```
% This short comment continues until the end of line.
disp(sin(42)); % Print the sine of 42.
%{
This is a block comment.
This comment spans multiple lines.
%}
```

- An identifier is a name for an entity such as a variable or function.
- Identifiers are case sensitive and may consist of uppercase and lowercase letters, underscores, and digits, but the first character cannot be a digit or an underscore.
- Although an identifier can be arbitrarily long, only the first $n$ characters are significant, where $n$ depends on the particular version of MATLAB being used.
- The namelengthmax function can be used to query the precise value of $n$.
- Some examples of valid and invalid identifiers are given below.

| Name | Valid Identifier |
| :--- | :--- |
| foo | yes |
| foo2 | yes |
| f_o_o | yes |
| _foo | no (due to leading underscore) |
| 2foo | no (due to leading digit) |
| foo\$ | no (due to dollar sign) |

- A reserved keyword (which has special meaning to MATLAB) cannot be used as an identifier (e.g., variable or function name).
- The iskeyword function can be used to obtain a list of all reserved keywords.
- The reserved keywords in MATLAB include those in the list below.

| break | end | persistent |
| :--- | :--- | :--- |
| case | for | return |
| catch | function | spmd |
| classdef | global | switch |
| continue | if | try |
| else | otherwise | while |
| elseif | parfor |  |

■ A comma or semicolon separates statements (or parts of statements) within a single line of code.

- For example, an if statement can be written on a single line using an approach like the following:

$$
\text { if } x==0 ; x=1 \text {; end }
$$

- In the case of a comma, the result computed from the statement (if any) is printed.
- In the case of a semicolon, the result is not printed.
- In order to avoid overly verbose program output, statements are typically ended by a semicolon instead of a comma.
- For example, the following line of code will not generate any output:

$$
x=1 ; y=2 \text { * } x ;
$$

- In contrast, the following line of code will print the values computed for $x$ and y :

$$
x=1, y=2 * x
$$

- Sometimes, it is desirable (for reasons of readability) to split a logical line of the source code across multiple physical lines in a program file.
- Some language constructs, however, must appear on the same logical line of the source code.
- For this reason, the language provides a mechanism for saying that a logical line of source code continues on the next physical line in the program file.
- In particular, an ellipsis consisting of at least three consecutive period (".") characters can be used to continue a logical line of the source code on the next physical line in the file.
- Example:

$$
\begin{aligned}
x & =1+2+3+\ldots \\
& 4+5+6
\end{aligned}
$$

## Variables

- A variable does not need to be declared prior to being used.
- A variable simply comes into existence when it is first assigned a value.
- The type of a variable is determined at run time based on the value assigned to the variable.
■ Essentially, every variable in MATLAB is deemed to be an array, where a scalar is simply a single-element array.
- Elements of an array are essentially complex numbers, but may be constrained to assume only real or integer values.


## Arrays

- Array indices start numbering at 1 (not 0 ).
- The specification of an array starts with a left (square) bracket and continues to the corresponding right (square) bracket, where these brackets need not be on the same line.
- The elements for each row are specified contiguously.
- Spaces or commas separate elements in the same row of an array.
- Semicolons or newlines separate elements on different rows.
- For example, the following specifies an array with 2 rows and 3 columns, where the first row has the elements 1,2 , and 3 , and the second row has the elements 4,5 , and 6 :

$$
\left[\begin{array}{llllll}
1 & 2 & 3 ; & 4 & 6
\end{array}\right]
$$

- The preceding matrix could also be written as each of the following:

```
[1, 2, 3; 4, 5, 6]
[1 2 3
4 5 6]
```


## Vectors With Equally-Spaced Elements

- MATLAB provides a compact syntax for specifying a vector with equally-spaced elements.
- In particular, an expression of the following form is employed:
start : step : end
- The above expression is equivalent to a row vector with its first element equal to start and each of the subsequent elements increasing in value by step (where step may be negative) until the value would surpass end.
■ The notation "start : 1 : end" can be further abbreviated to "start : end".
- Some examples of this abbreviated notation are shown below.

| Abbreviated Form | Long Form |
| :---: | :---: |
| 1: 4 | $\left[\begin{array}{lllll}1 & 2 & 3\end{array}\right]$ |
| 0: 0.2 : 1 | $\left[\begin{array}{lllllllll}0 & 0.2 & 0.4 & 0.6 & 0.8 & 1\end{array}\right]$ |
| 1 : -1 : -2 | $\left[\begin{array}{lllll}1 & 0 & -1 & -2\end{array}\right]$ |
| 0: 10: 25 | [0 10 20] |
| -1.5 : -1 : -4 | $\left[\begin{array}{llll}-1.5 & -2.5 & -3.5\end{array}\right]$ |

## Array Subscripting

- One or more elements of an array $x$ can be accessed by specifying the rows and columns in which the elements are contained using the subscripting operator.
- In particular, $x$ (rowspec, colspec) is the array consisting of the elements of $x$ that are in the rows specified by rowspec and columns specified by colspec.
- Here, rowspec is either a vector containing row indices or the special value ":" which means "all rows".
- Similarly, colspec is either a vector containing column indices or the special value ":" which means "all columns".
- Example:
- Suppose that $\mathrm{x}=\left[\begin{array}{llll}1 & 2 & 3 ; & 4 \\ \text { 6; } 78 & 9\end{array}\right]$.
- $x(1,3)$ is 3 (i.e., the row 1 column 3 element of $x$ ).
- $x(:, 1)$ is $[1 ; 4$; 7] (i.e., the 1 st column of $x$ ).
- $x(1,:)$ is $\left[\begin{array}{lll}1 & 2 & 3\end{array}\right]$ (i.e., the 1 st row of $\left.x\right)$.
- $x(2: 3,:)$ is $\left[\begin{array}{llll}4 & 5 & 6 ; & 8\end{array}\right]$ (i.e., the 2 nd and 3rd rows of $\left.x\right)$.


## Array Subscripting By Column-Major Index

- One or more elements of an array $x$ can also be accessed by specifying their index when numbered in column-major order (which orders elements first by column and then by row).
- In particular, $x$ (indspec) is an array that contains the elements of $x$ having the indices specified by indspec.
- The dimensions of $x$ (indspec) match the dimensions of indspec if $x$ is a matrix (i.e., not a row or column vector); otherwise, $x$ (indspec) is either a row or column vector to match $x$.
- Here, indspec is either an array containing element indices or the special value ":" which is equivalent to all indices of $x$ specified in a column vector.
- Example:
- Suppose that $\mathrm{x}=\left[\begin{array}{lllllllll}1 & 4 & 7 ; & 2 & 5 & 8 ; & 3 & 6 & 9\end{array}\right]$ (i.e., $\left[\begin{array}{ccc}1 & 4 & 7 \\ 2 & 5 & 8 \\ 3 & 9 & 9\end{array}\right]$.
- $x(3)$ is 3 (i.e., the 3rd element of $x$ in column-major order).
- $x\left(\left[\begin{array}{ll}1 & 4\end{array}\right]\right)$ is [14 7] (i.e., the 1st, 4th, and 7th elements of $x$ in column-major order).
- $x(:)$ is $[1 ; 2$; 3; 4; 5; 6; 7; 8; 9] (i.e., all elements of $x$ in column-major order placed in a column vector).


## Array-Subscripting Examples

Suppose that a is a $10 \times 10$ matrix and x is $10 \times 1$ vector.

| Expression | Meaning |
| :--- | :--- |
| $\mathrm{a}(1,:)$ | first row of $a$ |
| $\mathrm{a}(:, 1)$ | first column of $a$ |
| $\mathrm{a}(1: 50)$ | first 50 elements of a arranged in a row vector |
| $\mathrm{a}(1: 10)$ | first 10 elements of a arranged in a row vector (i.e., <br> the first column of a$)$ |
| $\mathrm{a}(1: 2: 10,:)$ | odd-indexed rows of a |
| $\mathrm{a}(:, 2: 2: 10)$ | even-indexed columns of a |
| $\mathrm{a}(1: 5,:)$ | rows 1 to 5 of a |
| $\mathrm{a}(:, 6: 10)$ | columns 6 to 10 of a |
| $\mathrm{a}(1: 2,9: 10)$ | submatrix consisting of elements that are in rows <br> 1,2 and also in columns 9,10 |
| $\mathrm{x}(1: 3)$ | first three elements of $\mathrm{x}($ arranged as a row or col- <br> umn vector to match x$)$ |

## Strings

- A string is simply an array whose elements are characters.
- A string literal starts with a single-quote character ("' ") and continues until the next single-quote character that is not part of a pair of consecutive single-quote characters.
- To include a single-quote character in a string, two consecutive single-quotes characters can be used.
- Some examples of string literals are as follows:

```
'Hello, World!'
'This is a test\n'
'Jane''s computer'
```

■ The following denotes a string that consists of one single-quote character:

■ For a string s , the expression s (length (s) : -1 : 1) corresponds to the string reversed (i.e., backwards).

## Section 15.5

## Operators and Expressions

## Operators

- Operators (such as addition and subtraction) are used to perform calculations.
- An operator that takes one operand is said to be unary.
- For example, the expression "-x" employs the unary minus (i.e., negation) operator, where the minus operator has the single operand $x$.
- An operator that takes two operands is said to be lbinary.

■ For example, the expression " $x$ - $y$ " employs the binary minus (i.e., subtraction) operator, where the operator has two operands x and y .

| Symbol | Description |
| :--- | :--- |
| + | unary plus (i.e., identity) and binary plus (i.e., addition) |
| - | unary minus (i.e., negation) and binary minus (i.e., subtraction) |
| $\star$ | multiplication |
| $/$ | right division (a / b means a $\left.{ }^{\star} \mathrm{b}{ }^{\wedge}(-1)\right)$ |
| $\backslash$ | left division (a $\backslash \mathrm{b}$ means a $\left.\wedge^{(-1)}{ }^{\star} \mathrm{b}\right)$ |
| $\wedge$ | exponentiation |
| $\prime$ | conjugate transpose |
| .$^{\star}$ | elementwise multiplication |
| .$/$ | elementwise division |
| .$^{\wedge}$ | elementwise exponentiation |
| $\cdot^{\prime}$ | transpose |

## Elementwise Versus Nonelementwise Operators

| First Operand | Second Operand | Operator | Result |
| :---: | :---: | :---: | :---: |
| $\begin{array}{rr} {[1} & -1 \\ 0 & 1] \end{array}$ | $\begin{gathered} {[1} \\ 3] \end{gathered}$ | * | $\begin{gathered} {[-2} \\ 3] \end{gathered}$ |
| $\begin{array}{rr} {[1} & -1 \\ 0 & 1] \end{array}$ | $\left[\begin{array}{ll} 1 & 2 \\ 3 & 4 \end{array}\right]$ | * | $\begin{array}{rc} {[-2} & -2 \\ 3 & 4] \end{array}$ |
| $\begin{array}{rr} {[1} & -1 \\ 0 & 1] \end{array}$ | $\left[\begin{array}{ll} 1 & 2 \\ 3 & 4 \end{array}\right]$ | .* | $\begin{array}{rr} {[1} & -2 \\ 0 & 4] \end{array}$ |
| $\left[\begin{array}{ll}1 & -1\end{array}\right]$ | $\left[\begin{array}{ll}2 & 2\end{array}\right]$ | * | $\left[\begin{array}{ll}2 & -2\end{array}\right]$ |

## Operator Precedence

- When an expression involves multiple operators, the order in which the operators are applied becomes significant.
- Operators are applied in order of decreasing precedence (i.e., priority).
- Operators with the same precedence are evaluated left to right.
- Parentheses can be used to force a particular ordering of operations.
- The precedence of operators in MATLAB is shown in the table below.

| Precedence Level | Operators |
| :---: | :---: |
| 1 (highest) | () |
| 2 | . ${ }^{\prime}$.^ ^ |
| 3 | unary+ unary- ~ |
| 4 | .* * ./ . \ / \} |
| 5 | binary+ binary- |
| 6 | : |
| 7 | \ll= \gg = == ~= |
| 8 |  |
| 9 | \| |
| 10 |  |
| 11 (lowest) | \|| |

## Operator Precedence/Associativity Example

| Expression | Fully-Parenthesized Expression |
| :---: | :---: |
| + b * x ^ 2 | $\mathrm{a}+(\mathrm{b}$ * (x ^ 2) ) |
| + b .* x .^ 2 | a + (b.* (x .^ 2) ) |
| + b | $(\mathrm{a}+\mathrm{b})+\mathrm{c}$ |
|  | $((a>b) \& \&(c<d))\|\mid(e==f)$ |
| $\mathrm{x}(\mathrm{m}+1: 2$ * $\mathrm{n}+1)$ | $x((m+1):((2 \times n)+1))$ |
| x * + x | (-x) * ( +x ) |
| $0<=\mathrm{x}$ \&\& $\mathrm{x}<=1$ | (0<ex) \&\& (x < $=1$ ) |
| $0<=x$ <= 1 (probably bug) | ( $0<=\mathrm{x}$ ) <= 1 (which is always 1) |

- A relational operator is an operator that tests an ordering or equivalence/equality relationship between two values.
- The following relational operators are provided by MATLAB:

| Symbol | Description |
| :--- | :--- |
| $<$ | less than |
| $<=$ | less than or equal |
| $>$ | greater than |
| $>=$ | greater than or equal |
| $==$ | equal |
| $\sim=$ | not equal |

## Logical Operators

- A logical operator is an operator that computes a primitive Boolean function.
- MATLAB considers any nonzero number to be true and zero to be false.
- The following logical operators are provided by MATLAB:

| Symbol | Description |
| :--- | :--- |
| $\&$ | elementwise AND |
| $I$ | elementwise OR |
| $\sim$ | NOT |
| $\& \&$ | short-circuit AND |
| $\\|\\|$ | short-circuit OR |

- The elementwise AND, elementwise OR, and NOT operators allow nonscalar (i.e., matrix/vector) operands, whereas the short-circuit AND and short-circuit OR operators require scalar operands.


## Examples of Expressions Involving Relational/Logical Operators

| Variable | Value |  |  |
| :--- | :--- | :--- | :--- |
| a | $\left[\begin{array}{lllll}0 & 1 & 2 & 3 & 4\end{array}\right]$ |  |  |
| b | $\left[\begin{array}{lllll}4 & 3 & 2 & 1 & 0\end{array}\right]$ |  |  |


| Expression | Value |
| :---: | :---: |
| $\mathrm{a}={ }^{\text {a }}$ b | $\left[\begin{array}{lllll}0 & 0 & 1 & 0 & 0\end{array}\right]$ |
| $\mathrm{a} \sim=\mathrm{b}$ | $\left[\begin{array}{lllll}1 & 1 & 0 & 1 & 1\end{array}\right]$ |
| $a>b$ | $\left[\begin{array}{lllll}0 & 0 & 0 & 1 & 1\end{array}\right]$ |
| $a>1$ | $\left[\begin{array}{lllll}0 & 0 & 1 & 1 & 1\end{array}\right]$ |
| $\mathrm{a}<\mathrm{b}$ | $\left[\begin{array}{lllll}1 & 1 & 0 & 0 & 0\end{array}\right]$ |
| a >= 1 \& a <= 3 | $\left[\begin{array}{lllll}0 & 1 & 1 & 1 & 0\end{array}\right]$ |
| $a<1 \mid a>3$ | $\left[\begin{array}{lllll}1 & 0 & 0 & 0 & 1\end{array}\right]$ |
| $\sim \mathrm{a}$ | $\left[\begin{array}{lllll}1 & 0 & 0 & 0 & 0\end{array}\right]$ |
| $0<=$ a \& a <= 1 | $\left[\begin{array}{lllll}1 & 1 & 0 & 0 & 0\end{array}\right]$ |
| $0<=a<=1$ (probably bug) | $\left[\begin{array}{lllll}1 & 1 & 1 & 1 & 1\end{array}\right]$ (always all ones) |

## Short-Circuited Evaluation of Logical Expressions

- The second operand of the short-circuit AND operator ( $\& \&$ ) is not evaluated if the first operand is false, since the second operand cannot affect the result in this case.
- The second operand of the OR operator ( $\mid \mathrm{l}$ ) is not evaluated if the first operand is true, since the second cannot affect the result in this case.
■ The above behavior is known as short-circuit evaluation.
- Example:
- Suppose that x and y are initialized as $\mathrm{x}=0$ and $\mathrm{y}=1$.
- If MATLAB is asked to evaluate the logical expression $\mathrm{x} \sim=0 \& \& \mathrm{y}<1 / \mathrm{x}$, only $\mathrm{x} \sim=0$ will be evaluated, since the result must be false regardless of the value of $y<1 / x$.
- Note that short circuit evaluation guarantees division by zero cannot occur during the evaluation of $x \sim=0 \& \& y<1 / x$.
- If MATLAB is asked to evaluate the logical expression $x \sim=42| | y>0$, only $\mathrm{x} \sim=42$ will be evaluated, since the result must be true regardless of the value of $\mathrm{y}>0$.

Section 15.6

## Selection Constructs

## if Statements

- An if statement allows groups of statements to be conditionally executed, and has a number of variants.
- The simplest variant (i.e., the if variant) has the form:

```
if expression
    statements
end
```

- If the expression expression is true, then the statements statements are executed.
- Example:

```
if x < 0
    disp('x is negative');
end
```

■ The next simplest variant of the if statement (i.e., the if-else variant) has the form:

```
if expression
    statements}
else
    statements}
end
```

- If the expression expression ${ }_{1}$ is true, then the statements statements $_{1}$ are executed; otherwise, the statements statements $s_{2}$ are executed.
- Example:

```
    if x < 0
        disp('x is negative');
    else
    disp('x is not negative');
    end
```

- Finally, the most general variant of the if statement has the following form (where the elseif and else clauses are optional):

```
if expression
    statements }
elseif expression}
    statements}
elseif expressionn-1
    statements}\mp@subsup{n}{-1}{
```

else
statements $_{n}$
end

- For each if/elseif clause in order of appearance, the expression expression $_{i}$ is evaluated, and if this expression is true, the statements statements $_{i}$ are executed and the remainder of the if statement is skipped.
- If none of the expressions in the if/elseif clauses are true and an else clause is present, the statements statements $_{n}$ in this clause are executed.


## if Statement Example

```
    \% Read a real number and then print information
    \% about its sign.
    x = input('Enter a real number: ');
    if \(\mathrm{x}>0\)
        disp('x is strictly positive');
    elseif \(\mathrm{x}<0\)
        disp('x is strictly negative');
    else
        disp('x is zero');
    end
```


## switch Statements

- The switch statement provides another means to conditionally execute groups of statements.
- The general form of a switch statement is as follows (where the otherwise clause is optional):
switch expression
case test_expression ${ }_{1}$
statements ${ }_{1}$
case test_expression 2
statements $_{2}$
case test_expression $n-1$ statements $_{n-1}$
otherwise
statements $_{n}$
end
- The first test expression, say test_expression ${ }_{k}$, matching the expression expression has its corresponding statements statements $_{k}$ executed.
- If none of the test expressions match the expression expression and an otherwise clause is present, the statements statements $n$ in this clause are executed.
- The switch expression expression must be either a scalar or string.

```
% Get a real value and print some information about it.
n = input('Enter real number: ');
switch mod(n, 2)
case 0
    disp('number is even integer');
case 1
    disp('number is odd integer');
case {0.5, 1.5}
    disp('number is odd multiple of one half');
otherwise
    disp('number is not multiple of one half');
end
```


## switch Statement Example [sting]

```
\% Get the type of a polygon and then print the number of
\% sides that the polygon has.
type = input('Enter polygon type: ', 's');
switch lower (type)
case 'triangle'
num_sides = 3;
case 'quadrilateral'
    num_sides = 4;
case 'pentagon'
    num_sides = 5;
case 'hexagon'
    num_sides = 6;
case 'heptagon'
    num_sides = 7;
case 'octagon'
    num_sides = 8;
case 'nonagon'
    num_sides = 9;
case 'decagon'
    num_sides = 10;
otherwise
    error('unknown polygon type');
    end
    fprintf('The polygon has \%d sides. \(\mathrm{n}^{\prime}\), num_sides);
```

Section 15.7

## Looping Constructs

- The for statement allows a group of statements to be repeated a fixed number of times.
- A for statement has the general form:

```
for variable = array
        statements
    end
```

- The statements statements are executed once for each column in the array array, where the variable variable is set to the corresponding array column each time.
- The following code fragment prints the integers $1,3,7$, and 9 :

```
for k = [lllll
    disp(k);
end
```


## for Statement Example

Code

```
1 fprintf('Degrees Radians \(\left.{ }^{\prime} n^{\prime}\right)\);
2 for theta_deg = -180 : 30 : 180
3 theta_rad = theta_deg * pi / 180;
fprintf('\%7.1f \(\% 7.4 f \backslash n^{\prime}\), theta_deg, theta_rad);
end
```

Output

| Degrees | Radians |
| ---: | ---: |
| -180.0 | -3.1416 |
| -150.0 | -2.6180 |
| -120.0 | -2.0944 |
| -90.0 | -1.5708 |
| -60.0 | -1.0472 |
| -30.0 | -0.5236 |
| 0.0 | 0.0000 |
| 30.0 | 0.5236 |
| 60.0 | 1.0472 |
| 90.0 | 1.5708 |
| 120.0 | 2.0944 |
| 150.0 | 2.6180 |
| 180.0 | 3.1416 |

## while Statements

- The while statement allows a group of statements to be executed an indefinite number of times.
- A while statement has the general form:
while expression
statements
end
- The statements statements are executed repeatedly as long as the condition expression is true.
- The following code fragments prints random numbers obtained via the rand function until a number not exceeding 0.1 is obtained:

```
\(\mathrm{x}=\) rand;
while x > 0.1
    disp(x);
    \(\mathrm{x}=\) rand;
end
```


## break Statements

- Sometimes, it may be necessary to terminate a loop by code executing in the body of the loop.
- This can be accomplished using a break statement.
- Typically, a break statement is used to terminate a loop from code that is somewhere in the middle of the loop body.
- The following code uses a break statement in order to terminate a while loop when a particular condition is satisfied:

```
values = [];
while 1
    x = rand;
    if x > 0.9
            break
    end
    values = [values x];
end
disp(values);
```

- Sometimes, it may necessary to start the next iteration of a loop without completing the remainder of the code in the loop body.
- This is accomplished using a cont inue statement.
- The following code fragment uses a continue statement to skip over the processing of any zero elements in the array $a$ :

```
\(a=\left[\begin{array}{lllll}1 & 0 & 3 & 2 & 0\end{array}\right] ;\)
for \(i=a\)
    if i == 0
        \% Skip over the processing of a
        \% zero element in the array.
        continue
    end
    \% Process the nonzero array element.
        disp(i);
end
```


## Section 15.8

## Using Functions

## Calling Functions

- A function that can be called with no arguments can be invoked by specifying either its name or its name followed by a pair of empty parentheses (with the former usually being preferred since it is less verbose).
- For example, to invoke the built-in $j$ function (which takes no arguments and returns $\sqrt{-1}$ ), we would write either " $j$ " or " $j()$ ".
- A function that requires one or more arguments is called by placing the list of arguments to the function (separated by commas) in parentheses (i.e., round brackets) after the name of the function.
- Examples:
- To invoke the sin function with the single argument 42 , we would write "sin(42)".
- To invoke the atan2 function with the arguments y and x , we would write "atan2 $(y, x)$ ".

■ MATLAB has many built-in (i.e., automatically predefined) functions.

- Although a built-in function can be overridden by a user-defined function or variable, doing this can often lead to confusing code.
- Consequently, it is recommended that the overriding of built-in functions be avoided in most cases.
- One exception to this rule might be overriding the built-in i and $j$ functions in code that does not use complex arithmetic, since this is unlikely to lead to confusion and $i$ and $j$ are often desirable names for looping variables.

| Name | Description |
| :--- | :--- |
| pi | $\pi$ |
| i | $\sqrt{-1}$ |
| j | $\sqrt{-1}$ |
| nan | not-a-number (NaN) |
| inf | infinity |
| ans | last expression evaluated that was not assigned to vari- <br> able |
| date | date |
| clock | wall clock |
| realmin | smallest usable positive real number |
| realmax | largest usable positive real number |

## Elementary Math Functions

| Name | Description |
| :--- | :--- |
| abs | magnitude of complex number |
| angle | principal argument of complex number |
| imag | imaginary part of complex number |
| real | real part of complex number |
| conj | conjugate of complex number |
| sign | signum function |
| rem | remainder (with same sign as dividend) |
| mod | remainder (with same sign as divisor) |


| Rounding Functions |  |
| :--- | :--- |
| Name Description <br> round round to nearest integer <br> fix round towards zero <br> floor round towards $-\infty$ <br> ceil round towards $\infty$ |  |

Basic Statistical Functions

| Name | Description |
| :--- | :--- |
| min | minimum value |
| max | maximum value |
| mean | mean value |
| std | standard deviation |
| median | median value |

Sum and Product Functions

| Name | Description |
| :--- | :--- |
| sum | sum of elements |
| prod | product of elements |
| cumsum | cumulative sum of elements |
| cumprod | cumulative product of elements |

## Exponential and Logarithmic Functions

| Name | Description |
| :--- | :--- |
| exp | exponential function |
| $\log$ | natural logarithmic function |
| $\log 10$ | base-10 logarithmic function |
| sqrt | square root function |

Trigonometric Functions

| Name | Description |
| :--- | :--- |
| $\sin$ | sine function |
| $\cos$ | cosine function |
| $\tan$ | tangent function |
| $\operatorname{asin}$ | arcsine function |
| $\operatorname{acos}$ | arccosine function |
| atan | arctangent function |
| $\operatorname{atan2} 2$ | two-argument form of arctangent function |

Radix Conversion Functions

| Name | Description |
| :--- | :--- |
| dec2bin | convert decimal to binary |
| bin2dec | convert binary to decimal |
| dec2hex | convert decimal to hexadecimal |
| hex2dec | convert hexadecimal to decimal |
| dec2base | convert decimal to arbitrary radix |
| base2dec | convert arbitrary radix to decimal |

## Other Math Functions

| Name | Description |
| :--- | :--- |
| sinc | normalized sinc function $\left(\right.$ i.e., $\left.f(t)=\frac{\sin (\pi t)}{\pi t}\right)$ |
| polyval | evaluate polynomial |
| cart2pol | Cartesian-to-polar coordinate conversion |
| pol2cart | polar-to-Cartesian coordinate conversion |

Matrix Computation Functions

| Name | Description |
| :--- | :--- |
| expm | compute matrix exponential |
| det | compute determinant of matrix |
| inv | compute inverse of matrix |
| eig | compute eigenvalues and eigenvectors of matrix |
| svd | compute singular-value decomposition (SVD) of matrix |
| lu | compute lower-upper (LU) decomposition of matrix |
| qr | compute QR decomposition of matrix |


| Name | Description |
| :--- | :--- |
| size | get row vector containing size of array in each di- <br> mension |
| width | get number of columns in array |
| height | get number of rows in array |
| length | get size of largest array dimension (e.g., for vector, <br> yields vector's length) |
| ndims | get number of array dimensions (always at least <br> two) |
| numel | get number of elements in array |
| isrow | tests if array is row vector |
| iscolumn | tests if array is column vector |
| isvector | tests if array is row or column vector |
| ismatrix | tests if array is matrix |
| isempty | tests if array is empty (i.e., contains no elements) |

## Array-Size Functions Example

```
x = input('Enter a matrix/vector: ');
if numel(x) < 2
    error('invalid matrix/vector (i.e., scalar or empty)');
end
if isrow(x)
    disp('array is row vector');
elseif iscolumn(x)
    disp('array is column vector');
elseif ismatrix(x)
    disp('array is matrix');
end
fprintf('array contains %d elements\n', numel(x));
fprintf('array elements in row-major order:\n');
for i = 1 : height(x)
            for j = 1 : width(x)
                fprintf("%.15f\n", x(i, j));
            end
end
```


## Special Matrix/Vector Functions

| Name | Description |
| :--- | :--- |
| eye | identity matrix |
| ones | matrix of ones |
| zeros | matrix of zeros |
| diag | diagonal matrix |
| rand | random matrix |
| linspace | vector with linearly spaced elements |
| logspace | vector with logarithmically spaced elements |

Basic Array Manipulation Functions

| Name | Description |
| :--- | :--- |
| rot 90 | rotate array by 90 degrees |
| fliplr | flip array horizontally |
| flipud | flip array vertically |
| reshape | change array dimensions |

## Array-Functions Example

| Expression | Result |  |
| :---: | ---: | :---: |
| zeros $(2,3)$ | $\left[\begin{array}{lll}0 & 0 & 0 \\ 0 & 0 & 0\end{array}\right]$ |  |
| ones $(2,3)$ | $\left[\begin{array}{lll}1 & 1 & 1 \\ 1 & 1 & 1\end{array}\right]$ |  |
|  | $\left[\begin{array}{lll}1 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & 1\end{array}\right]$ |  |
| eye (3) | $\left[\begin{array}{llll}0 & 1 & 2 & 3\end{array}\right]$ |  |


| Name | Description |
| :--- | :--- |
| any | any element nonzero |
| all | all elements nonzero |
| find | find nonzero elements |
| exist | check if variable/function exists |
| isfinite | detect finite values |
| isinf | detect infinities |
| isnan | detect NaNs |
| isempty | detect empty matrices |
| isstr | detect strings |
| strcmp | compare strings |

## Logical Functions Example

```
\(1 a=\left[\begin{array}{llllll}1 & 2 & 3 ; & 4 & 6 ; 7 & 9\end{array}\right] ;\)
2 if all(a > 0)
    disp('All matrix elements are positive.');
end
if \(\sim \operatorname{any}(a==0)\)
        disp('All matrix elements are nonzero.');
end
if all(real(a) == a)
disp('All matrix elements are real.');
end
i = find (a >= 8);
disp(['The following matrix elements are ', ...
    'greater than or equal to 8:']);
disp(a(i));
```


## Section 15.9

## Scripts and Functions

## M Files

- MATLAB code can be placed in a file, called an M file, for later execution.
- An M file must have file name extension ". m " (which is where term " M file" originates).
- M files are used for scripts and specifying user-defined functions.
- Depending on how the code in an M file is structured it can correspond to either a script or user-defined function.
- By using the MATLAB identifier name, one can execute the code stored in an M file named "name.m".

■ For example, the MATLAB identifier my_demo would refer to the MATLAB code in the M file named my_demo.m (assuming that the identifier my_demo is not used as a variable name).

## Scripts

- An M file that does not define a function is referred to as a script.
- Executing code in a script behaves as if the contents of the script were pasted into the place where script is being referenced.
- Since a script is executed in the context of its invoker, any changes made to variables are visible to the invoker.

Contents of file greet.m

```
1 % Print a greeting.
2 disp('Hello, World.');
```

- The code in the file greet.m shown above can be executed by simply typing the following in the MATLAB command window: greet
- Executing the code will result in the disp function being invoked to print the message "Hello, World!".


## Script Example: demo

Contents of file demo.m

$$
\begin{aligned}
& 1 \quad t=[0: .1: 8] ; \\
& 2 \mathrm{f}=\sin (\mathrm{t}) ; \\
& 3 \mathrm{plot}(\mathrm{t}, \mathrm{f}) ;
\end{aligned}
$$

- The code in the file demo.m shown above can be executed by simply typing the following in the MATLAB command window:


## demo

- Recall that a script is executed in the context of the invoker so any changes made to variables by the script are visible to the invoker.
- For example, if the code in the above script is invoked as follows, the changes to the variables $t$ and $f$ are visible to the invoker:

```
clear t f % Make t and f undefined.
demo % execute code in demo.m script
disp(t); % OK: t was defined by script
disp(f); % OK: f was defined by script
```

- The MATLAB search path is a list of directories that MATLAB examines in order to locate various files, such as MATLAB program files (i.e., M files) and files associated with toolboxes.
- When trying to locate a file, MATLAB will look for the file in each directory in the search path in order until the file is found.
- The order in which directories appear in the search path is important, since the directories are searched in that order.
- Normally, the search path contains the current working directory (which appears first) as well as any directories associated with files that are part of the MATLAB software installation.
- When MATLAB encounters an identifier, it must decide the specific entity to which the identifier refers (such as a variable, built-in function, or script or function specified in an M file).
- To do this for the identifier id, MATLAB employs an algorithm consisting of the following steps (in order):

1 If a user-defined (i.e., not built-in) variable exists with id as its name, id is deemed to refer to this variable.
2 Otherwise, if $i d$ has a corresponding M file (i.e., a file named id.m) in current directory, id is deemed to refer to the function or script associated with that M file.
3 Otherwise, if $i d$ has a corresponding M file (i.e., a file named $i d . \mathrm{m}$ ) in any directory on the MATLAB search path, id is deemed to refer to the function or script associated with the first such $M$ file found on this path.
4 Otherwise, if a built-in variable/function exists with id as its name, id is deemed to refer to this variable/function.

- If multiple entities are associated with the same identifier id, the order of the above steps becomes important in deciding which entity will be referenced by MATLAB when id is used.


## Naming Conflicts

- Generally, it is advisable to avoid naming variables and M files in ways that could result in conflicts.

■ When multiple entities are associated with the same identifier, one entity can hide (i.e., override) another.

- The use of an identifier that would conflict with any built-in function should be avoided in most cases.
- A user-defined variable will hide a function of the same name. For example:

```
sin = 42
    % can no longer call sin function
    sin(pi) % ERROR: trying to subscript with index pi
```

- An M file in the current directory will hide all other M files of the same name in the MATLAB search path.
- All of the directories in the MATLAB search path (i.e., all of the directories in which MATLAB will look for M-file scripts) can be printed with the path command:
path
- The working directory for MATLAB can be changed to dirname using the cd command:
cd dirname
- The current working directory can be printed with the pwd command:
pwd
- A function definition is introduced using the following syntax:

$$
\text { function }\left[\text { out }_{1}, \text { out }_{2}, \ldots, \text { out }_{n}\right]=\text { func }\left(\text { in }_{1}, \text { in }_{2}, \ldots, \text { in }_{m}\right)
$$

- This indicates that the function has:
- the name func,
- the $m$ input arguments $i n_{1}, i n_{2}, \ldots, i n_{m}$, and
- the $n$ output arguments out $_{1}$, out $_{2}, \ldots$, out $_{n}$.
- The body of a function extends from the function directive to a corresponding end directive or, if no such end directive is present, the end of the file.
- The code in a function executes until either the end of the function is reached or a return statement is encountered.
- In order for the function func to be callable from code outside the source file containing the function, the function must be placed in a file named "func. m " and be the first function defined in that file.
- Immediately following the line containing the function directive, one should provide comments to be printed in response to a help inquiry for the function.
- In MATLAB all input arguments to a function are passed by value.
- For this reason, changes to the input arguments made inside of a function will not propagate to the caller.
- Also, any variables accessed/manipulated inside of a function are local in scope to that function.
demo.m

```
function demo
    t = [0 : . 1 : 8];
    \(\mathrm{f}=\sin (\mathrm{t})\);
    plot(t, f);
end
```

Code to Invoke demo Function (e.g., in Command Window)

```
1 clear t f % Make t and f undefined.
2 demo
3 disp(t); % ERROR: t undefined
4 disp(f); % ERROR: f undefined
```

■ The above is identical to an earlier code example for scripts, except that the code in demo.m has been moved into the body of the function demo.
■ The above example illustrates a key difference between invoking scripts versus invoking functions, namely: a function introduces a new scope (i.e., variables defined/accessed in a function are local to that function).

## Function Example: my_sqr

my_sqr.m

```
function y = my_sqr(x)
    % Compute the square of each element in an array.
    % Use the elementwise product so that the code
    % works in both the scalar and non-scalar cases.
    y = x .^ 2;
end
```

    my_sqr_main.m
    1 a = my_sqr([llll 123\(])\);
    2 disp(a);
    Command Window Output

```
>> my_sqr_main 
```

>>

## Function Example: fahrenheit_to_celsius

fahrenheit_to_celsius.m

```
function \(c=\) fahrenheit_to_celsius(f)
    \% Convert Fahrenheit to Celsius.
    \(c=(5 / 9)\) * (f - 32);
end
```

fahrenheit_to_celsius_main.m

```
1 fprintf('%10s %10s\n', 'Fahrenheit', 'Celsius');
2 for f = 30 : 10 : 100
        c = fahrenheit_to_celsius(f);
    fprintf('%10.2f %10.2f\n', f, c);
end
```

Command Window Output

```
>> fahrenheit_to_celsius_main
Fahrenheit Celsius
    30.00 -1.11
    40.00 4.44
    50.00 10.00
    60.00 15.56
    70.00 21.11
    80.00 26.67
    90.00 32.22
    100.00 37.78
>>
```


## Function Example: roots_of_quadratic

roots_of_quadratic.m

```
function [r1 r2] = roots_of_quadratic (a, b, c)
    \% Find the two roots of the quadratic equation
    \% \(a^{*} \mathrm{x}^{\wedge} 2+\mathrm{b}\) * \(\mathrm{x}+\mathrm{c}=0\).
    if \(a=0\)
        error('leading coefficient is zero');
    end
    \(\mathrm{d}=\operatorname{sqrt}\left(\mathrm{b}{ }^{\wedge} 2-4\right.\) * a * c);
    \(r 1=(-b-d) /(2\) * a);
    \(r 2=(-b+d) /(2\) * \(a)\);
end
```

roots_of_quadratic_main.m
1 [a b] = roots_of_quadratic (1, 0, -16);
2 fprintf('The roots are $\% f$ and $\left.\% f . \backslash n^{\prime}, ~ a, ~ b\right) ;$

## Command Window Output

```
>> roots_of_quadratic_main
The roots are -4.000000 and 4.000000.
>>
```


## Function Example: my_mean

my_mean.m

```
1 function y = my_mean(x)
        % Get the average of the elements in an array.
        y = sum(x(:)) / numel(x);
end
```

my_mean_main.m

```
1 x = [llllllll
2 y = my_mean(x);
3 fprintf('The mean is %f.\n', y);
```

Command Window Output

```
>> my_mean_main
The mean is 4.666667.
>>
```


## Function Example: my_factorial [mulipele funcions Per filie]

## my_factorial.m

```
function \(\mathrm{f}=\) my_factorial(n)
    \% Compute the factorial of each element of an array.
    f = arrayfun(@my_factorial_scalar, n);
end
function \(f=\) my_factorial_scalar(n)
    \% Compute the factorial of a single number.
    if floor ( n ) ~= \(\mathrm{n}|\mid \mathrm{n}<0\)
        error('nonnegative integer required');
    end
    \(\mathrm{f}=1\);
    while \(\mathrm{n}>=2\)
        \(\mathrm{f}=\mathrm{f}\) * n ;
        \(\mathrm{n}=\mathrm{n}-1\);
    end
end
```

my_factorial_main.m

```
1 f = my_factorial(1 : 5);
2 fprintf('%d ', f); fprintf('\n');
```


## Command Window Output

```
>> my_factorial_main
12624120
>>
```


## Variadic Functions

- In MATLAB, functions may take a variable number of input arguments and may return a variable number of output arguments.
- In order for a function to determine the number of input and output arguments and access these arguments, several variables are automatically defined upon entry to a function, including those listed below.

| Name | Description |
| :--- | :--- |
| nargin | number of input arguments |
| nargout | number of output arguments |
| varargin | variable-length input argument |
| varargout | variable-length output argument |

## Variadic Function Example: mysum

## mysum.m

```
function y = mysum(a, b, c)
    % mysum - calculate the sum (of one to three quantities)
    if nargin == 1
            % function called with one argument
            y = a;
    elseif nargin == 2
            % function called with two arguments
            y = a + b;
    elseif nargin == 3
            % function called with three arguments
            y = a + b + c;
    else
            error('invalid number of arguments');
    end
end
```

mysum_main.m

```
1 fprintf('%f %f %f\n', mysum(1), mysum(1, 2), mysum(1, 2, 3));
```


## Command Window Output

```
>> mysum_main
1.000000 3.000000 6.000000
>>
```


## Variadic Function Example: mysum2

## mysum2.m

```
function \(y=\) mysum2(varargin)
    \% mysum2 - Compute the sum of the input arguments.
    if nargin == 0
        y = [];
        return
    end
    y = varargin\{1\};
    for \(i=2\) : nargin
        \(y=y+\) varargin\{i\};
    end
end
```

mysum2_main.m

```
1 fprintf('%f %f\n', mysum2(1, 1, 1, 1, 1, 1, 1, 1), ...
2 mysum(1, 2));
```

Command Window Output

```
>> mysum2_main
8.000000 3.000000
>>
```


## Function Handles and Anonymous Functions

- A function handle provides a means by which to refer to a function.
- A handle to the function $f$ is specified by the syntax $@ f$.
- For example:

```
func = @sin;
x = func(42); % set x to sin(42)
```

- An anonymous function is a function with no name that is identified only by a function handle and is specified using a special syntax.
- This syntax specifies the function parameters and the function body all in a single statement.
- An anonymous function is introduced by an at-character (i.e., "@") followed by the parameter list for the function enclosed in parentheses followed by the single-statement function body.
- If an anonymous function is assigned to a variable, that variable can then be used as a way invoke the anonymous function.
- For example:

```
square = @(x) x .^ 2;
x = square (2); % set x to 2 .^ 2
```

- The built-in MATLAB function sinc computes the normalized sinc function (i.e., $f(t)=\frac{\sin (\pi t)}{\pi t}$ ), not the sinc function (i.e., $\operatorname{sinc}(t)=\frac{\sin (t)}{t}$ ).
- With some care to avoid division by zero, the (unnormalized) sinc function can be implemented as the MATLAB function usinc as follows:

$$
\begin{aligned}
& \text { function } f=\text { usinc( } t) \\
& \qquad f=([t==0]+\sin (t)) . /([t==0]+t) ;
\end{aligned}
$$

end

- The above usinc function could be implemented more compactly using an anonymous function as follows:

$$
\text { usinc }=@(t)([t==0]+\sin (t)) . /([t==0]+t) ;
$$

■ Of course, the above usinc function could also be implemented in terms of MATLAB's built-in sinc function as follows:

$$
\text { usinc }=@(t) \text { sinc }(t / p i)
$$

## Section 15.10

## Data Visualization

## Data Visualization and Plotting

- MATLAB has extensive data-visualization capabilities with support for many types of plots, including:
- line plots in 2-D and 3-D (with support for such things as error bars, stacked plots, and staircase plots)
- scatter plots and bubble charts in 2-D and 3-D
- histograms and pie charts in 2-D and 3-D and heat maps
- discrete plots (such as stem, bar, and staircase plots) in 2-D and 3-D
- geographic plots
- polar plots
- contour plots
- quiver plots in 2-D and 3-D
- surface and mesh plots
- volume visualization (e.g., streams)
- animated plots
- images
- Data visualization output such as plots are generated in figure windows.
- A new figure window can be created with the figure function.
- A figure window is partitioned into an $m \times n$ rectangular grid of subfigures, where a subfigure is simply an individual data visualization, such as a plot.
- When a figure is first created, $m$ and $n$ are effectively defaulted to 1 .
- The values of $m$ and $n$ can be changed with the subplot function.
- The subplot function is also used to control which subfigure should be used when drawing new data visualizations (e.g., plots).
- Normally, plotting data in a subfigure first erases any data already being shown in the subfigure.
- The hold function can be used to override this behavior so that new data can be added to an already existing plot (e.g., to allow plotting multiple lines on a single plot).


## Basic 2-D Plotting Functions

| Name | Description |
| :--- | :--- |
| plot | linear $x-y$ plot |
| $\log \log$ | log log $x-y$ plot |
| semilogx | semi-log $x-y$ plot (x-axis logarithmic) |
| semilogy | semi-log $x$-y plot ( $y$-axis logarithmic) |
| polar | polar plot |
| bar | bar chart |
| stem | stem plot |
| pcolor | pseudocolor (checkerboard) plot |
| fplot | plot function associated with function handle |

Other Plotting-Related Functions

| Name | Description |
| :--- | :--- |
| axis | control axis scaling and appearance |
| hold | hold current plot |
| subplot | multiple axes in single figure |
| figure | create figure |

Plot Annotation Functions

| Name | Description |
| :--- | :--- |
| title | graph title |
| xlabel | x-axis label |
| ylabel | $y$-axis label |
| grid | grid lines |
| text | arbitrarily-positioned text |
| gtext | mouse-positioned text |

Figure Saving/Printing Functions

| Name | Description |
| :--- | :--- |
| print | print figure to file or device |
| saveas | save figure to file |

3-D Plotting Functions

| Name | Description |
| :--- | :--- |
| plot3 | plot lines/points in 3-D space |
| surf | plot surface in color |

Marker Styles

|  |  | Line Colors |  | Symbol | Marker Style |
| :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | . | point |
|  |  | Symbol | Line Color | $\bigcirc$ | circle |
| Line Styles |  |  |  | b | blue | x | cross |
| Symbol | Line Style | $g$ | green | + | plus sign |
| - | solid | r | red | * | asterisk |
| : | dotted | c | cyan | s | square |
| -. | dash dot | m | magenta | d | diamond |
| -- dashed |  | y | yellow | v | triangle (down) |
|  |  | k | black | ^ | triangle (up) |
|  |  | w | white | $<$ | triangle (left) |
|  |  | W |  | > | triangle (right) |
|  |  |  |  | p | pentagram |
|  |  |  |  | h | hexagram |

## Special Symbols for Plot Annotations

| String | Symbol |
| :---: | :---: |
| \alpha | $\alpha$ |
| \beta | $\beta$ |
| \delta | $\delta$ |
| \gamma | $\gamma$ |
| \omega | $\omega$ |
| \theta | $\theta$ |
| \Delta | $\Delta$ |
| \Omega | $\Omega$ |

## Example: Simple Line Plot [figure, plot]

```
1 figure
2 t = linspace(-4 * pi, 4 * pi, 500);
3 y = sin(t);
4 plot(t, y);
5 print('output.pdf', '-dpdf', '-bestfit');
```



## Example: Annotated Line Plot [title, xlabel, ylabel, grid]

```
1 figure
2 w = linspace (-10 * pi, 10 * pi, 500);
3 f = abs(w) .^ 2 .* sin(w);
4 plot(w, f);
5 title('Plot of f(\omega) Versus \omega');
6 xlabel('\omega');
7 ylabel('f(\omega)');
8 grid on
9 print('output.pdf', '-dpdf', '-bestfit', '-opengl');
```



## Example: Multiple Lines Per Plot [hold, legend]

1 figure
2 t = linspace (-2 * pi, 2 * pi, 500);
3 hold on
4 plot(t, sin(t), 'r--');
5 plot(t, $\left.\cos (t), \quad b^{\prime}\right)$;
6 hold off
7 xlabel('t');
8 legend('sin(t)', 'cos(t)');
9 print('output.pdf', '-dpdf', '-bestfit');


## Example: Multiple Plots Per Figure [subplot]

```
figure
2 t = linspace(-pi, pi, 500);
3 subplot(2, 1, 1);
4 plot(t, cos(t));
5 title('The Cosine Function')
6 subplot(2, 1, 2);
7 plot(t, sin(t));
8 title('The Sine Function');
9 print('output.pdf', '-dpdf', '-bestfit');
```



## Example: Multiple Plots Per Figure [Different subplot Grids]

```
figure
t = [-40 40];
subplot(2, 2, 1);
fplot(@(t) sin(1 / 8 * pi * t), t);
title('The Function f'); xlabel('t'); ylabel('f(t)');
subplot(2, 2, 2);
fplot(@(t) cos(1 / 16 * pi * t), t);
title('The Function g'); xlabel('t'); ylabel('g(t)');
subplot(2, 1, 2);
fplot(@(t) sinc(t / pi), t);
title('The Function h'); xlabel('t'); ylabel('h(t)');
print('output.pdf', '-dpdf', '-bestfit');
```



The Function h


## Example: Plot of Sequences [stem]

```
1 figure
2 n = 0 : 50;
3 hold on
4 stem(n, cos(3/32 * pi * n), 'filled', 'r', 'markersize', 5);
5 stem(n, 0.5*sin(1/16 * pi * n), 'filled', 'b', 'markersize', 5)
6 hold off
7 legend('cos(3{\pi}n/32)','sin({\pi}n/16)/2');
8 print('output.pdf', '-dpdf', '-bestfit', '-opengl');
```



## Example: Plot of Points/Path in the Plane [plot, axis]

```
figure
z=[1+j 0.5 1-j -0.5*j -1-j -0.5 -1+j 0.5*j 1+j];
hold on
plot(real(z), imag(z), 'b-', 'linewidth', 2);
plot(real(z), imag(z), 'r.', 'markersize', 60);
hold off
axis equal
grid on
print('output.pdf', '-dpdf', '-bestfit');
```



## Example: Plotting Directly From Function [fplot]

```
figure
2 hold on
3 t = [-8 8] * pi;
4 fplot(@(t) sinc(t / pi), t, 'b-', 'linewidth', 2);
5 fplot(@(t) sin(t) / 8, t, 'r:',' 'linewidth', 2);
6 hold off
7 grid on
8 xlabel('t');
9 legend('sinc(t)', 'sin(t)/8');
10 print('output.pdf', '-dpdf', '-bestfit');
```



- The print function is provided for saving or printing copies of figures.
- MATLAB supports the generation of figures in many formats, including:
- Portable Document Format (PDF)
- PostScript (PS)
- Encapsulated PostScript (EPS)
- ISO/IEC 10918 (JPEG)
- Portable Network Graphics (PNG)
- Scalable Vector Graphics (SVG)
- When saving figures that contain lines and/or text, vector formats (such as PDF) should be preferred over raster formats (such as JPEG and PNG), as vector formats allow distortion-free scaling.
- In many of the earlier plotting examples, the figures generated were saved to a file in PDF format using the print function.
- The print function can also be used to send output directly to a printer.
- For more information on the print function, use the help browser (i.e., "doc print").


## Additional Plotting Examples

- Several additional examples demonstrating various plotting capabilities of MATLAB can be found in the Examples section, including:
- A plot of a univariate function.
- A ploṭ of a bivariate function as a shaded surface.
- A ploṭ of a bivariate function as a surface with contour lines.
- A stem plopt of a sequence.
- A histogram.
- A pie chart.


## Section 15.11

## Symbolic Computation

- Symbolic computation is sometimes quite helpful in solving engineering problems.
- For example, a very complicated formula or equation involving several variables might need to be simplified without assuming specific values for the variables in the formula/equation.
- The Symbolic Math Toolbox provides MATLAB with such symbolic computation capabilities.
- Some computations that can be performed symbolically include:
- simplification of formulas
- integration
- differentiation
- solving systems of equations
- transforms (e.g., forward/inverse unilateral Laplace transform)
- plotting functions expressed in symbolic form
- The Symbolic Math Toolbox defines a new data type called a symbolic object.
- The toolbox uses symbolic objects to represent symbolic variables, constants, and expressions.
- A symbolic object can have as its value any valid mathematical expression.
- Symbolic objects can be used in many of the same ways as non-symbolic objects.
- One must, however, keep in mind that performing a computation symbolically is quite different than performing it non-symbolically.
- Generally speaking, symbolic computation is much slower than non-symbolic computation.


## Creating Symbolic Objects

- A symbolic variable is created using either the sym function or syms directive.
- For example, a symbolic variable $\times$ (whose value is simply itself) can be created using the sym function with the code:

$$
x=\operatorname{sym}\left({ }^{\prime} x^{\prime}\right) ;
$$

- The same result can be achieved in a less verbose manner using the syms directive as follows:
syms x;
- The syms directive allows multiple variables to be specified.
- For example, symbolic variables named $\mathrm{x}, \mathrm{y}$, and z can be created with the code:

```
syms x y z;
```

- The sym function can also be used to create symbolic constants.
- For example, we can create a symbolic constant $p$ that has the value $\pi$ with the code:
p = sym(pi);


## Creating Symbolic Objects (Continued)

■ The character string argument passed to the sym function can only contain a variable name or constant (i.e., it cannot be an arbitrary expression).

- For example, code like the following is not allowed:

```
f = sym('a * x ^ 2 + b * x + c');
    % ERROR: cannot pass arbitrary expression to
    % sym function
```

- From symbolic variables and constants, more complicated symbolic expressions can be constructed.
- For example, a symbolic expression $f$ with the value
a * $\mathrm{x}^{\wedge} 2+\mathrm{b}$ * $\mathrm{x}+\mathrm{c}$ can be created with the code:

```
    syms a b c x;
```

    \(f=a * x^{\wedge} 2+b * x+c ;\)
    - Symbolic objects can often be used in similar ways as non-symbolic objects.
- For example, we can do things like:

```
syms t;
f = t + 1;
g0 = f ^ 3 - 2 * f - 21;
g1 = cos(f) * sin(f / 3);
```

■ Although symbolic and non-symbolic objects can often be used in similar ways, they are very different things.

- For example, the following two lines of code have very different effects:

$$
\begin{aligned}
& x=p i ; \\
& x=\operatorname{sym}(p i) ;
\end{aligned}
$$

- To substitute some expression/variable for another variable, use the subs function.
■ For example, to substitute $t+1$ for $t$ in the expression $t \wedge 2$, we can use the following code:

```
syms t;
f = t ^ 2;
g = subs(f, t, t + 1)
```

- After executing the preceding code, $g$ is associated with the expression:

$$
(t+1) \wedge 2
$$

## Factoring of Symbolic Objects

- To factor a symbolic expression, use the factor function.
- For example, suppose that we want to factor the polynomial $t^{2}+3 t+2$.
- This could be accomplished with the following code:

```
syms t;
f = t^2 ^ 2 3*t + 2;
g= factor(f)
```

- After executing the preceding code, $g$ is associated with the (row-vector) expression:

$$
[t+2, t+1]
$$

- Note that, by default, the factor function will only produce factors with real roots.
- To obtain factors with complex roots, the 'FactorMode' option must be used; for example:

```
syms z;
g = factor(z ^ 2 + 4, 'FactorMode', 'complex')
```

- To simplify a symbolic expression, use the simplify function.

■ For example, suppose that we want to substitute 2 * $t+1$ for $t$ in the expression $t \wedge 2-1$ and simplify the result.

- This can be accomplished with the following code:

```
syms t;
f = t ^ 2 - 1;
g = simplify(subs(f, t, 2 * t + 1))
```

- After executing the preceding code, $g$ is associated with the expression (2*t+1) ^2-1


## Expansion of Symbolic Objects

- To expand an expression, use the expand function.
- For example, to compute $(t+1)^{5}$, we can use the following code:

```
syms t;
f=(t+1)^ ^ ;
g = expand(f)
```

- After executing the preceding code, $g$ is associated with the expression:

$$
t \wedge 5+5 * t \wedge 4+10 * t \wedge 3+10 * t \wedge 2+5 * t+1
$$

- To display an expression in a human-friendly format, use the pretty function.
- For example, to compute $\left[\frac{1}{2} t^{2}+\frac{1}{3}(t+1)\right]^{-4}$ in an expanded and beautified format, we can use the following code:

$$
\begin{aligned}
& \text { syms } t ; \\
& f=\left((1 / 2) * t^{\wedge} 2+(1 / 3) *(t+1)\right)^{\wedge}(-4) ; \\
& \operatorname{pretty}(\operatorname{expand}(f))
\end{aligned}
$$

- The output of the pretty function in this case might look something like:


■ To plot a symbolic expression in one variable, the ezplot function can be used.

- For example, to plot the function $f(t)=3 t^{2}-4 t+2$, we can use the code:

```
syms t;
ezplot(3 * t ^ 2 - 4 * t + 2);
```

- The range of the independent variable may optionally be specified.
- For example, to plot the function $f(t)=3 t^{2}-4 t+2$ over the interval $[-1,1]$, we can use the code:

```
syms t;
ezplot(3 * t ^ 2 - 4 * t + 2, [-1 1]);
```


## Simplification Example

sum_of_arithmetic_sequence.m

```
1 % Compute the formula for the sum of the arithmetic
2 % sequence:
3 % a, a + d, a + 2d, ... a + (n - 1)d.
4 clear
5 syms a d k n
6 f = simplify(symsum(a + k * d, k, 0, n - 1));
7 pretty(f);
```


## Command Window Output

```
>> sum_of_arithmetic_sequence
    \(d n(n-1)\)
a \(n+---------\)
>>
```


## Integration Example

integration_1.m

```
1 % Compute the integral, over the interval [a, b], of the
2 % function:
3% x(t) = t sin(t) e^(-t)
4 clear
5 syms t a b
6 f = int(t * sin(t) * exp(-t), a, b);
7 pretty(f);
```


## Command Window Output

```
>> integration_1
exp(-a) (cos(a) + a cos(a) + a sin(a))
    2
    exp(-b) (cos (b) + b cos(b) + b sin(b))
    2
    >>
```


## Differentiation Example

## differentiation_1.m

```
1 % Compute the derivative of the function:
2 % x (t) = t cos(t) e^(-t)
3 clear
4 syms t
5 f = diff(t * cos(t) * exp(-t), t);
6 pretty(f);
```


## Command Window Output

```
>> differentiation_1
exp(-t) cos(t) - t exp(-t) cos(t) - t exp(-t) sin(t)
>>
```


## Unilateral Laplace Transform Example

## ult_1.m

```
1 % Compute the (forward) unilateral Laplace transform of
2 % the function:
3% x(t) = t e^(-t) sin(t)
4 clear
5 syms t
6 f = laplace(t * exp(-t) * sin(t));
7 pretty(f);
```


## Command Window Output

```
>> ult_1
    2 s + 2
((s+1)}\mp@subsup{)}{}{2}+1\mp@subsup{)}{}{2
>>
```


## Inverse Unilateral Laplace Transform Example

iult_1.m

```
1 % Compute the inverse unilateral Laplace transform of
2 % the function:
3% X(s) = 1/[(s+1) (s + 2)^2]
4 clear
5 syms s
6 f = ilaplace(1 / ((s + 1) * (s + 2)^2));
7 pretty(f);
```


## Command Window Output

```
>> iult_1
exp(-t)- - exp (-2 t) - t exp(-2 t)
>>
```

Plotting Example

```
1 clear
2 syms t
3 ezplot(cos(2 * pi * t) * exp(-abs(t)/2), [-10 10]);
4 print('output.pdf', '-dpdf', '-bestfit');
```



## Section 15.12

## Continuous-Time Signal Processing

- Most CT LTI systems of practical interest are causal with rational transfer functions.
- For this reason, MATLAB has considerable functionality for working with such systems.
- Consider the rational transfer function $H$ with a RHP ROC, where

$$
H(s)=\frac{\sum_{k=1}^{m} a_{k} s^{m-k}}{\sum_{k=1}^{n} b_{k} s^{n-k}}=\frac{a_{1} s^{m-1}+a_{2} s^{m-2}+\ldots+a_{m}}{b_{1} s^{n-1}+b_{2} s^{n-2}+\ldots+b_{n}} .
$$

- Typically, MATLAB represents such a transfer function using two vectors of coefficients, one for the $\left\{a_{k}\right\}$ and one for the $\left\{b_{k}\right\}$.
- For example, $H(s)=\frac{s-1}{s^{2}+3 s+2}$ would be represented by the vectors [1-1] and [1 3 2 $]$.
- Consider a causal LTI (CT) system with a transfer function $H$ of the form $H(s)=\frac{\sum_{k=1}^{m} a_{k} s^{m-k}}{\sum_{k=1}^{n} b_{k} s^{n-k}}=\frac{a_{1} s^{m-1}+a_{2} s^{m-2}+\ldots+a_{m}}{b_{1} s^{n-1}+b_{2} s^{n-2}+\ldots+b_{n}}$.
- The frequency response of such a system can be computed and optionally plotted using the freqs function.
- The freqs function has the interface:

口 [freq_resp w] = freqs (a, b, w)
$\square \mathrm{a}$ is vector whose elements are $a_{1}, a_{2}, \ldots, a_{m}$

- b is vector whose elements are $b_{1}, b_{2}, \ldots, b_{n}$
- w is a vector whose elements are the frequencies at which to calculate the frequency response (and may be omitted in which case w is chosen by the freqs function)
- freq_resp is the vector of frequency response values
- If the return value of the freqs function is discarded, the function also plots magnitude and phase responses, with the magnitude and frequency plotted on logarithmic scales.


## Example: Plotting Frequency Response With freqs

```
1 freqs([1], [1 sqrt(2) 1]);
2 print('output.pdf', '-dpdf', '-bestfit');
```



## Example: Plotting Frequency Response With Linear Scale

```
function [freq_resp, omega] = my_freqs(tf_num, tf_denom, omega)
    % my_freqs - Compute and optionally plot the frequency response.
    % The optional plot uses a linear scale for both axes.
    if nargin >= 3
        [freq_resp, omega] = freqs(tf_num, tf_denom, omega);
    else
            [freq_resp, omega] = freqs(tf_num, tf_denom);
    end
    if nargout == 0
        mag_resp = abs(freq_resp);
        phase_resp = angle(freq_resp) / pi * 180;
        subplot(2, 1, 1);
        plot(omega, mag_resp);
        title('Magnitude Response');
        xlabel('Frequency (rad/s)');
        ylabel('Magnitude (unitless)');
        subplot(2, 1, 2);
        plot(omega, phase_resp);
        title('Phase Response');
        xlabel('Frequency (rad/s)');
        ylabel('Angle (degrees)');
    end
end
```


## Example: Plotting Frequency Response With my_freqs

```
1 my_freqs([1], [1 sqrt(2) 1]);
2 print('output.pdf', '-dpdf', '-bestfit');
```



| Name | Description |
| :--- | :--- |
| besself | design (CT) Bessel filter |
| butter | design (CT) Butterworth filter |
| cheby1 | design (CT/DT) Chebyshev Type 1 filter |
| cheby2 | design (CT/DT) Chebyshev Type 2 filter |
| ellip | design (CT/DT) elliptic or Cauer filter |

## Butterworth Lowpass Filter Design

```
% 16th-order Butterworth (CT) lowpass filter with
% cutoff frequency 500 Hz
[a b] = butter(16, 2 * pi * 500, 'low', 's');
```




## Butterworth Highpass Filter Design

```
% 16th-order Butterworth (CT) highpass filter with
% cutoff frequency 500 Hz
[a b] = butter(16, 2 * pi * 500, 'high', 's');
```




## Butterworth Bandpass Filter Design

```
% 16th-order Butterworth (CT) bandpass filter with
% cutoff frequencies 250 and 750 Hz
[a b] = butter(16, 2 * pi * [250 750], 'bandpass', 's');
```




## Butterworth Bandstop Filter Design

```
% 16th-order Butterworth (CT) bandstop filter with
% cutoff frequencies 250 and 750 Hz
[a b] = butter(16, 2 * pi * [250 750], 'stop', 's');
```





## Bessel Lowpass Filter Design

```
% 16th-order Bessel (CT) lowpass filter with
% cutoff frequency 500 Hz
[a b] = besself(16, 2 * pi * 500);
```




## Chebyshev Type-1 Lowpass Filter Design

| \% 16th-order Chebyshev Type-1 (CT) lowpass filter with \% cutoff frequency 500 Hz and 1 dB passband ripple [a b] = cheby1 (16, 1, 2 * pi * 500, 'low', 's'); |
| :---: |




## Chebyshev Type-2 Lowpass Filter Design

```
% 16th-order Chebyshev Type-2 (CT) lowpass filter with
% cutoff frequency 500 Hz and 20 dB stopband ripple
[a b] = cheby2(16, 20, 2 * pi * 500, 'low', 's');
```




## Elliptic Lowpass Filter Design

```
% 16th-order elliptic (CT) lowpass filter with
% cutoff frequency 500 Hz, 1 dB passband ripple, and
%20 dB stopband ripple
[a b] = ellip(10, 1, 20, 2 * pi * 500, 'low', 's');
```




## Section 15.13

## Discrete-Time Signal Processing

- Most DT LTI systems of practical interest are causal with rational transfer functions.
- For this reason, MATLAB has considerable functionality for working with such systems.
- Consider the rational transfer function

$$
H(z)=\frac{\sum_{k=1}^{m} a_{k} z^{-(k-1)}}{\sum_{k=1}^{n} b_{k} z^{-(k-1)}}=\frac{a_{1}+a_{2} z^{-1}+\ldots+a_{m} z^{-(m-1)}}{b_{1}+b_{2} z^{-1}+\ldots+b_{n} z^{-(n-1)}} .
$$

- Typically, MATLAB represents such a transfer function using two vectors of coefficients, one for the $\left\{a_{k}\right\}$ and one for the $\left\{b_{k}\right\}$.
- For example, $H(z)=\frac{1+z^{-1}}{6+5 z^{-1}+z^{-2}}$ would be represented by the vectors [1 1] and [6 5 1].

■ Consider a causal LTI (DT) system with a transfer function $H$ of the form $H(z)=\frac{\sum_{k=1}^{m} a_{k} z^{-(k-1)}}{\sum_{k=1}^{n} b_{k} z^{-(k-1)}}=\frac{a_{1}+a_{2} z^{-1}+\ldots+a_{m} z^{-(m-1)}}{b_{1}+b_{2} z^{-1}+\ldots+b_{n} z^{-(n-1)}}$.

- The frequency response of such a system can be computed and optionally plotted using the freqz function.
- The freqz function has the interface:
- [freq_resp w] = freqz (a, b, w)
- a is vector whose elements are $a_{1}, a_{2}, \ldots, a_{m}$
- b is vector whose elements are $b_{1}, b_{2}, \ldots, b_{n}$
$\square \mathrm{w}$ is a vector whose elements are the frequencies at which to calculate the frequency response (and may be omitted in which case $w$ is chosen by the freqz function)
- freq_resp is the vector of frequency response values
- If the return value of the freqz function is discarded, the function also plots magnitude and phase responses.


## Example: Plotting Frequency Response With freqz

```
1 freqz([1/6 1/2 1/2 1/6], [1 0 1/3 0]);
2 print('output.pdf', '-dpdf', '-bestfit');
```



| Name | Description |
| :--- | :--- |
| cheby1 | design (CT/DT) Chebyshev Type 1 filter |
| cheby2 | design (CT/DT) Chebyshev Type 2 filter |
| ellip | design (CT/DT) elliptic or Cauer filter |
| fir1 | design (DT) FIR filter using window method |
| fir2 | design (DT) FIR filter using frequency-sampling <br> method |
| cfirpm | complex and nonlinear-phase equiripple FIR filter <br> design |
| fircls | linear-phase FIR filter design using constrained <br> least-squares method |
| firls | linear-phase FIR filter design using least-squares <br> error minimization |
| firpm | Parks-McClellan optimal equiripple FIR filter de- <br> sign |

## Butterworth Lowpass Filter Design

```
% 16th-order Butterworth (DT) lowpass filter with
% normalized cutoff frequency 0.5
[a b] = butter(16, 0.5, 'low');
```




## Butterworth Highpass Filter Design

```
% 16th-order Butterworth (DT) highpass filter with
% normalized cutoff frequency 0.5
[a b] = butter(16, 0.5, 'high');
```




## Butterworth Bandpass Filter Design

```
% 16th-order Butterworth (DT) bandpass filter with
% normalized cutoff frequencies 0.4 and 0.6
[a b] = butter(16, [0.4 0.6], 'bandpass');
```




## Butterworth Bandstop Filter Design

```
% 16th-order Butterworth (DT) bandstop filter with
% normalized cutoff frequencies 0.4 and 0.6
[a b] = butter(16, [0.4 0.6], 'stop');
```




## Chebyshev Type-1 Lowpass Filter Design

```
\% 16th-order Chebyshev Type-1 (DT) lowpass filter with
\% normalized cutoff frequency 0.5 and 1 dB passband ripple
[a b] = chebyl (16, 1, 0.5, 'low');
```




## Chebyshev Type-2 Lowpass Filter Design

```
% 16th-order Chebyshev Type-2 (DT) lowpass filter with
% normalized cutoff frequency 0.5 and 20 dB stopband ripple
[a b] = cheby2(16, 20, 0.5, 'low');
```




## Section 15.14

## References

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4 H. Moore. MATLAB for Engineers. Pearson, Boston, MA, USA, 3rd edition, 2012.
5 A. Gilat. MATLAB - An Introduction with Applications. John Wiley \& Sons, Hoboken, NJ, USA, 6th edition, 2017.
๘ J. Sizemore and J. P. Mueller. MATLAB for Dummies. John Wiley \& Sons, Hoboken, NJ, USA, 2015.

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## Part 16

## Miscellany

## Sum of Arithmetic and Geometric Sequences

- The sum of the arithmetic sequence $a, a+d, a+2 d, \ldots, a+(n-1) d$ is given by

$$
\sum_{k=0}^{n-1}(a+k d)=\frac{n[2 a+d(n-1)]}{2} .
$$

- The sum of the geometric sequence $a, r a, r^{2} a, \ldots, r^{n-1} a$ is given by

$$
\sum_{k=0}^{n-1} r^{k} a=a \frac{r^{n}-1}{r-1} \quad \text { for } r \neq 1
$$

- The sum of the infinite geometric sequence $a, r a, r^{2} a, \ldots$ is given by

$$
\sum_{k=0}^{\infty} r^{k} a=\frac{a}{1-r} \quad \text { for }|r|<1
$$

## Part 17

## References

1 Michael Adams. 2020-05 ECE 260 Video Lectures Playlist on YouTube.
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s Iman Moazzen. Signal Processing Hacks With Iman.
https://www.sphackswithiman.com.
4 Iman Moazzen. YouTube Channel for Signal Processing Hacks With Iman. https://www.youtube. com/channel/UCVkatNMgkEdpWLhH0kBqqLw.
5 Wolfram Alpha Derivative Calculator.
https://www.wolframalpha.com/input/?i=derivative+.
๘ Wolfram Alpha Integral Calculator. https://www.wolframalpha.com/input/?i=integral+.
7 Wolfram Alpha Unilateral Laplace Transform Calculator. https://www. wolframalpha.com/input/?i=laplace+transform+calculator.

## Online Resources II

8 Wolfram Alpha Unilateral Z Transform Calculator. https:
//www.wolframalpha.com/input/?i=Z+transform+calculator.

- DSP Stack Exchange. https://dsp.stackexchange.com.

10 Math Stack Exchange. https://math.stackexchange.com.


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