Signals

→ flow of information

→ measured quantity that varies with time (or position)

→ electrical signal received from a transducer (microphone, thermometer, accelerometer, antenna, etc.)

→ electrical signal that controls a process

Continuous-time signals: voltage, current, temperature, speed, . . .

Discrete-time signals: daily minimum/maximum temperature, lap intervals in races, sampled continuous signals, . . .

Electronics can only deal easily with time-dependent signals, therefore spatial signals, such as images, are typically first converted into a time signal with a scanning process (TV, fax, etc.).
Signal processing

Signals may have to be transformed in order to

→ amplify or filter out embedded information
→ detect patterns
→ prepare the signal to survive a transmission channel
→ prevent interference with other signals sharing a medium
→ undo distortions contributed by a transmission channel
→ compensate for sensor deficiencies
→ find information encoded in a different domain

To do so, we also need

→ methods to measure, characterise, model and simulate transmission channels
→ mathematical tools that split common channels and transformations into easily manipulated building blocks

Analog electronics

Passive networks (resistors, capacities, inductivities, crystals, SAW filters), non-linear elements (diodes, ...), (roughly) linear operational amplifiers

Advantages:

- passive networks are highly linear over a very large dynamic range and large bandwidths
- analog signal-processing circuits require little or no power
- analog circuits cause little additional interference

\[
\frac{U_{\text{in}} - U_{\text{out}}}{R} = \frac{1}{L} \int_{-\infty}^{t} U_{\text{out}} d\tau + C \frac{dU_{\text{out}}}{dt}
\]
Digital signal processing
Analog/digital and digital/analog converter, CPU, DSP, ASIC, FPGA.

Advantages:

→ noise is easy to control after initial quantization
→ highly linear (within limited dynamic range)
→ complex algorithms fit into a single chip
→ flexibility, parameters can easily be varied in software
→ digital processing is insensitive to component tolerances, aging, environmental conditions, electromagnetic interference

But:

→ discrete time processing artifacts (aliasing)
→ can require significantly more power (battery, cooling)
→ digital clock and switching cause interference

Typical DSP applications

→ communication systems
  modulation/demodulation, channel equalization, echo cancellation

→ consumer electronics
  perceptual coding of audio and video on DVDs, speech synthesis, speech recognition

→ music
  synthetic instruments, audio effects, noise reduction

→ medical diagnostics
  magnetic-resonance and ultrasonic imaging, computer tomography, ECG, EEG, MEG, AED, audiology

→ geophysics
  seismology, oil exploration

→ astronomy
  VLBI, speckle interferometry

→ experimental physics
  sensor-data evaluation

→ aviation
  radar, radio navigation

→ security
  steganography, digital watermarking, biometric identification, surveillance systems, signals intelligence, electronic warfare

→ engineering
  control systems, feature extraction for pattern recognition
Syllabus

Signals and systems. Discrete sequences and systems, their types and properties. Linear time-invariant systems, convolution. Harmonic phasors are the eigen functions of linear time-invariant systems. Review of complex arithmetic.

MATLAB. Use of MATLAB on PWF machines to perform numerical experiments and visualise the results in homework exercises.

Fourier transform. Harmonic phasors as orthogonal base functions. Forms of the Fourier transform, convolution theorem, Dirac’s delta function, impulse combs in the time and frequency domain.

Discrete sequences and spectra. Periodic sampling of continuous signals, periodic signals, aliasing, sampling and reconstruction of low-pass and band-pass signals, spectral inversion.

Discrete Fourier transform. Continuous versus discrete Fourier transform, symmetry, linearity, review of the FFT, real-valued FFT.

Spectral estimation. Leakage and scalloping phenomena, windowing, zero padding.

Finite and infinite impulse-response filters. Properties of filters, implementation forms, window-based FIR design, use of frequency-inversion to obtain high-pass filters, use of modulation to obtain band-pass filters, FFT-based convolution, polynomial representation, z-transform, zeros and poles, use of analog IIR design techniques (Butterworth, Chebyshev I/II, elliptic filters).

(If time permits, we may also briefly visit of some modulation techniques.)

Objectives

By the end of the course, you should

→ be able to apply basic properties of time-invariant linear systems

→ understand sampling, aliasing, convolution, filtering, the pitfalls of spectral estimation

→ be able to explain the above in time and frequency domain representations

→ be competent to use filter-design software

→ be able to visualise and discuss digital filters in the z-domain

→ be able to use the FFT for convolution, deconvolution, filtering

→ be able to implement, apply and evaluate simple DSP applications in MATLAB
Textbooks


→ K. Steiglitz: *A digital signal processing primer – with applications to digital audio and computer music*. Addison-Wesley, 1996. (£40)


Sequences and systems

A *discrete sequence* \( \{x_n\} \) is a sequence of numbers

\[
\ldots, x_{-2}, x_{-1}, x_0, x_1, x_2, \ldots
\]

where \( x_n \) denotes the \( n \)-th number in the sequence \( (n \in \mathbb{Z}) \). A discrete sequence maps integer numbers onto real (or complex) numbers. The notation is not well standardized. Some authors write \( x[n] \) instead of \( x_n \), others \( x(n) \).

Where a discrete sequence \( \{x_n\} \) samples a continuous function \( x(t) \) as

\[
x_n = x(t_s \cdot n) = x(n/f_s),
\]

we call \( t_s \) the *sampling period* and \( f_s = 1/t_s \) the *sampling frequency*. A *discrete system* \( T \) receives as input a sequence \( \{x_n\} \) and transforms it into an output sequence \( \{y_n\} = T\{x_n\} \):

\[
\ldots, x_2, x_1, x_0, x_{-1}, \ldots \quad \xrightarrow{\text{discrete system } T} \quad \ldots, y_2, y_1, y_0, y_{-1}, \ldots
\]
Properties of sequences

A sequence \( \{x_n\} \) is

- absolutely summable \( \iff \sum_{n=-\infty}^{\infty} |x_n| < \infty \)
- square summable \( \iff \sum_{n=-\infty}^{\infty} |x_n|^2 < \infty \)
- periodic \( \iff \exists k > 0 : \forall n \in \mathbb{Z} : x_n = x_{n+k} \)

A square-summable sequence is also called an energy signal, and

\[ \sum_{n=-\infty}^{\infty} |x_n|^2 \]

its energy. This terminology reflects that if \( U \) is a voltage supplied to a load resistor \( R \), then \( P = UI = U^2 / R \) is the power consumed.

So even where we drop physical units (e.g., volts) for simplicity in sequence calculations, it is still customary to refer to the squared values of a sequence as power and to its sum or integral over time as energy.

A non-square-summable sequence is a power signal if its average power

\[ \lim_{k \to \infty} \frac{1}{1 + 2k} \sum_{n=-k}^{k} |x_n|^2 \]

exists.

Special sequences

Unit-step sequence:

\[ u_n = \begin{cases} 0, & n < 0 \\ 1, & n \geq 0 \end{cases} \]

Impulse sequence:

\[ \delta_n = \begin{cases} 1, & n = 0 \\ 0, & n \neq 0 \end{cases} = u_n - u_{n-1} \]
Types of discrete systems

A causal system cannot look into the future:

\[ y_n = f(x_n, x_{n-1}, x_{n-2}, \ldots) \]

A memory-less system depends only on the current input value:

\[ y_n = f(x_n) \]

A delay system shifts a sequence in time:

\[ y_n = x_{n-d} \]

\( T \) is a time-invariant system if for any \( d \)

\[ \{y_n\} = T\{x_n\} \iff \{y_{n-d}\} = T\{x_{n-d}\}. \]

\( T \) is a linear system if for any pair of sequences \( \{x_n\} \) and \( \{x'_n\} \)

\[ T\{a \cdot x_n + b \cdot x'_n\} = a \cdot T\{x_n\} + b \cdot T\{x'_n\}. \]

Examples:
The accumulator system

\[ y_n = \sum_{k=-\infty}^{n} x_k \]

is a causal, linear, time-invariant system with memory, as are the backward difference system

\[ y_n = x_n - x_{n-1}, \]

the \( M \)-point moving average system

\[ y_n = \frac{1}{M} \sum_{k=0}^{M-1} x_{n-k} = \frac{x_{n-M+1} + \cdots + x_{n-1} + x_n}{M} \]

and the exponential averaging system

\[ y_n = \alpha \cdot x_n + (1 - \alpha) \cdot y_{n-1} = \alpha \sum_{k=0}^{\infty} (1 - \alpha)^k \cdot x_{n-k}. \]
Examples for time-invariant non-linear memory-less systems are
\[ y_n = x_n^2, \quad y_n = \log_2 x_n, \quad y_n = \max\{\min\{\lfloor 256x_n \rfloor, 255\}, 0\}, \]
examples for linear but not time-invariant systems are
\[ y_n = \begin{cases} x_n, & n \geq 0 \\ 0, & n < 0 \end{cases} = x_n \cdot u_n \]
\[ y_n = x_{\lfloor n/4 \rfloor} \]
\[ y_n = x_n \cdot \Re(e^{\omega jn}) \]
and examples for linear time-invariant non-causal system are
\[ y_n = \frac{1}{2}(x_{n-1} + x_{n+1}) \]
\[ y_n = \sum_{k=-9}^{9} x_{n+k} \cdot \frac{\sin(\pi k\omega)}{\pi k\omega} \cdot [0.5 + 0.5 \cdot \cos(\pi k/10)] \]

**Constant-coefficient difference equations**

Of particular practical interest are causal linear time-invariant systems of the form
\[ y_n = b_0 \cdot x_n - \sum_{k=1}^{N} a_k \cdot y_{n-k} \]

Block diagram representation of sequence operations:

**Addition:** \[ x_n + x'_{n} \]

**Multiplication by constant:** \[ ax_n \]

**Delay:** \[ x_n \rightarrow z^{-1} \rightarrow x_{n-1} \]

The \( a_k \) and \( b_m \) are constant coefficients.
or

\[ y_n = \sum_{m=0}^{M} b_m \cdot x_{n-m} \]

or the combination of both:

\[ \sum_{k=0}^{N} a_k \cdot y_{n-k} = \sum_{m=0}^{M} b_m \cdot x_{n-m} \]

The MATLAB function \textit{filter} is an efficient implementation of the last variant.

### Convolution

All linear time-invariant (LTI) systems can be represented in the form

\[ y_n = \sum_{k=-\infty}^{\infty} a_k \cdot x_{n-k} \]

where \( \{a_k\} \) is a suitably chosen sequence of coefficients.

This operation over sequences is called \textit{convolution} and defined as

\[ \{p_n\} \ast \{q_n\} = \{r_n\} \iff \forall n \in \mathbb{Z} : r_n = \sum_{k=-\infty}^{\infty} p_k \cdot q_{n-k}. \]

If \( \{y_n\} = \{a_n\} \ast \{x_n\} \) is a representation of an LTI system \( T \), with \( \{y_n\} = T\{x_n\} \), then we call the sequence \( \{a_n\} \) the \textit{impulse response} of \( T \), because \( \{a_n\} = T\{\delta_n\} \).
**Convolution examples**

For arbitrary sequences \( \{ p_n \}, \{ q_n \}, \{ r_n \} \) and scalars \( a, b \):

→ Convolution is associative

\[
(\{ p_n \} \ast \{ q_n \}) \ast \{ r_n \} = \{ p_n \} \ast (\{ q_n \} \ast \{ r_n \})
\]

→ Convolution is commutative

\[
\{ p_n \} \ast \{ q_n \} = \{ q_n \} \ast \{ p_n \}
\]

→ Convolution is linear

\[
\{ p_n \} \ast (a \cdot \{ q_n \} + b \cdot \{ r_n \}) = a \cdot (\{ p_n \} \ast \{ q_n \}) + b \cdot (\{ p_n \} \ast \{ r_n \})
\]

→ The impulse sequence (slide 12) is neutral under convolution

\[
\{ p_n \} \ast \{ \delta_n \} = \{ \delta_n \} \ast \{ p_n \} = \{ p_n \}
\]

→ Sequence shifting is equivalent to convolving with a shifted impulse

\[
\{ p_{n-d} \} = \{ p_n \} \ast \{ \delta_{n-d} \}
\]
Can all LTI systems be represented by convolution?

Any sequence \( \{x_n\} \) can be decomposed into a weighted sum of shifted impulse sequences:

\[
\{x_n\} = \sum_{k=-\infty}^{\infty} x_k \cdot \{\delta_{n-k}\}
\]

Let’s see what happens if we apply a linear\(^*\) time-invariant\(^{**}\) system \( T \) to such a decomposed sequence:

\[
T\{x_n\} = T\left( \sum_{k=-\infty}^{\infty} x_k \cdot \{\delta_{n-k}\} \right) = \sum_{k=-\infty}^{\infty} x_k \cdot T\{\delta_{n-k}\}
\]

\[
= \sum_{k=-\infty}^{\infty} x_k \cdot \{\delta_{n-k}\} \ast T\{\delta_n\} = \left( \sum_{k=-\infty}^{\infty} x_k \cdot \{\delta_{n-k}\} \right) \ast T\{\delta_n\} = \{x_n\} \ast T\{\delta_n\} \quad \text{q.e.d.}
\]

\( \Rightarrow \) The impulse response \( T\{\delta_n\} \) fully characterizes an LTI system.

**Exercise 1** What type of discrete system (linear/non-linear, time-invariant/non-time-invariant, causal/non-causal, causal, memory-less, etc.) is:

(a) \( y_n = |x_n| \)

(b) \( y_n = -x_{n-1} + 2x_n - x_{n+1} \)

(c) \( y_n = \prod_{i=0}^{8} x_{n-i} \)

(d) \( y_n = \frac{1}{2}(x_{2n} + x_{2n+1}) \)

(e) \( y_n = \frac{3x_{n-1} + x_{n-2}}{x_{n-3}} \)

(f) \( y_n = x_n \cdot e^{n/14} \)

(g) \( y_n = x_n \cdot u_n \)

(h) \( y_n = \sum_{i=-\infty}^{\infty} x_i \cdot \delta_{i-n+2} \)

**Exercise 2**

Prove that convolution is (a) commutative and (b) associative.
**Exercise 3** A finite-length sequence is non-zero only at a finite number of positions. If \( m \) and \( n \) are the first and last non-zero positions, respectively, then we call \( n - m + 1 \) the length of that sequence. What maximum length can the result of convolving two sequences of length \( k \) and \( l \) have?

**Exercise 4** The length-3 sequence \( a_0 = -3, \ a_1 = 2, \ a_2 = 1 \) is convolved with a second sequence \( \{ b_n \} \) of length 5.

(a) Write down this linear operation as a matrix multiplication involving a matrix \( A \), a vector \( \vec{b} \in \mathbb{R}^5 \), and a result vector \( \vec{c} \).

(b) Use MATLAB to multiply your matrix by the vector \( \vec{b} = (1, 0, 0, 2, 2) \) and compare the result with that of using the conv function.

(c) Use the MATLAB facilities for solving systems of linear equations to undo the above convolution step.

**Exercise 5** (a) Find a pair of sequences \( \{ a_n \} \) and \( \{ b_n \} \), where either contains at least three different values and where the convolution \( \{ a_n \} * \{ b_n \} \) results in an all-zero sequence.

(b) Does every LTI system \( T \) have an inverse LTI system \( T^{-1} \) such that \( \{ x_n \} = T^{-1}T\{ x_n \} \) for all sequences \( \{ x_n \} \)? Why?

---

**Direct form I and II implementations**

The block diagram representation of the constant-coefficient difference equation on slide 17 is called the *direct form I implementation*. The number of delay elements can be halved by using the commutativity of convolution to swap the two feedback loops, leading to the *direct form II implementation* of the same LTI system.
Convolution: optics example

If a projective lens is out of focus, the blurred image is equal to the original image convolved with the aperture shape (e.g., a filled circle):

$$I \ast h(\text{disk}, r = \frac{a/2}{f})$$

Point-spread function $h$:  

$$h(x, y) = \begin{cases} \frac{1}{\pi r^2}, & x^2 + y^2 \leq r^2 \\ 0, & x^2 + y^2 > r^2 \end{cases}$$

Original image $I$, blurred image $B = I \ast h$, i.e.

$$B(x, y) = \int \int I(x-x', y-y') \cdot h(x', y') \cdot dx'dy'$$

Convolution: electronics example

Any passive network ($R, L, C$) convolves its input voltage $U_{in}$ with an impulse response function $h$, leading to $U_{out} = U_{in} \ast h$, that is

$$U_{out}(t) = \int_{-\infty}^{\infty} U_{in}(t-\tau) \cdot h(\tau) \cdot d\tau$$

In this example:

$$\frac{U_{in} - U_{out}}{R} = C \cdot \frac{dU_{out}}{dt}$$

$$\quad h(t) = \begin{cases} \frac{1}{RC} \cdot e^{-t/RC}, & t \geq 0 \\ 0, & t < 0 \end{cases}$$
Why are sine waves useful?

1) Adding together sine waves of equal frequency, but arbitrary amplitude and phase, results in another sine wave of the same frequency:

\[ A_1 \cdot \sin(\omega t + \varphi_1) + A_2 \cdot \sin(\omega t + \varphi_2) = A \cdot \sin(\omega t + \varphi) \]

with

\[ A = \sqrt{A_1^2 + A_2^2 + 2A_1A_2 \cos(\varphi_2 - \varphi_1)} \]

\[ \tan \varphi = \frac{A_1 \sin \varphi_1 + A_2 \sin \varphi_2}{A_1 \cos \varphi_1 + A_2 \cos \varphi_2} \]

Also

\[ A \cdot \sin(\omega t + \varphi) = \]
\[ a \cdot \sin(\omega t) + b \cdot \cos(\omega t) \]

with \( A = \sqrt{a^2 + b^2} \) and \( \tan \varphi = \frac{b}{a} \).

Note: Convolution of a discrete sequence \( \{x_n\} \) with another sequence \( \{y_n\} \) is nothing but adding together scaled and delayed copies of \( \{x_n\} \).
(Think of \( \{y_n\} \) as decomposition into a sum of impulses.)

If \( \{x_n\} \) is a sampled sine wave of frequency \( f \), so is \( \{x_n\} * \{y_n\} \)!

\[ \implies \text{Sine-wave sequences form a family of discrete sequences} \]
\[ \text{that is closed under convolution with arbitrary sequences.} \]

The same applies for continuous sine waves and convolution.

2) Sine waves are orthogonal to each other:

\[ \int_{-\infty}^{\infty} \sin(\omega_1 t + \varphi_1) \cdot \sin(\omega_2 t + \varphi_2) \, dt = 0 \]

\[ \iff \omega_1 \neq \omega_2 \lor \varphi_1 - \varphi_2 = (2k + 1)\pi \quad (k \in \mathbb{Z}) \]

They can be used to form an orthogonal function basis for a transform.
Why are exponential functions useful?

Adding together two exponential functions with the same base \( z \), but different scale factor and offset, results in another exponential function with the same base:

\[
A_1 \cdot z^{t+\varphi_1} + A_2 \cdot z^{t+\varphi_2} = A_1 \cdot z^t \cdot z^{\varphi_1} + A_2 \cdot z^t \cdot z^{\varphi_2} = (A_1 \cdot z^{\varphi_1} + A_2 \cdot z^{\varphi_2}) \cdot z^t = A \cdot z^t
\]

Likewise, if we convolve a sequence \( \{x_n\} \) of values \( \ldots, z^{-3}, z^{-2}, z^{-1}, 1, z, z^2, z^3, \ldots \)

\[x_n = z^n\] with an arbitrary sequence \( \{h_n\} \), we get \( \{y_n\} = \{z^n\} * \{h_n\}, \)

\[y_n = \sum_{k=-\infty}^{\infty} x_{n-k} \cdot h_k = \sum_{k=-\infty}^{\infty} z^{n-k} \cdot h_k = z^n \cdot \sum_{k=-\infty}^{\infty} z^{-k} \cdot h_k = z^n \cdot H(z)\]

where \( H(z) \) is independent of \( n \).

Exponential sequences are closed under convolution with arbitrary sequences. The same applies in the continuous case.

Why are complex numbers so useful?

1) They give us all \( n \) solutions ("roots") of equations involving polynomials up to degree \( n \) (the \( \sqrt{-1} = j \) story).

2) They form the "great unifying theory" that combines sine functions and exponential functions:

\[
\cos(\omega t) = \frac{1}{2} \left( e^{j\omega t} + e^{-j\omega t} \right)
\]

\[
\sin(\omega t) = \frac{1}{2j} \left( e^{j\omega t} - e^{-j\omega t} \right)
\]

or

\[
\cos(\omega t + \varphi) = \frac{1}{2} \left( e^{j\omega t+\varphi} + e^{-j\omega t-\varphi} \right)
\]

or

\[
\cos(\omega n + \varphi) = \Re\{e^{j\omega n+\varphi}\} = \Re\{(e^{j\omega})^n \cdot e^{j\varphi}\}
\]

\[
\sin(\omega n + \varphi) = \Im\{e^{j\omega n+\varphi}\} = \Im\{(e^{j\omega})^n \cdot e^{j\varphi}\}
\]
We can now represent sine waves as projections of a rotating complex vector. This allows us to represent sine-wave sequences as exponential sequences with basis $e^{j\omega}$.

A phase shift in such a sequence corresponds to a rotation of a complex vector.

3) Complex multiplication allows us to modify the amplitude and phase of a complex rotating vector using a single operation and value.

Rotation of a 2D vector in $(x, y)$-form is notationally slightly messy, but fortunately $j^2 = -1$ does exactly what is required here:

$$
\begin{pmatrix}
  x_3 \\
  y_3
\end{pmatrix}
= \begin{pmatrix}
  x_2 & -y_2 \\
  y_2 & x_2
\end{pmatrix}
\cdot
\begin{pmatrix}
  x_1 \\
  y_1
\end{pmatrix}
= \begin{pmatrix}
  x_1x_2 - y_1y_2 \\
  x_1y_2 + x_2y_1
\end{pmatrix}
\cdot
\begin{pmatrix}
  -y_2 \\
  x_2
\end{pmatrix}
$$

$$z_1 = x_1 + jy_1, \quad z_2 = x_2 + jy_2$$

$$z_1 \cdot z_2 = x_1x_2 - y_1y_2 + j(x_1y_2 + x_2y_1)$$

### Complex phasors

Amplitude and phase are two distinct characteristics of a sine function that are inconvenient to keep separate notationally.

Complex functions (and discrete sequences) of the form

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

(where $j^2 = -1$) are able to represent both amplitude and phase in one single algebraic object.

Thanks to complex multiplication, we can also incorporate in one single factor both a multiplicative change of amplitude and an additive change of phase of such a function. This makes discrete sequences of the form

$$x_n = e^{j\omega n}$$

*eigensequences* with respect to an LTI system $T$, because for each $\omega$, there is a complex number (eigenvalue) $H(\omega)$ such that

$$T\{x_n\} = H(\omega) \cdot \{x_n\}$$

In the notation of slide 29, where the argument of $H$ is the base, we would write $H(e^{j\omega})$. 
Recall: Fourier transform

The Fourier integral transform and its inverse are defined as

\[ \mathcal{F}\{g(t)\}(\omega) = G(\omega) = \alpha \int_{-\infty}^{\infty} g(t) \cdot e^{\pm j\omega t} \, dt \]

\[ \mathcal{F}^{-1}\{G(\omega)\}(t) = g(t) = \beta \int_{-\infty}^{\infty} G(\omega) \cdot e^{\pm j\omega t} \, d\omega \]

where \( \alpha \) and \( \beta \) are constants chosen such that \( \alpha \beta = 1/(2\pi) \).

Many equivalent forms of the Fourier transform are used in the literature, and there is no strong consensus on whether the forward transform uses \( e^{-j\omega t} \) and the backwards transform \( e^{j\omega t} \), or vice versa. Some authors set \( \alpha = 1 \) and \( \beta = 1/(2\pi) \), to keep the convolution theorem free of a constant prefactor; others use \( \alpha = \beta = 1/\sqrt{2\pi} \), in the interest of symmetry.

The substitution \( \omega = 2\pi f \) leads to a form without prefactors:

\[ \mathcal{F}\{h(t)\}(f) = H(f) = \int_{-\infty}^{\infty} h(t) \cdot e^{\pm 2\pi jft} \, dt \]

\[ \mathcal{F}^{-1}\{H(f)\}(t) = h(t) = \int_{-\infty}^{\infty} H(f) \cdot e^{\pm 2\pi jft} \, df \]

Another notation is in the continuous case

\[ \mathcal{F}\{h(t)\}(\omega) = H(e^{j\omega}) = \int_{-\infty}^{\infty} h(t) \cdot e^{-j\omega t} \, dt \]

\[ \mathcal{F}^{-1}\{H(e^{j\omega})\}(t) = h(t) = \frac{1}{2\pi} \int_{-\infty}^{\infty} H(e^{j\omega}) \cdot e^{j\omega t} \, d\omega \]

and in the discrete-sequence case

\[ \mathcal{F}\{h_n\}(\omega) = H(e^{j\omega}) = \sum_{n=-\infty}^{\infty} h_n \cdot e^{-j\omega n} \]

\[ \mathcal{F}^{-1}\{H(e^{j\omega})\}(t) = h_n = \frac{1}{2\pi} \int_{-\pi}^{\pi} H(e^{j\omega}) \cdot e^{j\omega n} \, d\omega \]

which treats the Fourier transform as a special case of the \( z \)-transform (to be introduced shortly).
Convolution theorem

Continuous form:

\[ F\{(f * g)(t)\} = F\{f(t)\} \cdot F\{g(t)\} \]

\[ F\{f(t) \cdot g(t)\} = F\{f(t)\} \ast F\{g(t)\} \]

Discrete form:

\( \{x_n\} \ast \{y_n\} = \{z_n\} \iff X(e^{j\omega}) \cdot Y(e^{j\omega}) = Z(e^{j\omega}) \)

Convolution in the time domain is equivalent to (complex) scalar multiplication in the frequency domain.

Convolution in the frequency domain corresponds to scalar multiplication in the time domain.

Proof:
\[ z(r) = \int_s x(s)y(r-s)ds \iff \int_r z(r)e^{-j\omega r}dr = \int_r \int_s x(s)y(r-s)e^{-j\omega r}dsdr = \int_s x(s)\int_r y(r-s)e^{-j\omega r}drds = \int_s x(s)e^{-j\omega s}\int_r y(r-s)e^{-j\omega(r-s)}drds = \int_s x(s)e^{-j\omega s}\int_t y(t)e^{-j\omega t}dt ds \]

Dirac’s delta function

The continuous equivalent of the impulse sequence \( \{\delta_n\} \) is known as Dirac’s delta function \( \delta(x) \). It is a generalized function, defined such that

\[ \delta(x) = \begin{cases} 0, & x \neq 0 \\ \infty, & x = 0 \end{cases} \]

\[ \int_{-\infty}^{\infty} \delta(x) dx = 1 \]

and can be thought of as the limit of function sequences such as

\[ \delta(x) = \lim_{n \to \infty} \begin{cases} 0, & |x| \geq 1/n \\ n/2, & |x| < 1/n \end{cases} \]

or

\[ \delta(x) = \lim_{n \to \infty} \frac{n}{\sqrt{\pi}} e^{-n^2x^2} \]

The delta function is mathematically speaking not a function, but a distribution, that is an expression that is only defined when integrated.
Some properties of Dirac’s delta function:

\[
\int_{-\infty}^{\infty} f(x) \delta(x - a) \, dx = f(a)
\]

\[
\int_{-\infty}^{\infty} e^{\pm 2\pi j ft} \, df = \delta(t)
\]

\[
\frac{1}{2\pi} \int_{-\infty}^{\infty} e^{\pm j \omega t} \, d\omega = \delta(t)
\]

Fourier transform:

\[
\mathcal{F}\{\delta(t)\}(\omega) = \int_{-\infty}^{\infty} \delta(t) \cdot e^{-j\omega t} \, dt = e^{0} = 1
\]

\[
\mathcal{F}^{-1}\{1\}(t) = \frac{1}{2\pi} \int_{-\infty}^{\infty} 1 \cdot e^{j\omega t} \, d\omega = \delta(t)
\]


**Sine and cosine in the frequency domain**

\[
\cos(2\pi ft) = \frac{1}{2} e^{2\pi j ft} + \frac{1}{2} e^{-2\pi j ft} \quad \sin(2\pi ft) = \frac{1}{2j} e^{2\pi j ft} - \frac{1}{2j} e^{-2\pi j ft}
\]

As any real-valued signal \( x(t) \) can be represented as a combination of sine and cosine functions, the spectrum of any real-valued signal will show the symmetry \( X(e^{j\omega}) = [X(e^{-j\omega})]^* \), where * denotes the complex conjugate (i.e., negated imaginary part).
Sampling using a Dirac comb

The loss of information in the sampling process that converts a continuous function \( x(t) \) into a discrete sequence \( \{x_n\} \) defined by

\[
x_n = x(t_s \cdot n) = x(n/f_s)
\]

can be modelled through multiplying \( x(t) \) by a comb of Dirac impulses

\[
s(t) = \sum_{n=-\infty}^{\infty} \delta(t - t_s \cdot n)
\]

to obtain the sampled function

\[
\hat{x}(t) = x(t) \cdot s(t)
\]

The function \( \hat{x}(t) \) now contains exactly the same information as the discrete sequence \( \{x_n\} \), but is still in a form that can be analysed using the Fourier transform on continuous functions.

The Fourier transform of a Dirac comb

\[
s(t) = \sum_{n=-\infty}^{\infty} \delta(t - t_s \cdot n)
\]

is another Dirac comb

\[
S(f) = \mathcal{F}\left\{ \sum_{n=-\infty}^{\infty} \delta(t - t_s n) \right\}(f) = \int_{-\infty}^{\infty} \sum_{n=-\infty}^{\infty} \delta(t - t_s n) e^{2\pi jft} dt = \frac{1}{t_s} \sum_{n=-\infty}^{\infty} \delta \left( f - \frac{n}{t_s} \right).
\]
A wave \( \cos(2\pi tf) \) sampled at frequency \( f_s \) cannot be distinguished from \( \cos(2\pi t(kf_s \pm f)) \) for any \( k \in \mathbb{Z} \). Therefore, ensure \( |f| < f_s/2 \).
Reconstruction of a continuous band-limited waveform

The ideal anti-aliasing filter for eliminating any frequency content above $f_s/2$ before sampling with a frequency of $f_s$ has the Fourier transform

$$H(f) = \begin{cases} 
1 & \text{if } -\frac{f_s}{2} < f \leq \frac{f_s}{2} \\
0 & \text{otherwise}
\end{cases}.$$ 

This leads, after an inverse Fourier transform, to the impulse response

$$h(t) = \frac{\sin \pi t f_s}{\pi t f_s}.$$ 

The original band-limited signal can be reconstructed by convolving this with the sampled signal $\hat{x}(t)$, which eliminates the periodicity of the frequency domain introduced by the sampling process:

$$x(t) = h(t) \ast \hat{x}(t)$$

Note that sampling $h(t)$ gives the impulse function: $h(t) \cdot s(t) = \delta(t)$.

Impulse response of rectangular filter with cut–off frequency $f_s/2$
Bandpass sampling

Sampled signals can be reconstructed as long as their doublesided bandwidth (counting both the positive and negative frequency range) does not exceed the sampling frequency. The anti-aliasing filter used can also be a band-pass filter instead of a low-pass filter.

A suitably bandlimited signal $x(t)$ with center frequency $f_c$ sampled at frequency $f_s$ can be reconstructed by convolution with

$$h(t) = \frac{\sin \pi t f_s}{\pi t f_s} \cdot \cos(2\pi t f_c).$$

Exercise 6

- Generate a one second long Gaussian noise sequence $\{r_n\}$ (using MATLAB function `randn`) with a sampling rate of 300 Hz.

- Use the `fir1(50, 45/150)` function to design a finite impulse response low-pass filter with a cut-off frequency of 45 Hz. Use the `filtfilt` function in order to apply that filter to the generated noise signal, resulting in the filtered noise signal $\{x_n\}$.

- Then sample $\{x_n\}$ at 100 Hz by setting all but every third sample value to zero, resulting in sequence $\{y_n\}$.

- Generate another lowpass filter with a cut-off frequency of 50 Hz and apply it to $\{y_n\}$, in order to interpolate the reconstructed filtered noise signal $\{z_n\}$. Multiply the result by three, to compensate the energy lost during sampling.

- Plot $\{x_n\}$, $\{y_n\}$, and $\{z_n\}$. Finally compare $\{x_n\}$ and $\{z_n\}$.

Why should the first filter have a lower cut-off frequency than the second?
**Spectrum of a periodic signal**

A signal $x(t)$ that is periodic with frequency $f_p$ can be factored into a single period $\hat{x}(t)$ convolved with an impulse comb $p(t)$. This corresponds in the frequency domain to the multiplication of the spectrum of the single period with a comb of impulses spaced $f_p$ apart.

$$x(t) \leftrightarrow X(f)$$

$$\hat{x}(t) \leftrightarrow \hat{X}(f)$$

$$p(t) \leftrightarrow P(f)$$

**Spectrum of a sampled signal**

A signal $x(t)$ that is sampled with frequency $f_s$ has a spectrum that is periodic with a period of $f_s$.

$$x(t) \leftrightarrow X(f)$$

$$\hat{x}(t) \leftrightarrow \hat{X}(f)$$

$$S(f) \leftrightarrow \hat{X}(f)$$
Continuous vs discrete Fourier transform

- Sampling a **continuous** signal makes its spectrum **periodic**
- A **periodic** signal has a **sampled** spectrum

We **sample** a signal \( x(t) \) with \( f_s \), getting \( \hat{x}(t) \). We take \( n \) consecutive samples of \( \hat{x}(t) \) and **repeat** these periodically, getting a new signal \( \ddot{x}(t) \) with period \( n/f_s \). Its spectrum \( \dddot{X}(f) \) is **sampled** (i.e., has non-zero value) at frequency intervals \( f_s/n \) and **repeats** itself with a period \( f_s \).

Now both \( \ddot{x}(t) \) and its spectrum \( \dddot{X}(f) \) are **finite** vectors of length \( n \).

Properties of the Fourier transform

If
\[
x(t) \quad \longleftrightarrow\quad X(f)
\]
and
\[
y(t) \quad \longleftrightarrow\quad Y(f)
\]
are pairs of functions that are mapped onto each other by the Fourier transform, then so are the following pairs.

**Linearity:**
\[
a x(t) + b y(t) \quad \longleftrightarrow\quad a X(f) + b Y(f)
\]

**Time scaling:**
\[
x(at) \quad \longleftrightarrow\quad \frac{1}{|a|} X \left( \frac{f}{a} \right)
\]

**Frequency scaling:**
\[
\frac{1}{|a|} x \left( \frac{t}{a} \right) \quad \longleftrightarrow\quad X(af)
\]
Time shifting:

\[ x(t - \Delta t) \quad \bullet \circ \quad X(f) \cdot e^{-2\pi j f \Delta t} \]

Frequency shifting:

\[ x(t) \cdot e^{2\pi j f \Delta t} \quad \bullet \circ \quad X(f - \Delta f) \]

Parseval’s theorem (total power):

\[ \int_{-\infty}^{\infty} |x(t)|^2 dt = \int_{-\infty}^{\infty} |X(f)|^2 df \]

---

**Fourier transform symmetries**

We call a function \( x(t) \)

odd: \( x(-t) = -x(t) \)
even: \( x(-t) = x(t) \)

and \( \cdot^* \) is the complex conjugate, such that \( (a + jb)^* = (a - jb) \).

Then

- \( x(t) \) is real \( \iff \) \( X(-f) = [X(f)]^* \)
- \( x(t) \) is imaginary \( \iff \) \( X(-f) = -[X(f)]^* \)
- \( x(t) \) is even \( \iff \) \( X(f) \) is even
- \( x(t) \) is odd \( \iff \) \( X(f) \) is odd
- \( x(t) \) is real and even \( \iff \) \( X(f) \) is real and even
- \( x(t) \) is real and odd \( \iff \) \( X(f) \) is imaginary and odd
- \( x(t) \) is imaginary and even \( \iff \) \( X(f) \) is imaginary and even
- \( x(t) \) is imaginary and odd \( \iff \) \( X(f) \) is real and odd

---

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Discrete Fourier Transform (DFT)

\[ X_k = \sum_{i=0}^{n-1} x_i \cdot e^{-2\pi j k i / n} \quad x_k = \frac{1}{n} \sum_{i=0}^{n-1} X_i \cdot e^{2\pi j k i / n} \]

The \( n \)-point DFT multiplies a vector with an \( n \times n \) matrix

\[
F_n = \begin{pmatrix}
1 & 1 & 1 & \cdots & 1 \\
1 & e^{-2\pi j \frac{1}{n}} & e^{-2\pi j \frac{2}{n}} & \cdots & e^{-2\pi j \frac{n-1}{n}} \\
1 & e^{-2\pi j \frac{2}{n}} & e^{-2\pi j \frac{4}{n}} & \cdots & e^{-2\pi j \frac{2(n-1)}{n}} \\
\vdots & \vdots & \vdots & \ddots & \vdots \\
1 & e^{-2\pi j \frac{n-1}{n}} & e^{-2\pi j \frac{2(n-1)}{n}} & \cdots & e^{-2\pi j \frac{(n-1)(n-1)}{n}} \\
\end{pmatrix}
\]

\[
F_n \cdot \begin{pmatrix} x_0 \\ x_1 \\ x_2 \\ \vdots \\ x_{n-1} \end{pmatrix} = \begin{pmatrix} X_0 \\ X_1 \\ X_2 \\ \vdots \\ X_{n-1} \end{pmatrix}, \quad \frac{1}{n} \cdot F_n^* \cdot \begin{pmatrix} X_0 \\ X_1 \\ X_2 \\ \vdots \\ X_{n-1} \end{pmatrix} = \begin{pmatrix} x_0 \\ x_1 \\ x_2 \\ \vdots \\ x_{n-1} \end{pmatrix}
\]

Discrete Fourier Transform visualized

The \( n \)-point DFT of a signal \( \{x_i\} \) sampled at frequency \( f_s \) contains in the elements \( X_0 \) to \( X_{n/2} \) of the resulting frequency-domain vector the frequency components 0, \( f_s/n, 2f_s/n, 3f_s/n, \ldots, f_s/2 \), and contains in \( X_{n-1} \) downto \( X_{n/2} \) the corresponding negative frequencies. Note that for a real-valued input vector, both \( X_0 \) and \( X_{n/2} \) will be real, too.

Why is there no phase information recovered at \( f_s/2 \)?
Fast Fourier Transform (FFT)

\[
\mathcal{F}_n\{x_i\}_{i=0}^{n-1} = \sum_{i=0}^{n-1} x_i \cdot e^{2\pi j \frac{ik}{n}}
\]

\[
= \sum_{i=0}^{\frac{n}{2}-1} x_{2i} \cdot e^{2\pi j \frac{ik}{n/2}} + e^{2\pi j \frac{k}{n}} \sum_{i=0}^{\frac{n}{2}-1} x_{2i+1} \cdot e^{2\pi j \frac{ik}{n/2}}
\]

\[
= \begin{cases} 
\mathcal{F}_{\frac{n}{2}}\{x_{2i}\}_{i=0}^{\frac{n}{2}-1} & + e^{2\pi j \frac{k}{n}} \cdot \left( \mathcal{F}_{\frac{n}{2}}\{x_{2i+1}\}_{i=0}^{\frac{n}{2}-1} \right)_k, & k < \frac{n}{2} \\
\mathcal{F}_{\frac{n}{2}}\{x_{2i}\}_{i=0}^{\frac{n}{2}-1} & + e^{2\pi j \frac{k}{n}} \cdot \left( \mathcal{F}_{\frac{n}{2}}\{x_{2i+1}\}_{i=0}^{\frac{n}{2}-1} \right)_{k-\frac{n}{2}}, & k \geq \frac{n}{2} 
\end{cases}
\]

The DFT over \(n\)-element vectors can be reduced to two DFTs over \(n/2\)-element vectors plus \(n\) multiplications and \(n\) additions, leading to \(\log_2 n\) rounds and \(n \log_2 n\) additions and multiplications overall, compared to \(n^2\) for the equivalent matrix multiplication.

A high-performance FFT implementation in C with many processor-specific optimizations and support for non-power-of-2 sizes is available at [http://www.fftw.org/](http://www.fftw.org/).

Efficient real-valued FFT

The symmetry properties of the Fourier transform applied to the discrete Fourier transform \(\{X_i\}_{i=0}^{n-1} = \mathcal{F}_n\{x_i\}_{i=0}^{n-1}\) have the form

\[\forall i : x_i = \Re(x_i) \iff \forall i : X_{n-i} = X_i^*\]

\[\forall i : x_i = j \cdot \Im(x_i) \iff \forall i : X_{n-i} = -X_i^*\]

These two symmetries, combined with the linearity of the DFT, allows us to calculate two real-valued \(n\)-point DFTs

\(\{X'_i\}_{i=0}^{n-1} = \mathcal{F}_n\{x'_i\}_{i=0}^{n-1}\)

\(\{X''_i\}_{i=0}^{n-1} = \mathcal{F}_n\{x''_i\}_{i=0}^{n-1}\)

simultaneously in a single complex-valued \(n\)-point DFT, by composing its input as

\[x_i = x'_i + j \cdot x''_i\]

and decomposing its output as

\[X'_i = \frac{1}{2}(X_i + X_{n-i}^*)\]

\[X''_i = \frac{1}{2}(X_i - X_{n-i}^*)\]

To optimize the calculation of a single real-valued FFT, use this trick to calculate the two half-size real-value FFTs that occur in the first round.
Fast complex multiplication

Calculating the product of two complex numbers as

\[(a + jb) \cdot (c + jd) = (ac - bd) + j(ad + bc)\]

involves four (real-valued) multiplications and two additions.

The alternative calculation

\[(a + jb) \cdot (c + jd) = (\alpha - \beta) + j(\alpha + \gamma) \text{ with }\]

\[
\alpha = a(c + d) \\
\beta = d(a + b) \\
\gamma = c(b - a)
\]

provides the same result with three multiplications and five additions.

The latter may perform faster on CPUs where multiplications take three or more times longer than additions.

This trick is most helpful on simpler microcontrollers. Specialized signal-processing CPUs (DSPs) feature 1-clock-cycle multipliers. High-end desktop processors use pipelined multipliers that stall where operations depend on each other.

FFT-based convolution

Calculating the convolution of two finite sequences \[\{x_i\}_{i=0}^{m-1}\] and \[\{y_i\}_{i=0}^{n-1}\] of lengths \(m\) and \(n\) via

\[z_i = \min\{m-1,i\} \sum_{j=\max\{0,i-(n-1)\}} x_j \cdot y_{i-j}, \quad 0 \leq i < m + n - 1\]

takes \(mn\) multiplications.

Can we apply the FFT and the convolution theorem to calculate the convolution faster, in just \(O(m \log m + n \log n)\) multiplications?

\[\{z_i\} = \mathcal{F}^{-1}(\mathcal{F}\{x_i\} \cdot \mathcal{F}\{y_i\})\]

There is obviously no problem if this condition is fulfilled:

\[\{x_i\} \text{ and } \{y_i\} \text{ are periodic, with equal period lengths}\]

In this case, the fact that the DFT interprets its input as a single period of a periodic signal will do exactly what is needed, and the FFT and inverse FFT can be applied directly as above.
In the general case, measures have to be taken to prevent a wrap-over:

\[ F^{-1}[F(A) \cdot F(B)] \]

Both sequences are padded with zero values to a length of at least \( m + n - 1 \). This ensures that the start and end of the resulting sequence do not overlap.

Zero padding is typically applied to extend the length of both sequences to the next larger power of two \( 2^{\lceil \log_2 (m+n-1) \rceil} \), to facilitate use of the FFT.

With a causal sequence, simply append the padding zeros at the end. With a non-causal sequence, values with a negative index number are wrapped around the DFT block boundaries and appear at the right end. In this case, zero-padding is applied in the center of the block, between the last and first element of the sequence.

Thanks to the periodic nature of the DFT, zero padding at both ends has the same effect as padding only at one end.

If both sequences can be loaded entirely into RAM, the FFT can be applied to them in one step. However, one of the sequences might be too large for that. It could also be a realtime waveform (e.g., a telephone signal) that must not be delayed until the end of the transmission. In such cases, the sequence has to be split into shorter blocks that are separately convolved and then added together with a suitable overlap.
Each block is zero-padded at both ends and then convolved as before:

The regions originally added as zero padding are, after convolution, aligned to overlap with the unpadded ends of their respective neighbour blocks. The overlapping parts of the blocks are then added together.

Deconvolution

A signal $u(t)$ was distorted by convolution with a known impulse response $h(t)$ (e.g., through a transmission channel or a sensor problem). The “smeared” result $s(t)$ was recorded.

Can we undo the damage and restore (or at least estimate) $u(t)$?
The convolution theorem turns the problem into one of multiplication:

\[ s(t) = \int u(t - \tau) \cdot h(\tau) \cdot d\tau \]

\[ s = u \ast h \]

\[ \mathcal{F}\{s\} = \mathcal{F}\{u\} \cdot \mathcal{F}\{h\} \]

\[ \mathcal{F}\{u\} = \mathcal{F}\{s\}/\mathcal{F}\{h\} \]

\[ u = \mathcal{F}^{-1}\{\mathcal{F}\{s\}/\mathcal{F}\{h\}\} \]

In practice, we also record some noise \( n(t) \) (quantization, etc.):

\[ c(t) = s(t) + n(t) = \int u(t - \tau) \cdot h(\tau) \cdot d\tau + n(t) \]

**Problem** – At frequencies \( f \) where \( \mathcal{F}\{h\}(f) \) approaches zero, the noise will be amplified (potentially enormously) during deconvolution:

\[ \tilde{u} = \mathcal{F}^{-1}\{\mathcal{F}\{c\}/\mathcal{F}\{h\}\} = u + \mathcal{F}^{-1}\{\mathcal{F}\{n\}/\mathcal{F}\{h\}\} \]

**Typical workarounds:**

- Modify the Fourier transform of the impulse response, such that \( |\mathcal{F}\{h\}(f)| > \epsilon \) for some experimentally chosen threshold \( \epsilon \).

- If estimates of the signal spectrum \( |\mathcal{F}\{s\}(f)| \) and the noise spectrum \( |\mathcal{F}\{n\}(f)| \) can be obtained, then we can apply the “Wiener filter” ("optimal filter")

\[ W(f) = \frac{|\mathcal{F}\{s\}(f)|^2}{|\mathcal{F}\{s\}(f)|^2 + |\mathcal{F}\{n\}(f)|^2} \]

before deconvolution:

\[ \tilde{u} = \mathcal{F}^{-1}\{W \cdot \mathcal{F}\{c\}/\mathcal{F}\{h\}\} \]

**Exercise 7** Use MATLAB to deconvolve the blurred stars from slide 25. The files stars-blurred.png with the blurred-stars image and stars-psf.png with the impulse response (point-spread function) are available on the course-material web page. You may find the MATLAB functions `imread`, `double`, `imagesc`, `circshift`, `fft2`, `ifft2` of use. Try different ways to control the noise (see above) and distortions near the margins (windowing). [The MATLAB image processing toolbox provides ready-made “professional” functions `deconvwnr`, `deconvreg`, `deconvlucy`, `edgetaper`, for such tasks. Do not use these, except perhaps to compare their outputs with the results of your own attempts.]
We introduced the DFT as a special case of the continuous Fourier transform, where the input is sampled and periodic.

If the input is sampled, but not periodic, the DFT can still be used to calculate an approximation of the Fourier transform of the original continuous signal. However, there are two effects to consider. They are particularly visible when analysing pure sine waves.

Sine waves whose frequency is a multiple of the base frequency \( \frac{f_s}{n} \) of the DFT are identical to their periodic extension beyond the size of the DFT. They are therefore represented exactly by a single sharp peak in the DFT. All their energy falls into one single frequency “bin” in the DFT result.

Sine waves with other frequencies that do not match exactly one of the output frequency bins of the DFT are still represented by a peak at the output bin that represents the nearest integer multiple of the DFT’s base frequency. However, such a peak is distorted in two ways:

- Its amplitude is lower (down to 63.7%)
- Much signal energy has “leaked” to other frequencies.
The leakage of energy to other frequency bins not only blurs the estimated spectrum. The peak amplitude also changes significantly as the frequency of a tone changes from that associated with one output bin to the next, a phenomenon known as scalloping. In the above graphic, an input sine wave gradually changes from the frequency of bin 15 to that of bin 16 (only positive frequencies shown).

**Windowing**

![Sine wave](image1)

![Discrete Fourier Transform](image2)

![Sine wave multiplied with window function](image3)

![Discrete Fourier Transform](image4)
The reason for the leakage and scalloping losses is easy to visualize with the help of the convolution theorem:

The operation of cutting a sequence with the size of the DFT input vector out of a longer original signal (the one whose continuous Fourier spectrum we try to estimate) is equivalent to multiplying this signal with a rectangular function. This destroys all information and continuity outside the “window” that is fed into the DFT.

Multiplication with a rectangular window of length $T$ in the time domain is equivalent to convolution with $\sin(\pi f T)/(\pi f T)$ in the frequency domain.

The subsequent interpretation of this window as a periodic sequence by the DFT leads to sampling of this convolution result (sampling meaning multiplication with a Dirac comb whose impulses are spaced $f_s/n$ apart).

Where the window length was an exact multiple of the original signal period, sampling of the $\sin(\pi f T)/(\pi f T)$ leads to a single Dirac pulse, and the windowing causes no distortion. In all other cases, the effects of the convolution become visible in the frequency domain as leakage and scalloping losses.

Some better window functions

![Some better window functions graph](image)

All these functions are 0 outside the interval [0,1].
Numerous alternatives to the rectangular window have been proposed that reduce leakage and scalloping in spectral estimation. These are vectors multiplied element-wise with the input vector before applying the DFT to it. They all force the signal amplitude smoothly down to zero at the edge of the window, thereby avoiding the introduction of sharp jumps in the signal when it is extended periodically by the DFT.

Three examples of such window vectors \( \{w_i\}_{i=0}^{n-1} \) are:

**Triangular window** (Bartlett window):

\[
w_i = 1 - \left| 1 - \frac{i}{n/2} \right|
\]

**Hanning window** (raised-cosine window, Hann window):

\[
w_i = 0.5 - 0.5 \times \cos \left( 2\pi \frac{i}{n-1} \right)
\]

**Hamming window**:

\[
w_i = 0.54 - 0.46 \times \cos \left( 2\pi \frac{i}{n-1} \right)
\]
Zero padding increases DFT resolution

The two figures below show two spectra of the 16-element sequence

\[ s_i = \cos(2\pi \cdot 3i/16) + \cos(2\pi \cdot 4i/16), \quad i \in \{0, \ldots, 15\}. \]

The left plot shows the DFT of the windowed sequence

\[ x_i = s_i \cdot w_i, \quad i \in \{0, \ldots, 15\} \]

and the right plot shows the DFT of the zero-padded windowed sequence

\[ x'_i = \begin{cases} s_i \cdot w_i, & i \in \{0, \ldots, 15\} \\ 0, & i \in \{16, \ldots, 63\} \end{cases} \]

where \( w_i = 0.54 - 0.46 \times \cos(2\pi i/15) \) is the Hamming window.

Applying the discrete Fourier transform to an \( n \)-element long real-valued sequence leads to a spectrum consisting of only \( n/2 + 1 \) discrete frequencies.

Since the resulting spectrum has already been distorted by multiplying the (hypothetically longer) signal with a windowing function that limits its length to \( n \) non-zero values and forces the waveform smoothly down to zero at the window boundaries, appending further zeros outside the window will not distort the signal further.

The frequency resolution of the DFT is the sampling frequency divided by the block size of the DFT. Zero padding can therefore be used to increase the frequency resolution of the DFT.

Note that zero padding does not add any additional information to the signal. The spectrum has already been “low-pass filtered” by being convolved with the spectrum of the windowing function. Zero padding in the time domain merely samples this spectrum blurred by the windowing step at a higher resolution, thereby making it easier to visually distinguish spectral lines and to locate their peak more precisely.
Frequency inversion

In order to turn the spectrum $X(f)$ of a real-valued signal $x_i$ sampled at $f_s$ into an inverted spectrum $X'(f) = X(f_s/2 - f)$, we merely have to shift the periodic spectrum by $f_s/2$:

This can be accomplished by multiplying the sampled sequence $x_i$ with $y_i = \cos \pi f_s t = \cos \pi i$, which is nothing but multiplication with the sequence

$$\ldots, 1, -1, 1, -1, 1, -1, 1, -1, \ldots$$

So in order to design a discrete high-pass filter that attenuates all frequencies $f$ outside the range $f_c < |f| < f_s/2$, we merely have to design a low-pass filter that attenuates all frequencies outside the range $-f_c < f < f_c$, and then multiply every second value of its impulse response with $-1$.

Window-based design of FIR filters

Recall that the ideal continuous low-pass filter with cut-off frequency $f_c$ has the frequency characteristic

$$H(f) = \begin{cases} 
1 & \text{if } -f_c < f \leq f_c \\
0 & \text{otherwise}
\end{cases}$$

and the impulse response

$$h(t) = \frac{\sin 2\pi t f_c}{2\pi t f_c}.$$ 

Sampling this impulse response with the sampling frequency $f_s$ of the signal to be processed will lead to a periodic frequency characteristic, that matches the periodic spectrum of the sampled signal.

There are two problems though:

$\rightarrow$ the impulse response is infinitely long

$\rightarrow$ this filter is not causal, that is $h(t) \neq 0$ for $t < 0$
Solutions:

→ Make the impulse response finite by multiplying the sampled \( h(t) \) with a windowing function

→ Make the impulse response causal by adding a delay of half the window size

The impulse response of an \( n \)-th order low-pass filter is then chosen as

\[
h_i = \frac{\sin[2\pi(i - n/2)f_c/f_s]}{2\pi(i - n/2)f_c/f_s} \cdot w_i
\]

where \( \{w_i\} \) is a windowing sequence, such as the Hamming window

\[
w_i = 0.54 - 0.46 \times \cos\left(\frac{2\pi i}{n}\right)
\]

with \( w_i = 0 \) for \( i < 0 \) and \( i > n \).

Note that for \( f_c = f_s/4 \), we have \( h_i = 0 \) for all even values of \( i \). Therefore, this special case requires only half the number of multiplications during the convolution. Such “half-band” FIR filters are used, for example, as anti-aliasing filters wherever a sampling rate needs to be halved.

---

**FIR lowpass filter design example**

order: 30, cutoff frequency (-6 dB): \( 0.25 \times f_s/2 \)
We truncate the ideal, infinitely-long impulse response by multiplication with a window sequence. In the frequency domain, this will convolve the rectangular frequency response of the ideal low-pass filter with the frequency characteristic of the window. The width of the main lobe determines the width of the transition band, and the side lobes cause ripples in the passband and stopband.

To design a passband filter that attenuates all frequencies \( f \) outside the range \( f_l < f < f_h \), we first design a low-pass filter with a cut-off frequency \( (f_h - f_l)/2 \) and multiply its impulse response with a sine wave of frequency \( (f_h + f_l)/2 \), before applying the usual windowing:

\[
h_i = \frac{\sin[\pi(i - n/2)(f_h - f_l)/f_s]}{\pi(i - n/2)(f_h - f_l)/f_s} \cdot \sin[\pi(f_h + f_l)] \cdot w_i
\]

**Exercise 8** Explain the difference between the DFT, FFT, and FFTW.

**Exercise 9** Push-button telephones use a combination of two sine tones to signal, which button is currently being pressed:

<table>
<thead>
<tr>
<th>Button Key</th>
<th>1209 Hz</th>
<th>1336 Hz</th>
<th>1477 Hz</th>
<th>1633 Hz</th>
</tr>
</thead>
<tbody>
<tr>
<td>697 Hz</td>
<td>1</td>
<td>2</td>
<td>3</td>
<td>A</td>
</tr>
<tr>
<td>770 Hz</td>
<td>4</td>
<td>5</td>
<td>6</td>
<td>B</td>
</tr>
<tr>
<td>852 Hz</td>
<td>7</td>
<td>8</td>
<td>9</td>
<td>C</td>
</tr>
<tr>
<td>941 Hz</td>
<td>*</td>
<td>0</td>
<td>#</td>
<td>D</td>
</tr>
</tbody>
</table>

(a) You receive a digital telephone signal with a sampling frequency of 8 kHz. You cut a 256-sample window out of this sequence, multiply it with a windowing function and apply a 256-point DFT. What are the indices where the resulting vector \( (X_0, X_1, \ldots, X_{255}) \) will show the highest amplitude if button [9] was pushed at the time of the recording?

(b) Use MATLAB to determine, which button sequence was typed in the touch tones recorded in the file touchtone.wav on the course-material web page.
Polynomial representation of sequences

We can represent sequences \( \{x_n\} \) as polynomials:

\[
X(v) = \sum_{n=-\infty}^{\infty} x_n v^n
\]

Example of polynomial multiplication:

\[
\begin{align*}
(1 + 2v + 3v^2) \cdot (2 + 1v) &= 2 + 4v + 6v^2 \\
&\quad + 1v + 2v^2 + 3v^3 \\
&= 2 + 5v + 8v^2 + 3v^3
\end{align*}
\]

Compare this with the convolution of two sequences (in MATLAB):

\[
\text{conv([1 2 3], [2 1]) equals [2 5 8 3]}
\]

Convolution of sequences is equivalent to polynomial multiplication:

\[
\{h_n\} \ast \{x_n\} = \{y_n\} \Rightarrow y_n = \sum_{k=-\infty}^{\infty} h_k \cdot x_{n-k}
\]

\[
H(v) \cdot X(v) = \left( \sum_{n=-\infty}^{\infty} h_n v^n \right) \cdot \left( \sum_{n=-\infty}^{\infty} x_n v^n \right) = \sum_{n=-\infty}^{\infty} \sum_{k=-\infty}^{\infty} h_k \cdot x_{n-k} \cdot v^n
\]

Note how the Fourier transform of a sequence can be accessed easily from its polynomial form:

\[
X(e^{-j\omega}) = \sum_{n=-\infty}^{\infty} x_n e^{-j\omega n}
\]
Example of polynomial division:

\[ \frac{1}{1 - av} = 1 + av + a^2v^2 + a^3v^3 + \cdots = \sum_{n=0}^{\infty} a^n v^n \]

\[
\begin{array}{c|cccc}
& 1 & + & av & + a^2v^2 & + \cdots \\
\hline
1 - av & 1 & - av \\
\end{array}
\]

\[
\begin{array}{cccc}
& av & - a^2v^2 \\
av & a^2v^2 & - a^3v^3 & \cdots \\
\end{array}
\]

Rational functions (quotients of two polynomials) can provide a convenient closed-form representations for infinitely-long exponential sequences, in particular the impulse responses of IIR filters.

**The z-transform**

The z-transform of a sequence \( \{x_n\} \) is defined as:

\[ X(z) = \sum_{n=-\infty}^{\infty} x_n z^{-n} \]

Note that it differs only in the sign of the exponent from the polynomial representation discussed on the preceding slides.

Recall that the above \( X(z) \) is exactly the factor with which an exponential sequence \( \{z^n\} \) is multiplied, if it is convolved with \( \{x_n\} \):

\[
\{z^n\} \ast \{x_n\} = \{y_n\}
\]

\[
\Rightarrow \quad y_n = \sum_{k=\infty}^{\infty} z^{n-k} x_k = z^n \cdot \sum_{k=\infty}^{\infty} z^{-k} x_k = z^n \cdot X(z)
\]
The z-transform defines for each sequence a continuous complex-valued surface over the complex plane \( \mathbb{C} \). For finite sequences, its value is always defined across the entire complex plane.

For infinite sequences, it can be shown that the z-transform converges only for the region

\[
\lim_{n \to \infty} \frac{|x_{n+1}|}{|x_n|} < |z| < \lim_{n \to -\infty} \frac{|x_{n+1}|}{|x_n|}
\]

The z-transform identifies a sequence unambiguously only in conjunction with a given region of convergence. In other words, there exist different sequences, that have the same expression as their z-transform, but that converge for different amplitudes of \( z \).

The z-transform is a generalization of the Fourier transform, which it contains on the complex unit circle (\(|z| = 1\)):

\[
\mathcal{F}\{x_n\}(\omega) = X(e^{j\omega}) = \sum_{n=-\infty}^{\infty} x_n e^{-j\omega n}
\]

The z-transform of the impulse response \( \{h_n\} \) of the causal LTI system defined by

\[
\sum_{k=0}^{N} a_k \cdot y_{n-k} = \sum_{m=0}^{M} b_m \cdot x_{n-m}
\]

with \( \{y_n\} = \{h_n\} \ast \{x_n\} \) is the rational function

\[
H(z) = \frac{b_0 + b_1 z^{-1} + b_2 z^{-2} + \cdots + b_m z^{-m}}{a_0 + a_1 z^{-1} + a_2 z^{-2} + \cdots + a_k z^{-k}}
\]

which can also be written as

\[
H(z) = \frac{z^n \sum_{l=0}^{m} b_l z^{m-l}}{z^m \sum_{l=0}^{n} a_l z^{n-l}}
\]

\( H(z) \) has \( m \) zeros and \( n \) poles at non-zero locations in the \( z \) plane, plus \( n - m \) zeros (if \( n > m \)) or \( m - n \) poles (if \( m > n \)) at \( z = 0 \).
This function can be converted into the form

\[ H(z) = \frac{b_0}{a_0} \cdot \frac{\prod_{l=1}^{m} (1 - c_l \cdot z^{-1})}{\prod_{l=1}^{n} (1 - d_l \cdot z^{-1})} = \frac{b_0}{a_0} \cdot z^{n-m} \cdot \frac{\prod_{l=1}^{m} (z - c_l)}{\prod_{l=1}^{n} (z - d_l)} \]

where the \(c_l\) are the non-zero positions of zeros (\(H(c_l) = 0\)) and the \(d_l\) are the non-zero positions of the poles (i.e., \(z \rightarrow d_l \Rightarrow |H(z)| \rightarrow \infty\)) of \(H(z)\). Except for a constant factor, \(H(z)\) is entirely characterized by the position of these zeros and poles.

As with the Fourier transform, convolution in the time domain corresponds to complex multiplication in the \(z\)-domain:

\[ \{x_n\} \bullet - \circ X(z), \{y_n\} \bullet - \circ Y(z) \Rightarrow \{x_n\} \ast \{y_n\} \bullet - \circ X(z) \cdot Y(z) \]

Delaying a sequence by one corresponds in the \(z\)-domain to multiplication with \(z^{-1}\):

\[ \{x_{n-\Delta n}\} \bullet - \circ X(z) \cdot z^{-\Delta n} \]

This example is an amplitude plot of \(H(z) = 0.8 \frac{z}{1 - 0.2 \cdot z^{-1}}\) which features a zero at 0 and a pole at 0.2.
\[ H(z) = \frac{z}{z-0.7} = \frac{1}{1-0.7z^{-1}} \]

\[ H(z) = \frac{z}{z-0.9} = \frac{1}{1-0.9z^{-1}} \]

\[ H(z) = \frac{z}{z-1} = \frac{1}{1-z^{-1}} \]

\[ H(z) = \frac{z}{z-1.1} = \frac{1}{1-1.1z^{-1}} \]
\[ H(z) = \frac{z^2}{(z-0.9e^{j\pi/6})(z-0.9e^{-j\pi/6})} = \frac{1}{1-1.8\cos(\pi/6)z^{-1}+0.9^2z^{-2}} \]

\[ H(z) = \frac{z^2}{(z-e^{j\pi/6})(z-e^{-j\pi/6})} = \frac{1}{1-2\cos(\pi/6)z^{-1}+z^{-2}} \]

\[ H(z) = \frac{z^2}{(z-0.9e^{j\pi/2})(z-0.9e^{-j\pi/2})} = \frac{1}{1-1.8\cos(\pi/2)z^{-1}+0.9^2z^{-2}} = \frac{1}{1+0.9^2z^{-2}} \]

\[ H(z) = \frac{z}{z+1} = \frac{1}{1+z^{-1}} \]
IIR Filter design techniques

The design of a filter starts with specifying the desired parameters:

→ The passband is the frequency range where we want to approximate a gain of one.

→ The stopband is the frequency range where we want to approximate a gain of zero.

→ The order of a filter is the number of poles it uses in the $z$-domain, and equivalently the number of delay elements necessary to implement it.

→ Both passband and stopband will in practice not have gains of exactly one and zero, respectively, but may show several deviations from these ideal values, and these ripples may have a specified maximum quotient between the highest and lowest gain.

→ There will in practice not be an abrupt change of gain between passband and stopband, but a transition band where the frequency response will gradually change from its passband to its stopband value.

The designer can then trade off conflicting goals such as a small transition band, a low order, a low ripple amplitude, or even an absence of ripples.

Design techniques for making these tradeoffs for analog filters (involving capacitors, resistors, coils) can also be used to design digital IIR filters:

**Butterworth filters**

Have no ripples, gain falls monotonically across the pass and transition band. Within the passband, the gain drops slowly down to $1 - \sqrt{1/2}$ ($-3$ dB). Outside the passband, it drops asymptotically by a factor $2^N$ per octave ($N \cdot 20$ dB/decade).
Chebyshev type I filters
Distribute the gain error uniformly throughout the passband (equiripples) and drop off monotonically outside.

Chebyshev type II filters
Distribute the gain error uniformly throughout the stopband (equiripples) and drop off monotonically in the passband.

Elliptic filters (Cauer filters)
Distribute the gain error as equiripples both in the passband and stopband. This type of filter is optimal in terms of the combination of the passband-gain tolerance, stopband-gain tolerance, and transition-band width that can be achieved at a given filter order.

All these filter design techniques are implemented in the MATLAB Signal Processing Toolbox in the functions butter, cheby1, cheby2, and ellip, which output the coefficients $a_n$ and $b_n$ of the difference equation that describes the filter. These can be applied with filter to a sequence, or can be visualized with zplane as poles/zeros in the $z$-domain, with impz as an impulse response, and with freqz as an amplitude and phase spectrum. The commands sptool and fdatool provide interactive GUIs to design digital filters.

Butterworth filter design example

order: 1, cutoff frequency ($-3$ dB): $0.25 \times f_s/2$
Butterworth filter design example

order: 5, cutoff frequency (−3 dB): 0.25 × \( f_s/2 \)

Chebyshev type I filter design example

order: 5, cutoff frequency: 0.5 × \( f_s/2 \), pass-band ripple: −3 dB
Chebyshev type II filter design example

order: 5, cutoff frequency: $0.5 \times f_s/2$, stop-band ripple: $-20$ dB

Elliptic filter design example

order: 5, cutoff frequency: $0.5 \times f_s/2$, pass-band ripple: $-3$ dB, stop-band ripple: $-20$ dB
**Exercise 10** Draw the direct form II block diagrams of the causal infinite-impulse response filters described by the following $z$-transforms and write down a formula describing their time-domain impulse responses:

(a) $H(z) = \frac{1}{1 - \frac{1}{2}z^{-1}}$

(b) $H'(z) = \frac{1 - \frac{1}{4}z^{-4}}{1 - \frac{1}{4}z^{-1}}$

(c) $H''(z) = \frac{1}{2} + \frac{1}{4}z^{-1} + \frac{1}{2}z^{-2}$

**Exercise 11** (a) Perform the polynomial division of the rational function given in exercise 10 (a) until you have found the coefficient of $z^{-5}$ in the result.

(b) Perform the polynomial division of the rational function given in exercise 10 (b) until you have found the coefficient of $z^{-10}$ in the result.

(c) Use its $z$-transform to show that the filter in exercise 10 (b) has actually a finite impulse response and draw the corresponding block diagram.

**Exercise 12** Consider the system $h : \{x_n\} \rightarrow \{y_n\}$ with $y_n + y_{n-1} = x_n - x_{n-4}$.

(a) Draw the direct form I block diagram of a digital filter that realises $h$.

(b) What is the impulse response of $h$?

(c) What is the step response of $h$ (i.e., $h * u$)?

(d) Apply the $z$-transform to (the impulse response of) $h$ to express it as a rational function $H(z)$.

(e) Can you eliminate a common factor from nominator and denominator? What does this mean?

(f) For what values $z \in \mathbb{C}$ is $H(z) = 0$?

(g) How many poles does $H$ have in the complex plane?

(h) Write $H$ as a fraction using the position of its poles and zeros and draw their location in relation to the complex unit circle.

(i) If $h$ is applied to a sound file with a sampling frequency of 8000 Hz, sine waves of what frequency will be eliminated and sine waves of what frequency will be quadrupled in their amplitude?
Random sequences and noise

A discrete random sequence \( \{x_n\} \) is a sequence of numbers

\[ \ldots, x_{-2}, x_{-1}, x_0, x_1, x_2, \ldots \]

where each value \( x_n \) is the outcome of a random variable \( x_n \) in a corresponding sequence of random variables

\[ \ldots, X_{-2}, X_{-1}, X_0, X_1, X_2, \ldots \]

Such a collection of random variables is called a random process. Each individual random variable \( x_n \) is characterized by its probability distribution function

\[ P_{x_n}(a) = \text{Prob}(x_n \leq a) \]

and the entire random process is characterized completely by all joint probability distribution functions

\[ P_{x_{n_1}, \ldots, x_{n_k}}(a_1, \ldots, a_k) = \text{Prob}(x_{n_1} \leq a_1 \land \ldots \land x_{n_k} \leq a_k) \]

for all possible sets \( \{x_{n_1}, \ldots, x_{n_k}\} \).

Two random variables \( x_n \) and \( x_m \) are called independent if

\[ P_{x_n, x_m}(a, b) = P_{x_n}(a) \cdot P_{x_m}(b) \]

and a random process is called stationary if

\[ P_{x_{n_1+l}, \ldots, x_{n_k+l}}(a_1, \ldots, a_k) = P_{x_{n_1}, \ldots, x_{n_k}}(a_1, \ldots, a_k) \]

for all \( l \), in other words, if the probability distributions are time invariant.

The derivative \( p_{x_n}(a) = P'_{x_n}(a) \) is called the probability density function, and helps us to define quantities such as

\[ \rightarrow \text{the expected value } \mathcal{E}(x_n) = \int a p_{x_n}(a) \, da \]

\[ \rightarrow \text{the mean-square value} \ (\text{average power}) \ \mathcal{E}(\|x_n\|^2) \]

\[ \rightarrow \text{the variance} \ \text{Var}(x_n) = \mathcal{E}[\|x_n - \mathcal{E}(x_n)\|^2] = \mathcal{E}(\|x_n\|^2) - \|\mathcal{E}(x_n)\|^2 \]

\[ \rightarrow \text{the correlation} \ \text{Cor}(x_n, x_m) = \mathcal{E}(x_n \cdot x_m^*) \]

Remember that \( \mathcal{E}(\cdot) \) is linear, that is \( \mathcal{E}(ax) = a\mathcal{E}(x) \) and \( \mathcal{E}(x + y) = \mathcal{E}(x) + \mathcal{E}(y) \). Also, \( \text{Var}(ax) = a^2\text{Var}(x) \) and, if \( x \) and \( y \) are independent, \( \text{Var}(x + y) = \text{Var}(x) + \text{Var}(y) \).
A stationary random process \( \{x_n\} \) can be characterized by its mean value
\[
m_x = \mathbb{E}(x_n),
\]
its variance
\[
\sigma_x^2 = \mathbb{E}(|x_n - m_x|^2) = \gamma_{xx}(0)
\]
(\( \sigma_x \) is also called standard deviation), its autocorrelation sequence
\[
\phi_{xx}(k) = \mathbb{E}(x_{n+k} \cdot x_n^*)
\]
and its autocovariance sequence
\[
\gamma_{xx}(k) = \mathbb{E}[(x_{n+k} - m_x) \cdot (x_n - m_x)^*] = \phi_{xx}(k) - |m_x|^2
\]
A stationary pair of random processes \( \{x_n\} \) and \( \{y_n\} \) can, in addition, be characterized by its crosscorrelation sequence
\[
\phi_{xy}(k) = \mathbb{E}(x_{n+k} \cdot y_n^*)
\]
and its crosscovariance sequence
\[
\gamma_{xy}(k) = \mathbb{E}[(x_{n+k} - m_x) \cdot (y_n - m_y)^*] = \phi_{xy}(k) - m_x m_y^*
\]

### Deterministic crosscorrelation sequence

For deterministic sequences \( \{x_n\} \) and \( \{y_n\} \), the crosscorrelation sequence is
\[
c_{xy}(k) = \sum_{i=-\infty}^{\infty} x_{i+k} y_i.
\]
After dividing through the overlapping length of the finite sequences involved, \( c_{xy}(k) \) can be used to estimate from a finite sample of a stationary random sequence the underlying \( \phi_{xy}(k) \). MATLAB’s \texttt{xcorr} function does that with option \texttt{unbiased}.

If \( \{x_n\} \) is similar to \( \{y_n\} \), but lags \( l \) elements behind \( (x_n \approx y_{n-l}) \), then \( c_{xy}(l) \) will be a peak in the crosscorrelation sequence. It is therefore widely calculated to locate shifted versions of a known sequence in another one.

The crosscorrelation is a close cousin of the convolution, with just the second input sequence mirrored:
\[
\{c_{xy}(n)\} = \{x_n\} * \{y_{-n}\}
\]
It can therefore be calculated equally easily via the Fourier transform:
\[
C_{xy}(f) = X(f) \cdot Y^*(f)
\]
Swapping the input sequences mirrors the output sequence: \( c_{xy}(k) = c_{yx}(-k) \)
Equivalently, we define the autocorrelation sequence in the time-domain as
\[ c_{xx}(k) = \sum_{i=-\infty}^{\infty} x_{i+k} x_i. \]
which corresponds in the frequency domain to
\[ C_{xx}(f) = X(f) \cdot X^*(f) = |X(f)|^2. \]

In other words, the Fourier transform \( C_{xx}(f) \) of the autocorrelation sequence \( \{c_{xx}(n)\} \) of a sequence \( \{x_n\} \) is identical to the squared amplitudes of the Fourier transform or power spectrum of \( \{x_n\} \).

This suggests, that the Fourier transform of the autocorrelation sequence of a random process might be a suitable way for defining the power spectrum of that random process.

**Filtered random sequences**

Let \( \{x_n\} \) be a random sequence from a stationary random process. The output
\[ y_n = \sum_{k=-\infty}^{\infty} h_k \cdot x_{n-k} = \sum_{k=-\infty}^{\infty} h_{n-k} \cdot x_k \]
of an LTI applied to it will then be another random sequence, characterized by
\[ m_y = m_x \sum_{k=-\infty}^{\infty} h_k \]
and
\[ \phi_{yy}(k) = \sum_{i=-\infty}^{\infty} \phi_{xx}(k-i)c_{hh}(i), \quad \text{where} \quad \phi_{xx}(k) = \mathcal{E}(x_{n+k} \cdot x_n^*) \]
\[ c_{hh}(k) = \sum_{i=-\infty}^{\infty} h_{i+k} h_i. \]
In other words:

\[
\{ y_n \} = \{ h_n \} \ast \{ x_n \} \quad \Rightarrow \quad \{ \phi_{yy}(n) \} = \{ c_{hh}(n) \} \ast \{ \phi_{xx}(n) \}
\]

\[
\Phi_{yy}(f) = |H(f)|^2 \cdot \Phi_{xx}(f)
\]

Similarly:

\[
\{ y_n \} = \{ h_n \} \ast \{ x_n \} \quad \Rightarrow \quad \{ \phi_{xy}(n) \} = \{ h_n \} \ast \{ \phi_{xx}(n) \}
\]

\[
\Phi_{xy}(f) = H(f) \cdot \Phi_{xx}(f)
\]

**White noise**

A random sequence \( \{ x_n \} \) is a *white noise* signal, if \( m_x = 0 \) and

\[
\phi_{xx}(k) = \sigma_x^2 \delta_k.
\]

The power spectrum of a white noise signal is flat:

\[
\Phi_{xx}(f) = \sigma_x^2.
\]

Application example:

Where an LTI \( \{ y_n \} = \{ h_n \} \ast \{ x_n \} \) can be observed to operate on white noise \( \{ x_n \} \) with \( \phi_{xx}(k) = \sigma_x^2 \delta_k \), the crosscorrelation between input and output will reveal the impulse response of the system:

\[
\phi_{xy}(n) = \sigma_x^2 \cdot h_n.
\]
Averaging and noise reduction

Often an original signal \( \{x_i\} \) is only accessible with some added noise
\[
\{y_i\} = \{x_i\} + \{n_i\}
\]
which turns a deterministic sequence into a random sequence. The signal-to-noise ratio (SNR) of the received signal \( \{y_i\} \) is the square root of the power ratio of these components:
\[
\text{SNR}_y = \sqrt{\mathcal{E}(|x_i|^2)/\mathcal{E}(|n_i|^2)}.
\]

As an SNR might also be given in terms of a power ratio, it is commonly expressed in decibels, to avoid any confusion between power and voltage ratios:
\[
10 \text{ dB} \cdot \log_{10} \frac{\mathcal{E}(|x_i|^2)}{\mathcal{E}(|n_i|^2)} = 20 \text{ dB} \cdot \log_{10} \sqrt{\mathcal{E}(|x_i|^2)/\mathcal{E}(|n_i|^2)}.
\]

The simplest noise reduction technique is averaging. If we know that the \( k \) signal values \( x_1, \ldots , x_k \) are identical, and the noise average is \( m_n = 0 \), then we can calculate an average value
\[
\bar{y} = \frac{1}{k} \sum_{i=1}^{k} y_i
\]
as an approximation for the true value \( x_1 = \cdots = x_k \).

What noise level remains after averaging?

The \( k \) identical signal values \( x_1, \ldots , x_k \) are characterized by \( m_x = x_i \) \((i = 1, \ldots , k)\) and \( \sigma_x^2 = 0 \). The average signal power is \( \mathcal{E}(|x_i|^2) = m_x^2 \).

We assume that the \( k \) noise values \( n_1, \ldots , n_k \) are statistically independent and are the output of a stationary process with mean \( m_n = 0 \) and variance \( \sigma_n^2 \). The average noise power is \( \mathcal{E}(|n_i|^2) = \sigma_n^2 \).

The averaging result \( \bar{y} \) can be split up into a signal component \( \bar{x} \) and a noise component \( \bar{n} \):
\[
\bar{y} = \frac{1}{k} \sum_{i=1}^{k} y_i = \frac{1}{k} \sum_{i=1}^{k} x_i + \frac{1}{k} \sum_{i=1}^{k} n_i = \bar{x} + \bar{n}
\]
The corresponding average power values are
\[
\mathcal{E}(|\bar{x}|^2) = \mathcal{E} \left( \left| \frac{1}{k} \sum_{i=1}^{k} x_i \right|^2 \right) = \mathcal{E} \left( \left| \frac{1}{k} \sum_{i=1}^{k} m_x \right|^2 \right) = m_x^2
\]
and

\[ E(|\bar{n}|^2) = \text{Var}(\bar{n}) = \text{Var}\left(\frac{1}{k} \sum_{i=1}^{k} n_i\right) = \frac{1}{k^2} \sum_{i=1}^{k} \text{Var}(n_i) = \frac{1}{k} \sigma_n^2. \]

We can now compare the signal-to-noise ratio of the original noisy sequence

\[ \text{SNR}_y = \sqrt{\frac{E(|x_i|^2)}{E(|n_i|^2)}} = \frac{m_x}{\sigma_n} \]

with that of the averaging result

\[ \text{SNR}_{\bar{y}} = \sqrt{\frac{E(|\bar{x}|^2)}{E(|\bar{n}|^2)}} = \sqrt{k} \cdot \frac{m_x}{\sigma_n}. \]

Averaging \( k \) samples of identical signal values with added independent zero-mean noise values will increase the signal-to-noise ratio by the factor \( \sqrt{k} \).

Remember that adding identical values \( x_1 \) and \( x_2 \) will double their value and therefore quadruple their power. On the other hand, adding independent zero-mean noise values \( n_1 \) and \( n_2 \) will – on average – only double their power (\( = \text{variance} \)).

### Noise reduction filters

Added independent noise values with \( \phi_{nn}(k) = \sigma_n^2 \delta_k \) have a flat spectrum \( \Phi_{nn}(f) = \sigma_n^2 \), which is added in \( \{y_i\} = \{x_i\} + \{n_i\} \) across the spectrum of the noise-free signal \( \{x_i\} \).

Knowledge of the power spectrum \( X(f) \) of the original signal can help to design a filter that attenuates much of the noise in \( \{y_i\} \), without degrading the wanted signal \( \{x_i\} \) too much. If \( \{x_i\} \) changes very slowly, most of its energy will be at low frequencies, and noise can be reduced with a suitably chosen low-pass filter.

If there are no significant changes in \( \{x_i\} \) during \( k \) consecutive samples, then convolution with a \( k \) samples wide rectangular window will not affect the wanted signal.

We have already seen that this will reduce the noise amplitude by a factor \( \sqrt{k} \), or equivalently the noise power by a factor of \( k \).
How does a general LTI filter
\[
\bar{y}_i = \sum_{k=-\infty}^{\infty} h_k \cdot y_{i-k} = \sum_{k=-\infty}^{\infty} h_k \cdot (x_{i-k} + n_{i-k}) = \bar{x}_i + \bar{n}_i
\]
affect the signal-to-noise ratio \( \text{SNR}_y = \sqrt{\mathcal{E}(|x_i|^2)/\mathcal{E}(|n_i|^2)} \)?
With \( m_x \approx x_i \), we get
\[
\mathcal{E}(|\bar{x}_i|^2) = \mathcal{E} \left( \left| \sum_{k=-\infty}^{\infty} h_k x_{i-k} \right|^2 \right) = m_x^2 \left| \sum_{k=-\infty}^{\infty} h_k \right|^2
\]
and
\[
\mathcal{E}(|\bar{n}_i|^2) = \text{Var}(\bar{n}_i) = \text{Var} \left( \sum_{k=-\infty}^{\infty} h_k n_{i-k} \right) = \sum_{k=-\infty}^{\infty} h_k^2 \cdot \text{Var}(n_{i-k}) = \sigma_n^2 \sum_{k=-\infty}^{\infty} h_k^2
\]
\[\Rightarrow \text{SNR}_{\bar{y}_i}/\text{SNR}_{y_i} = \left| \sum_{k=-\infty}^{\infty} h_k \right| / \sqrt{\sum_{k=-\infty}^{\infty} h_k^2}\]

**Exponential averaging**

A particularly easy to implement IIR low-pass filter, with only one delay element, is the exponential averaging filter
\[
\tilde{y}_i = (1 - \alpha)y_i + \alpha \tilde{y}_{i-1} = (1 - \alpha)\sum_{k=0}^{\infty} \alpha^k y_{i-k}, \quad (0 \leq \alpha < 1).
\]

When applied as a noise filter on a sequence \( \{y_i\} = \{x_i\} + \{n_i\} \) with \( \sigma_x^2 = 0, m_n = 0, \) and mutually independent noise values \( \{n_i\} \), we get
\[
\sigma_{\tilde{y}_i}^2 = (1 - \alpha)^2 \sum_{k=0}^{\infty} \alpha^{2k} \sigma_{y_{i-k}}^2 = \sigma_n^2 (1 - \alpha)^2 \sum_{k=0}^{\infty} \alpha^{2k} = \sigma_n^2 (1 - \alpha)^2 \cdot \frac{1}{1 - \alpha^2} = \sigma_n^2 \cdot \frac{1 - \alpha}{1 + \alpha}
\]
and therefore
\[
\text{SNR}_{\tilde{y}_i}/\text{SNR}_{y_i} = \frac{m_x}{\sigma_{\tilde{y}_i}} / \frac{m_x}{\sigma_n} = \sqrt{\frac{1 + \alpha}{1 - \alpha}}.
\]
The above diagrams show different type of spectral estimates of a sequence 
\[ x_i = \sin(2\pi j \times 8/64) + \sin(2\pi j \times 14.32/64) + n_i \] 
with \( \phi_{nn}(i) = 4\delta_i \).

Left is a single 64-element DFT of \( \{x_i\} \) (with rectangular window). The flat spectrum of white noise is only an expected value. In a single discrete Fourier transform of such a sequence, the significant variance of the noise spectrum becomes visible. It almost drowning the two peaks from sine waves.

After cutting \( \{x_i\} \) into 1000 windows of 64 elements each, calculating their DFT, and plotting the average of their absolute values, the centre figure shows an approximation of the expected value of the amplitude spectrum, with a flat noise floor. Taking the absolute value before averaging is called \textit{incoherent averaging}, as the phase information is thrown away.

The rightmost figure was generated from the same set of 1000 windows, but this time the complex values of the DFTs were averaged before the absolute value was taken. This is called \textit{coherent averaging} and, because of the linearity of the DFT, identical to first averaging the 1000 windows and then applying a single DFT and taking its absolute value. The windows start 64 samples apart. Only periodic waveforms with a period that divides 64 are not averaged away. This periodic averaging step suppresses both the noise and the second sine wave.

**Periodic averaging**

If a zero-mean signal \( \{x_i\} \) has a periodic component with period \( p \), the periodic component can be isolated by \textit{periodic averaging}:

\[ \bar{x}_i = \lim_{k \to \infty} \frac{1}{2k + 1} \sum_{n=-k}^{k} x_{i+pn} \]

Periodic averaging corresponds in the time domain to convolution with a Dirac comb \( \sum_n \delta_{i-pn} \). In the frequency domain, this means multiplication with a Dirac comb that eliminates all frequencies but multiples of \( 1/p \).
Exercise 13  A tumble dryer measures the air humidity $x_i$ in the drum twice per second, but the control software uses only a smoothed value $y_i = \sum_{k=0}^{15} x_{i-k}$ to decide when to stop the drying programme.

(a) Assuming that humidity fluctuations in a tumble dryer have components that are periodic with the rotational frequency of the drum, for which rotational speeds would the above smoothing filter suppress these periodic components particularly well?

(b) It is the year 2020, and a cartel of semiconductor manufacturers recently managed to double the cost of memory every 18 months. Your job is to redesign products to eliminate unnecessary use of precious memory. Replace the above data smoothing algorithm with one that requires the least amount of memory possible, but that reduces the random fluctuations caused in the measurements by tumbling clothes by exactly the same factor. You now can assume that these fluctuations are not correlated between different measurements $x_i$.

(c) First predict the approximate shape of the frequency characteristic of both the old and the new smoothing filter, then use MATLAB to plot it.

Outlook

Further topics that we have not covered in this brief introductory tour through DSP, but for the understanding of which you should now have a good theoretical foundation:

→ multirate systems
→ orthogonal filterbanks and wavelets
→ effects of rounding errors
→ adaptive filters
→ DSP hardware architectures
→ modulation and symbol detection techniques
→ sound effects
→ audio-visual coding techniques

Your feedback on the course is very much appreciated:

http://www.cl.cam.ac.uk/cgi-bin/lr/login