Microphone-Array Signal Processing

José A. Apolinário Jr. and Marcello L. R. de Campos

{apolin}, {mcampos}@ieee.org



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1. Introduction and Fundamentals



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- 2. Sensor Arrays and Spatial Filtering

Outline

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- 5. DoA Estimation with Microphone Arrays

1. Introduction and Fundamentals

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1.2 Signals in Space and Time

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$$= \begin{bmatrix} \frac{\partial(\cdot)}{\partial x} & \frac{\partial(\cdot)}{\partial y} & \frac{\partial(\cdot)}{\partial z} \end{bmatrix}^{T}$$

and

$$\nabla_{\mathbf{x}}^{2}(\cdot) = \frac{\partial^{2}(\cdot)}{\partial x^{2}} + \frac{\partial^{2}(\cdot)}{\partial y^{2}} + \frac{\partial^{2}(\cdot)}{\partial z^{2}}$$

respectively.

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or, for $s(\mathbf{x}, t)$ a general scalar field,

$$\frac{\partial^2 s}{\partial x^2} + \frac{\partial^2 s}{\partial y^2} + \frac{\partial^2 s}{\partial z^2} = \frac{1}{c^2} \frac{\partial^2 s}{\partial t^2}$$

where: *c* is the propagation speed, \overrightarrow{E} is the electric field intensity, and $\mathbf{x} = [x \ y \ z]^T$ is a position vector.

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"vector" wave equation

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Note: From this point onwards the terms wave and field will be used interchangeably.

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$$s(\mathbf{x},t) = Ae^{j(\omega t - k_x x - k_y y - k_z z)}$$

where A is a complex constant and k_x , k_y , k_z , and $\omega \ge 0$ are real constants.

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<u>constraints</u> to be satisfied by the parameters of the scalar field

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 $k_x x + k_y y + k_z z = C$

where C is a constant.

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The planes where $s(\mathbf{x}, t)$ is constant are perpendicular to the wavenumber vector \mathbf{k}

As the plane wave propagates, it advances a distance δx in δt seconds.

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Therefore,

$$s(\mathbf{x}, t) = s(\mathbf{x} + \delta \mathbf{x}, t + \delta t)$$
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$$\implies \omega \delta t - \mathbf{k}^T \delta \mathbf{x} = 0$$

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After $T = 2\pi/\omega$ seconds, the plane wave has completed one cycle and it appears as it did before, but its *wavefront* has advanced a distance of one *wavelength*, λ . After $T = 2\pi/\omega$ seconds, the plane wave has completed one cycle and it appears as it did before, but its *wavefront* has advanced a distance of one *wavelength*, λ .

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The wavenumber vector, \mathbf{k} , may be considered a spatial frequency variable, just as ω is a temporal frequency variable.

We may rewrite the wave equation as

$$s(\mathbf{x}, t) = Ae^{j(\omega t - \mathbf{k}^T \mathbf{x})}$$
$$= Ae^{j\omega(t - \boldsymbol{\alpha}^T \mathbf{x})}$$

where $\alpha = \mathbf{k}/\omega$ is the slowness vector.

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As $c = \omega/||\mathbf{k}||$, vector α has a magnitude which is the reciprocal of c.

Any arbitrary periodic waveform $s(\mathbf{x}, t) = s(t - \boldsymbol{\alpha}^T \mathbf{x})$ with fundamental period ω_0 can be represented as a sum:

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The coefficients are given by

$$S_n = \frac{1}{T} \int_0^T s(u) e^{-jn\omega_0 u} du$$

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- The various components of $s(\mathbf{x}, t)$ have different frequencies $\omega = n\omega_0$ and different wavenumber vectors, **k**.
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We will come back to this later...

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2. Sensor Arrays and Spatial Filtering

- 3. Optimal Beamforming
- 4. Adaptive Beamforming
- 5. DoA Estimation with Microphone Arrays

2. Sensor Arrays and Spatial Filtering

2.1 Wavenumber-Frequency Space

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$$s(\mathbf{x},t) = \frac{1}{(2\pi)^4} \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} S(\mathbf{k},\omega) e^{j(\omega t - \mathbf{k}^T \mathbf{x})} d\mathbf{k} \, d\omega$$

We have already concluded that if the space-time signal is a propagating waveform such that $s(\mathbf{x}, t) = s(t - \boldsymbol{\alpha}_0^T \mathbf{x})$, then its Fourier transform is equal to

$$S(\mathbf{k},\omega) = S(\omega)\delta(\mathbf{k}-\omega\boldsymbol{\alpha}_0)$$

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Remember the nonperiodic propagating wave Fourier transform?

This means that $s(\mathbf{x}, t)$ only has energy along the direction of $\mathbf{k} = \mathbf{k}_0 = \omega \boldsymbol{\alpha}_0$ in the wavenumber-frequency space. 2.2 Frequency-Wavenumber (WN) Response and Beam patterns (BP)

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$$\boldsymbol{f}(t,\boldsymbol{p}) = \begin{bmatrix} f(t,\boldsymbol{p}_0) \\ f(t,\boldsymbol{p}_1) \\ \vdots \\ f(t,\boldsymbol{p}_{N-1}) \end{bmatrix}$$

Array output



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$$y(t) = \sum_{n=0}^{N-1} \int_{-\infty}^{\infty} h_n(t-\tau) f_n(\tau, \boldsymbol{p}_n) d\tau$$
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where $h(t) = [h_o(t) \ h_1(t) \ \cdots \ h_{N-1}(t)]^T$

In the frequency domain, ····

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Plane wave propagating ····

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If f(t) is the signal that would be received at the origin, then:

$$\boldsymbol{f}(t,\boldsymbol{p}) = \begin{bmatrix} f(t-\tau_0) \\ f(t-\tau_1) \\ \vdots \\ f(t-\tau_{N-1}) \end{bmatrix}$$



























$$\therefore \tau_n = -\frac{\mathbf{u}^T \mathbf{p}_n}{c} = \frac{\mathbf{a}^T \mathbf{p}_n}{c}$$



n is the time since the plane wave hits the sensor at location \mathbf{p}_n until it reaches point (0,0).

Back to the frequency domain

Then, we have:

$$\boldsymbol{F}(\omega) = \begin{bmatrix} \int_{-\infty}^{\infty} e^{-j\omega t} f(t-\tau_0) dt \\ \int_{-\infty}^{\infty} e^{-j\omega t} f(t-\tau_1) dt \\ \vdots \\ \int_{-\infty}^{\infty} e^{-j\omega t} f(t-\tau_{N-1}) dt \end{bmatrix}$$
$$= \begin{bmatrix} e^{-j\omega \tau_0} \\ e^{-j\omega \tau_1} \\ \vdots \\ e^{-j\omega \tau_{N-1}} \end{bmatrix} F(\omega)$$

$$\boldsymbol{k} = \frac{\omega}{c} \boldsymbol{a} = \frac{2\pi}{c/f} \boldsymbol{a} = \frac{2\pi}{\lambda} \boldsymbol{a} = -\frac{2\pi}{\lambda} \boldsymbol{u}$$

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Wavenumber Vector ("spatial frequency")



$$\mathbf{k} \neq \frac{\omega}{c} \mathbf{a} = \frac{2\pi}{c/f} \mathbf{a} = \frac{2\pi}{\lambda} \mathbf{a} = -\frac{2\pi}{\lambda} \mathbf{u}$$

Wavenumber Vector ("spatial frequency")

- Note that $|\boldsymbol{k}| = \frac{2\pi}{\lambda}$
- Therefore

$$\omega au_n = rac{\omega}{c} \boldsymbol{a}^T \boldsymbol{p}_n = \boldsymbol{k}^T \boldsymbol{p}_n$$

Array Manifold Vector

And we have

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Array Manifold Vector

• In this particular example, we can use $h_n(t) = \frac{1}{N}\delta(t+\tau_n)$ such that

$$y(t) = f(t)$$

Following, we have the delay-and-sum beamformer.





A common delay is added in each channel to make the operations physically realizable



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Delay-and-sum Beamformer



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Array Manifold Vector

LTI System



LTI System

$$e^{j\omega t} \longrightarrow h(t) \longrightarrow H(\omega)e^{j\omega t}$$

Space-time signals (base functions):

$$f_n(t, \boldsymbol{p}) = e^{j\omega(t-\tau_n)} = e^{j(\omega t - \boldsymbol{k}^T \boldsymbol{p}_n)}$$
Note that $\omega \tau_n = \boldsymbol{k}^T \boldsymbol{p}_n$

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Note that $\omega \tau_n = \boldsymbol{k}^T \boldsymbol{p}_n$

$$f(t, \boldsymbol{p}) = e^{j\omega t} \boldsymbol{v}_{\boldsymbol{k}}(\boldsymbol{k})$$

Frequency-Wavenumber Response Function

The response of the array to this plane wave is:

$$y(t, \boldsymbol{k}) = \boldsymbol{H}^T(\omega) \boldsymbol{v}_{\boldsymbol{k}}(\boldsymbol{k}) e^{j\omega t}$$

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And we define the Frequency-Wavenumber Response Function:

Upsilon
$$(\boldsymbol{\hat{\boldsymbol{\omega}}},\boldsymbol{k}) \triangleq \boldsymbol{H}^{T}(\boldsymbol{\omega})\boldsymbol{v}_{\boldsymbol{k}}(\boldsymbol{k})$$

 $\Upsilon(\omega, \mathbf{k})$ describes the complex gain of an array to an input plane wave with wavenumber \mathbf{k} and temporal frequency ω .

Beam Pattern and Bandpass Signal

BEAM PATTERN is the Frequency Wavenumber Response Function evaluated versus the direction:

$$B(\omega:\theta,\phi) = \Upsilon(\omega, \boldsymbol{k})$$

Note that $\mathbf{k} = \frac{2\pi}{\lambda} \mathbf{a}(\theta, \phi)$, and \mathbf{a} is the unit vector with spherical coordinates angles θ and ϕ

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Let's write a bandpass signal:

$$f(t, \boldsymbol{p}_n) = \sqrt{2}Re\{\tilde{f}(t, \boldsymbol{p}_n)e^{j\omega_c t}\}, n = 0, 1, \cdots, N-1$$

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Let's write a bandpass signal:

$$f(t, \boldsymbol{p}_n) = \sqrt{2Re}\{\tilde{f}(t, \boldsymbol{p}_n)e^{j\omega_c t}\}, n = 0, 1, \cdots, N-1$$

• ω_c corresponds to the carrier frequency and the complex envelope $\tilde{f}(t, \boldsymbol{p}_n)$ is bandlimited to the region $|\underbrace{\omega - \omega_c}_{\omega_L}| \leq 2\pi B_s/2$

Bandlimited plane wave: $f(t, \boldsymbol{p}_n) = \sqrt{2}Re\{\tilde{f}(t - \tau_n)e^{j\omega_c(t - \tau_n)}\}, n = 0, 1, \cdots, N - 1$

- **Bandlimited plane wave:** $f(t, \boldsymbol{p}_n) = \sqrt{2}Re\{\tilde{f}(t \tau_n)e^{j\omega_c(t \tau_n)}\}, n = 0, 1, \cdots, N 1$
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- Assuming the origin is at the array's center of gravity: $\sum_{n=0}^{N-1} \boldsymbol{p}_n = 0 \Rightarrow \tau_n \leq \Delta T_{max}$
- In Narrowband (NB) signals, $B_s \Delta T_{max} \ll 1$ $\Rightarrow \tilde{f}(t - \tau_n) \simeq \tilde{f}(t)$ and $f(t, \mathbf{p}_n) = \sqrt{2Re} \{ \tilde{f}(t) e^{-j\omega_c \tau_n} e^{j\omega_c t} \}$

Phased-Array

For NB signals, the delay is approximated by a phase-shift:

 \Rightarrow delay&sum beamformer \equiv PHASED ARRAY

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The phased array can be implemented adjusting the gain and phase to achieve a desired beam pattern

NB Beamformers

In narrowband beamformers: $y(t, k) = w^H v_k(k) e^{j\omega t}$

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•
$$\Upsilon(\omega, \mathbf{k}) = \underbrace{\mathbf{w}^{H}}_{\mathbf{H}^{T}(\omega)} \mathbf{v}_{\mathbf{k}}(\mathbf{k})$$

2.3 Uniform Linear Arrays (ULA)

Uniformly Spaced Linear Arrays



Microphone-Array Signal Processing. C Apolinárioi & Campos – p. 43/115



• An ULA along axis z:



Location of the elements:

$$\begin{cases} p_{zn} = (n - \frac{N-1}{2})d, \text{ for } n = 0, 1, \cdots, N-1 \\ p_{xn} = p_{yn} = 0 \end{cases}$$



Location of the elements:

$$\begin{cases} p_{zn} = (n - \frac{N-1}{2})d, \text{ for } n = 0, 1, \cdots, N-1 \\ p_{xn} = p_{yn} = 0 \end{cases}$$

• Therefore,
$$p_n = \begin{bmatrix} 0 \\ 0 \\ (n - \frac{N-1}{2})d \end{bmatrix}$$

Array manifold vector:

$$\boldsymbol{v}_{\boldsymbol{k}}(\boldsymbol{k}) = \begin{bmatrix} e^{-j\boldsymbol{k}^{T}\boldsymbol{p}_{0}} \\ \vdots \\ e^{-j\boldsymbol{k}^{T}\boldsymbol{p}_{n}} \\ \vdots \\ e^{-j\boldsymbol{k}^{T}\boldsymbol{p}_{N-1}} \end{bmatrix}_{[k_{x} \ k_{y} \ k_{z}]} \begin{bmatrix} 0 \\ 0 \\ [n-\frac{N-1}{2}]d \end{bmatrix}$$

Array manifold vector:

$$\boldsymbol{v}_{\boldsymbol{k}}(\boldsymbol{k}) = \begin{bmatrix} e^{-j\boldsymbol{k}^{T}\boldsymbol{p}_{0}} & e^{-j\boldsymbol{k}^{T}\boldsymbol{p}_{1}} & \dots & e^{-j\boldsymbol{k}^{T}\boldsymbol{p}_{N-1}} \end{bmatrix}^{T}$$
$$\therefore \boldsymbol{v}_{\boldsymbol{k}}(\boldsymbol{k}) = \boldsymbol{v}_{\boldsymbol{k}}(k_{z}) = \begin{bmatrix} e^{+j\frac{(N-1)}{2}k_{z}d} \\ e^{+j(\frac{N-1}{2}-1)k_{z}d} \\ \vdots \\ e^{-j(\frac{N-1}{2})k_{z}d} \end{bmatrix}$$

Outline

- 1. Introduction and Fundamentals
- 2. Sensor Arrays and Spatial Filtering
- 3. Optimal Beamforming
- 4. Adaptive Beamforming
- 5. DoA Estimation with Microphone Arrays

3. Optimal Beamforming

3.1 Introduction

Scope: use the statistical representation of signal and noise to design array processors that are optimal in a statistical sense.

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- We assume that the appropriate statistics are known.
- Our objective of interest is to estimate the waveform of a plane-wave impinging on the array in the presence of noise and interfering signals.
- Even if a particular beamformer developed in this chapter has good performance, it does not guarantee that its adaptive version (next chapter) will. However, if the performance is poor, it is unlikely that the adaptive version will be useful.

3.2 Optimal Beamformers

MVDR Beamformer

Snapshot model in the frequency domain:
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In many applications, we implement a beamforming in the frequency domain ($\omega_m = \omega_c + m \frac{2\pi}{\Delta T}$ and M varies from $-\frac{M-1}{2}$ to $\frac{M-1}{2}$ if odd and from $-\frac{M}{2}$ to $\frac{M}{2} - 1$ if even).



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In order to generate these vectors, divide the observation interval *T* in *K* disjoint intervals of duration ΔT : $(k-1)\Delta T \leq t < k\Delta T, k = 1, \dots, K$.

• ΔT must be significantly greater than the propagation time across the array.

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- ΔT also depends on the bandwidth of the input signal.
- Assume an input signal with BW B_s centered in f_c
- In order to develop the frequency-domain snapshot model for the case in which the desired signals and the interfering signals can de modeled as plane waves, we have two cases: desired signals are deterministic or samples of a random process.

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- Frequency-domain snapshot consists of signal plus noise: $X(\omega) = X_s(\omega) + N(\omega)$
- The signal vector can be written as $X_s(\omega) = F(\omega)v(\omega : k_s)$ where $F(\omega)$ is the frequency-domain snapshot of the source signal and $v(\omega : k_s)$ is the array manifold vector for a plane-wave with wavenumber k_s .

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- The noise snapshot is a zero-mean random vector $N(\omega)$ with spectral matrix given by $S_n(\omega) = S_c(\omega) + \sigma_{\omega}^2 I$

• We process $X(\omega)$ with the $1 \times N$ operator $W^H(\omega)$:



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Distortionless criterion (in the absence of noise):

 $Y(\omega) = F(\omega)$ = $W^{H}(\omega)X_{s}(\omega) = F(\omega)W^{H}(\omega)v(\omega:k_{s})$ $\Longrightarrow W^{H}(\omega)v(\omega:k_{s}) = 1$

In the presence of noise, we have:

 $Y(\omega) = F(\omega) + Y_n(\omega)$

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The mean square of the output noise is:

$$E[|Y_n(\omega)|^2] = \boldsymbol{W}^H(\omega)\boldsymbol{S}_n(\omega)\boldsymbol{W}(\omega)$$

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 $E[|Y_n(\omega)|^2]$ subject to $W^H(\omega)v(\omega:k_s)=1$

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Using the method of Lagrange multipliers, we define the following cost function to be minimized

$$F = \mathbf{W}^{H}(\omega)\mathbf{S}_{n}(\omega)\mathbf{W}\omega$$
$$+ \lambda \left[\mathbf{W}^{H}(\omega)\mathbf{v}(\omega:\mathbf{k}_{s}) - 1\right] + \lambda^{*} \left[\mathbf{v}^{H}(\omega:\mathbf{k}_{s})\mathbf{W}(\omega) - 1\right]$$

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 \checkmark ...and the result (suppressing ω and k_s) is

$$oldsymbol{W}_{mvdr}^{H}=\Lambda_{s}oldsymbol{v}^{H}oldsymbol{S}_{n}^{-1}$$
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This result is referred to as MVDR or Capon Beamformer.

If the gradient of ξ with respect to w (real case):

$$\boldsymbol{\nabla}_{\boldsymbol{w}} \boldsymbol{\xi} = \begin{bmatrix} \frac{\partial \boldsymbol{\xi}}{\partial w_0} \\ \frac{\partial \boldsymbol{\xi}}{\partial w_1} \\ \vdots \\ \frac{\partial \boldsymbol{\xi}}{\partial w_{N-1}} \end{bmatrix}$$

• The gradient of ξ with respect to w (real case):

$$\nabla w \xi = \begin{bmatrix} \frac{\partial \xi}{\partial w_0} \\ \frac{\partial \xi}{\partial w_1} \\ \vdots \\ \frac{\partial \xi}{\partial w_{N-1}} \end{bmatrix}$$

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which, when R is symmetric, leads to $\nabla w(w^T R w) = 2Rw$

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- which corresponds to $\nabla w \xi = \nabla a \xi + j \nabla_b \xi$
- Let us define the *derivative* $\frac{\partial}{\partial w}$ (with respect to w):

$$\frac{\partial}{\partial \mathbf{w}} = \frac{1}{2} \begin{bmatrix} \frac{\partial}{\partial a_0} - j \frac{\partial}{\partial b_0} \\ \frac{\partial}{\partial a_1} - j \frac{\partial}{\partial b_1} \\ \vdots \\ \frac{\partial}{\partial a_{N-1}} - j \frac{\partial}{\partial b_{N-1}} \end{bmatrix}$$

● The conjugate derivative with respect to w is

$$\frac{\partial}{\partial \mathbf{w}^*} = \frac{1}{2} \begin{bmatrix} \frac{\partial}{\partial a_0} + j \frac{\partial}{\partial b_0} \\ \frac{\partial}{\partial a_1} + j \frac{\partial}{\partial b_1} \\ \vdots \\ \frac{\partial}{\partial a_{N-1}} + j \frac{\partial}{\partial b_{N-1}} \end{bmatrix}$$

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- Therefore, $\nabla_{\mathbf{w}} \xi = \nabla_{\mathbf{a}} \xi + j \nabla_{\mathbf{b}} \xi$ is equivalent to $2 \frac{\partial \xi}{\partial \mathbf{w}^*}$.
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- The complex gradient may be slightly tricky if compared to the simple real gradient. For this reason, we exemplify the use of the complex gradient by calculating $\nabla_{\mathbf{w}} E[|e(k)|^2]$.

•
$$\nabla_{\mathbf{w}} E[e(k)e^{*}(k)] = E\{e^{*}(k)[\nabla_{\mathbf{w}}e(k)] + e(k)[\nabla_{\mathbf{w}}e^{*}(k)]\}$$

We compute each gradient ...

$$\nabla_{\mathbf{w}} e(k) = \nabla_{\mathbf{a}} [d(k) - \mathbf{w}^{H} \mathbf{x}(k)] + j \nabla_{\mathbf{b}} [d(k) - \mathbf{w}^{H} \mathbf{x}(k)]$$
$$= -\mathbf{x}(k) - \mathbf{x}(k) = -2\mathbf{x}(k)$$

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$$\nabla_{\mathbf{w}} e^*(k) = \nabla_{\mathbf{a}} [d^*(k) - \mathbf{w}^T \mathbf{x}^*(k)] + j \nabla_{\mathbf{b}} [d^*(k) - \mathbf{w}^T \mathbf{x}^*(k)]$$
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$$= -\mathbf{x}^*(k) + \mathbf{x}^*(k) = \mathbf{0}$$

such that the final result is

$$\nabla_{\mathbf{w}} E[e(k)e^{*}(k)] = -2E[e^{*}(k)\mathbf{x}(k)]$$
$$= -2E[\mathbf{x}(k)[d(k) - \mathbf{w}^{H}\mathbf{x}(k)]^{*}]$$
$$= -2\underbrace{E[\mathbf{x}(k)d^{*}(k)]}_{\mathbf{p}} + 2\underbrace{E[\mathbf{x}(k)\mathbf{x}^{H}(k)]}_{\mathbf{R}}\mathbf{w}$$

Which results in the Wiener solution $w = R^{-1}p$.

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- When a set of linear constraints involving the coefficient vector of an adaptive filter is imposed, the resulting problem (LCAF)—admitting the MSE as the objective function—can be stated as minimizing $E[|e(k)|^2]$ subject to $C^H w = f$.

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- The output of the processor is $y(k) = \mathbf{w}^H \mathbf{x}(k)$.
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- When a set of linear constraints involving the coefficient vector of an adaptive filter is imposed, the resulting problem (LCAF)—admitting the MSE as the objective function—can be stated as minimizing $E[|e(k)|^2]$ subject to $\mathbf{C}^H \mathbf{w} = \mathbf{f}$.
- The output of the processor is $y(k) = \mathbf{w}^H \mathbf{x}(k)$.
- It is worth mentioning that the most general case corresponds to having a reference signal, d(k). It is, however, usual to have no reference signal as in Linearly-Constrained Minimum-Variance (LCMV) applications. In LCMV, if f = 1, the system is often referred to as Minimum-Variance Distortionless Response (MVDR).

• Using Lagrange multipliers, we form $\xi(k) = E[e(k)e^*(k)] + \mathcal{L}_R^T Re[\mathbf{C}^H \mathbf{w} - \mathbf{f}] + \mathcal{L}_I^T Im[\mathbf{C}^H \mathbf{w} - \mathbf{f}]$

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- We can also represent the above expression with a complex \mathcal{L} given by $\mathcal{L}_R + j\mathcal{L}_I$ such that

$$\xi(k) = E[e(k)e^*(k)] + Re[\mathcal{L}^H(\mathbf{C}^H\mathbf{w} - \mathbf{f})]$$
$$= E[e(k)e^*(k)] + \frac{1}{2}\mathcal{L}^H(\mathbf{C}^H\mathbf{w} - \mathbf{f}) + \frac{1}{2}\mathcal{L}^T(\mathbf{C}^T\mathbf{w}^* - \mathbf{f}^*)$$

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$$= E[e(k)e^*(k)] + \frac{1}{2}\mathcal{L}^H(\mathbf{C}^H\mathbf{w} - \mathbf{f}) + \frac{1}{2}\mathcal{L}^T(\mathbf{C}^T\mathbf{w}^* - \mathbf{f}^*)$$

• Noting that $e(k) = d(k) - \mathbf{w}^H \mathbf{x}(k)$, we compute:

$$\nabla_{\mathbf{w}}\xi(k) = \nabla_{\mathbf{w}}\left\{E[e(k)e^{*}(k)] + \frac{1}{2}\mathcal{L}^{H}(\mathbf{C}^{H}\mathbf{w} - \mathbf{f}) + \frac{1}{2}\mathcal{L}^{T}(\mathbf{C}^{T}\mathbf{w}^{*} - \mathbf{f}^{*})\right\}$$
$$= E[-2\mathbf{x}(k)e^{*}(k)] + \mathbf{0} + \mathbf{C}\mathcal{L}$$
$$= -2E[\mathbf{x}(k)d^{*}(k)] + 2E[\mathbf{x}(k)\mathbf{x}^{H}(k)]\mathbf{w} + \mathbf{C}\mathcal{L}$$

• By using $\mathbf{R} = E[\mathbf{x}(k)\mathbf{x}^{H}(k)]$ and $\mathbf{p} = E[d^{*}(k)\mathbf{x}(k)]$, the gradient is equated to zero and the results can be written as (note that stationarity was assumed for the input and reference signals): $-2\mathbf{p} + 2\mathbf{R}\mathbf{w} + \mathbf{C}\mathcal{L} = \mathbf{0}$

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- Symplex replacing \mathcal{L} , we obtain the Wiener solution for the linearly constrained adaptive filter: $\mathbf{w}_{opt} = \mathbf{R}^{-1}\mathbf{p} + \mathbf{R}^{-1}\mathbf{C}(\mathbf{C}^{H}\mathbf{R}^{-1}\mathbf{C})^{-1}(\mathbf{f} \mathbf{C}^{H}\mathbf{R}^{-1}\mathbf{p})$

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- For this case, d(k) = 0, the cost function is termed minimum output energy (MOE) and is given by $E[|e(k)|^2] = \mathbf{w}^H \mathbf{R} \mathbf{w}$
- Also note that in case we do not have constraints (C and f are nulls), the optimal solution above becomes the unconstrained Wiener solution R⁻¹p.

We start by doing a transformation in the coefficient vector.

 $\textbf{ Let } \mathbf{T} = [\mathbf{C} \ \mathbf{B}] \text{ such that }$

$$\mathbf{w} = \mathbf{T} \bar{\mathbf{w}} = [\mathbf{C} \ \mathbf{B}] \begin{bmatrix} \bar{\mathbf{w}}_U \\ - \bar{\mathbf{w}}_L \end{bmatrix} = \mathbf{C} \bar{\mathbf{w}}_U - \mathbf{B} \bar{\mathbf{w}}_L$$

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- $\bar{\mathbf{w}}_U$ is fixed and termed the quiescent weight vector; the minimization process will be carried out only in the lower part, also designated $\mathbf{w}_{GSC} = \bar{\mathbf{w}}_L$.

It is shown below how to split the transformation matrix into two parts: a fixed path and an *adaptive* path.



This structure (detailed below) was named the Generalized Sidelobe Canceller (GSC).



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- It is always possible to have the overall equivalent coefficient vector which is given by $\mathbf{w} = \mathbf{F} \mathbf{B}\mathbf{w}_{GSC}$.
- If we pre-multiply last equation by \mathbf{B}^H and isolate \mathbf{w}_{GSC} , we find $\mathbf{w}_{GSC} = -(\mathbf{B}^H \mathbf{B})^{-1} \mathbf{B}^H \mathbf{w}$.
- Knowing that $\mathbf{T} = [\mathbf{C} \mathbf{B}]$ and that $\mathbf{T}^H \mathbf{T} = \mathbf{I}$, it follows that $\mathbf{P} = \mathbf{I} \mathbf{C}(\mathbf{C}^H \mathbf{C})^{-1}\mathbf{C}^H = \mathbf{B}(\mathbf{B}^H \mathbf{B})\mathbf{B}^H$.

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- The cross-correlation vector is given as:

$$\mathbf{p}_{GSC} = E[d_{GSC}^* \mathbf{x}_{GSC}]$$

= $E\{[\mathbf{F}^H \mathbf{x} - d]^* [\mathbf{B}^H \mathbf{x}]\}$
= $E[-\mathbf{B}^H d^* \mathbf{x} + \mathbf{B}^H \mathbf{x} \mathbf{x}^H \mathbf{F}]$
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• • • • and $\mathbf{w}_{GSC-OPT} = (\mathbf{B}^H \mathbf{R} \mathbf{B})^{-1} (-\mathbf{B}^H \mathbf{p} + \mathbf{B}^H \mathbf{R} \mathbf{F})$





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- In this case, the optimum filter \mathbf{w}_{OPT} is: $\mathbf{F} - \mathbf{B}\mathbf{w}_{GSC-OPT} = \mathbf{F} - \mathbf{B}(\mathbf{B}^{H}\mathbf{R}\mathbf{B})^{-1}\mathbf{B}^{H}\mathbf{R}\mathbf{F} =$ $\mathbf{R}^{-1}\mathbf{C}(\mathbf{C}^{H}\mathbf{R}^{-1}\mathbf{C})^{-1}\mathbf{f}$ (LCMV solution)

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With simple constraint matrices, simple blocking matrices satisfying $B^TC = 0$ are possible.

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```
[U,S,V]=svd(C);
B3=U(:,p+1:M*N); % p=N in this case
```

The GSC

SVD: the blocking matrix can be produced with the following Matlab command lines,

[U,S,V]=svd(C);

B3=U(:,p+1:M*N); % p=N in this case

\mathbf{B}_3^T is given by:

-0.50	-0.17	-0.17	0.83	0.00	0.00	0.00	0.00	0.00	0.00	0.00	0.00
0.25	-0.42	0.08	0.08	0.75	-0.25	-0.25	-0.25	0.00	0.00	0.00	0.00
0.25	-0.42	0.08	0.08	-0.25	0.75	-0.25	-0.25	0.00	0.00	0.00	0.00
0.25	-0.42	0.08	0.08	-0.25	-0.25	0.75	-0.20	0.00	0.00	0.00	0.00
0.25	-0.42	0.08	0.08	-0.25	-0.25	-0.25	0.75	0.00	0.00	0.00	0.00
0.25	0.08	-0.42	0.08	0.00	0.00	0.00	0.00	0.75	-0.25	-0.25	-0.25
0.25	0.08	-0.42	0.08	0.00	0.00	0.00	0.00	-0.25	0.75	-0.25	-0.25
0.25	0.08	-0.42	0.08	0.00	0.00	0.00	0.00	-0.25	-0.25	0.75	-0.25
0.25	0.08	-0.42	0.08	0.00	0.00	0.00	0.00	-0.25	-0.25	-0.25	0.75



[Q,R]=qr(C); B4=Q(:,p+1:M*N);



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- Two other possibilities are: the one presented in [Tseng Griffiths 88] where a decomposition procedure is introduced in order to offer an effective implementation structure and the other one concerned to a narrowband BF implemented with GSC where B is combined with a wavelet transform [Chu Fang 99].

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- Two other possibilities are: the one presented in [Tseng Griffiths 88] where a decomposition procedure is introduced in order to offer an effective implementation structure and the other one concerned to a narrowband BF implemented with GSC where B is combined with a wavelet transform [Chu Fang 99].
- Finally, a new efficient linearly constrained adaptive scheme which can also be visualized as a GSC structure can be found in [Campos&Werner&Apolinário IEEE-TSP Sept. 2002].

Outline

- 1. Introduction and Fundamentals
- 2. Sensor Arrays and Spatial Filtering
- 3. Optimal Beamforming
- 4. Adaptive Beamforming
- 5. DoA Estimation with Microphone Arrays

4. Adaptive Beamforming

4.1 Introduction

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- Different algorithms may be employed for iteratively approximating the desired solution.
- We will briefly cover a small subset of algorithms for constrained adaptive filters.

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 - For example, if direction of arrival of the signal of interest is known, jammer suppression can take place through spatial filtering without the need of training signal, or in systems with constant-envelope modulation (e.g., M-PSK), a constant-modulus constraint can mitigate multipath propagation effects.

4.2 Constrained FIR Filters

Broadband Array Beamformer



Optimal Constrained MSE Filter

We look for

$$\min_{\mathbf{w}} \xi(k) \quad \text{s.t. } \mathbf{C}^H \mathbf{w} = \mathbf{f},$$

where

- $\xi(k) = E[|e(k)|^2]$
- \blacksquare C is the $MN \times p$ constraint matrix
- **f** is the $p \times 1$ gain vector

Optimal Constrained MSE Filter

The optimal beamformer is

$$\mathbf{w}(k) = \mathbf{R}^{-1}\mathbf{p} + \mathbf{R}^{-1}\mathbf{C}\left(\mathbf{C}^{H}\mathbf{R}^{-1}\mathbf{C}\right)^{-1}\left(\mathbf{f} - \mathbf{C}^{H}\mathbf{R}^{-1}\mathbf{p}\right)$$

where:

•
$$\mathbf{R} = E\left[\mathbf{x}(k)\mathbf{x}^{H}(k)\right] \text{ and } \mathbf{p} = E\left[d^{*}(k)\mathbf{x}(k)\right]$$

• $\mathbf{w}(k) = \left[\mathbf{w}_{1}^{T}(k) \mathbf{w}_{2}^{T}(k) \cdots \mathbf{w}_{M}^{T}(k)\right]^{T}$
• $\mathbf{x}(k) = \left[\mathbf{x}_{1}^{T}(k) \mathbf{x}_{2}^{T}(k) \cdots \mathbf{x}_{M}^{T}(k)\right]^{T}$
• $\mathbf{x}_{i}^{T}(k) = \left[x_{i}(k) x_{i}(k-1) \cdots x_{i}(k-N+1)\right]$

The Constrained LS Beamformer

In the absence of statistical information, we may choose

$$\min_{\mathbf{w}} \left[\xi(k) = \sum_{i=0}^{k} \lambda^{k-i} |d(i) - \mathbf{w}^H \mathbf{x}(i)|^2 \right] \text{ s.t. } \mathbf{C}^H \mathbf{w} = \mathbf{f}$$

with $\lambda \in (0, 1]$,

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with $\lambda \in (0, 1]$, which gives, as solution,

$$\mathbf{w}(k) = \mathbf{R}^{-1}(k)\mathbf{p}(k) + \mathbf{R}^{-1}(k)\mathbf{C} \left(\mathbf{C}^{H}\mathbf{R}^{-1}(k)\mathbf{C}\right)^{-1} \left[\mathbf{f} - \mathbf{C}^{H}\mathbf{R}^{-1}(k)\mathbf{p}(k)\right],$$

where

$$\mathbf{R}(k) = \sum_{i=0}^{k} \lambda^{k-i} \mathbf{x}(i) \mathbf{x}^{H}(i), \text{ and } \mathbf{p}(k) = \sum_{i=0}^{k} \lambda^{k-i} d^{*}(i) \mathbf{x}(i).$$

A (cheaper) alternative cost function is

$$\min_{\mathbf{w}} \left[\xi(k) = \| \mathbf{w}(k) - \mathbf{w}(k-1) \|^2 + \mu |e(k)|^2 \right] \text{ s.t. } \mathbf{C}^H \mathbf{w}(k) = \mathbf{f},$$

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which gives, as solution,

$$\mathbf{w}(k) = \mathbf{w}(k-1) + \mu e^*(k) \left[\mathbf{I} - \mathbf{C} \left(\mathbf{C}^H \mathbf{C} \right)^{-1} \mathbf{C}^H \right] \mathbf{x}(k),$$

where $e(k) = d(k) - \mathbf{w}^{H}(k-1)\mathbf{x}(k)$, μ is a positive small constant called step size.

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which gives, as solution,

$$\mathbf{w}(k) = \mathbf{P}\left[\mathbf{w}(k-1) + \mu e^*(k)\mathbf{x}(k)\right] + \mathbf{F},$$

where $e(k) = d(k) - \mathbf{w}^{H}(k-1)\mathbf{x}(k)$, μ is a positive small constant called step size, $\mathbf{P} = \mathbf{C} (\mathbf{C}^{H}\mathbf{C})^{-1} \mathbf{C}^{H}$, and $\mathbf{F} = \mathbf{C} (\mathbf{C}^{H}\mathbf{C})^{-1} \mathbf{f}$.

We may wish to trade complexity for speed of convergence:

$$\min_{\mathbf{w}} \left[\xi(k) = \|\mathbf{w}(k) - \mathbf{w}(k-1)\|^2 \right] \text{ s.t. } \begin{cases} \mathbf{X}^T(k)\mathbf{w}^*(k) = \mathbf{d}(k) \\ \mathbf{C}^H \mathbf{w}(k) = \mathbf{f}, \end{cases}$$

where

•
$$\mathbf{d}(k) = [d(k) \ d(k-1) \ \cdots \ d(k-L+1)]^T$$

• $\mathbf{X}(k) = [\mathbf{x}(k) \ \mathbf{x}(k-1) \ \cdots \ \mathbf{x}(k-L+1)]^T$

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which gives, as solution,

$$\mathbf{w}(k) = \mathbf{P}\left[\mathbf{w}(k-1) + \mu \mathbf{X}(k)\mathbf{t}(k)\right] + \mathbf{F}$$

where

•
$$\mathbf{e}(k) = \mathbf{d}(k) - \mathbf{X}^T(k)\mathbf{w}^*(k-1)$$

• $\mathbf{t}(k) = \left[\mathbf{X}^H(k)\mathbf{P}\mathbf{X}(k)\right]^{-1}\mathbf{e}^*(k)$

Outline

- 1. Introduction and Fundamentals
- 2. Sensor Arrays and Spatial Filtering
- 3. Optimal Beamforming
- 4. Adaptive Beamforming
- 5. DoA Estimation with Microphone Arrays

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5.0 Signal Preparation

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- But, most importantly, the delay is well represented only if the signal is also analytic, i. e., having only non-negative frequency components.

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- First thing to note: when this is the case, the signal is narrow band with a center frequency in ω_0 (in the continuous-time domain, it corresponds to a carrier frequency $\Omega_0 = f_s \omega_0$)
- But, most importantly, the delay is well represented only if the signal is also analytic, i. e., having only non-negative frequency components.
- An analytic signal, mathematically, can be obtained by multiplying its Fourier transform by the continuous Heaviside step function:

$$X_a(e^{j\omega}) = 2X(e^{j\omega})u(\omega), u(\omega) = \begin{cases} 0, \omega < 0\\ 1, \omega = 0\\ 1, \omega > 0 \end{cases}$$

■ Let $x(n) = s(n) \cos(\omega_0 n)$, s(n) having a maximum frequency component (ω_m) much lower than ω_0 :



If
$$x(n) = s(n)e^{j\omega_0 n}$$
, then
 $x(n)e^{-j\omega_0 \tau} = s(n)e^{j\omega_0(n-\tau)} \approx x(n-\tau)$ if $\tau \ll 1/\omega_m$

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• We can make $x(n) = s(n)cos(\omega_0 n) = \underbrace{\frac{s(n)}{2}e^{j\omega_0 n}}_{x_+(n)} + \underbrace{\frac{s(n)}{2}e^{-j\omega_0 n}}_{x_-(n)}$ such that $x(n-\tau) \approx x_+e^{-j\omega_0 \tau} + x_-(n)e^{+j\omega_0 \tau} = s(n)cos(\omega_0(n-\tau))$

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- ... but, how to obtain $x_{+}(n)$ or a scaled copy? Using the Hilbert Transform $x_{H}(n) = \mathcal{HT}\{x(n)\}$ where $\int jX(e^{j\omega}), -\pi < \omega < 0$

$$X_H(e^{j\omega}) = \begin{cases} X(e^{j\omega}), \omega = 0\\ -jX(e^{j\omega}), 0 < \omega < \pi \end{cases}$$

Knowing that

$$x(n) = x_{-}(n) + x_{+}(n) = \mathcal{F}^{-1} \{X_{-}(e^{j\omega}) + X_{+}(e^{j\omega})\}, \text{ we compute } y(n) = x(n) + jx_{H}(n)$$

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$$y(n) =$$

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$$= \mathcal{F}^{-1} \{ X_{-}(e^{j\omega}) + X_{+}(e^{j\omega}) - X_{-}(e^{j\omega}) + X_{+}(e^{j\omega}) \}$$

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$$= \mathcal{F}^{-1} \{ X_{-}(e^{j\omega}) + X_{+}(e^{j\omega}) - X_{-}(e^{j\omega}) + X_{+}(e^{j\omega}) \}$$

• Therefore $y(n) = 2\mathcal{F}^{-1}\{X_+(e^{j\omega})\} = s(n)e^{j\omega_0 n}$ which is analytic!

Signal Model

• Consider $x_m(t)$ the signal from the *m*-th microphone (prior to the A/D converter) corresponding to audio from *D* sources (directions θ_1 to θ_D) plus noise: $x_m(t) = s_1(t - \bar{\tau}_m(\theta_1)) + \dots + s_D(t - \bar{\tau}_m(\theta_D)) + n_m(t)$

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- Assuming $\bar{\tau}_m(\theta_d) = T\tau_m(\theta_d)$ in s ($\tau_m(\theta_d)$ in number of samples), after the A/D converter and $\{.\} + j\mathcal{HT}\{.\}$ to make it an analytic signal, we could write $x_m(n) = s_1(n)e^{-j\omega_0\tau_m(\theta_1)} + \cdots + s_D(n)e^{-j\omega_0\tau_m(\theta_D)} + n_m(n)$

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- Assuming $\bar{\tau}_m(\theta_d) = T\tau_m(\theta_d)$ in s ($\tau_m(\theta_d)$ in number of samples), after the A/D converter and $\{.\} + j\mathcal{HT}\{.\}$ to make it an analytic signal, we could write $x_m(n) = s_1(n)e^{-j\omega_0\tau_m(\theta_1)} + \cdots + s_D(n)e^{-j\omega_0\tau_m(\theta_D)} + n_m(n)$
- **•** For an array with M microphones, we would have:

$$\underbrace{\mathbf{x}(n)}_{M \times 1} = \underbrace{\mathbf{A}}_{M \times D} \underbrace{\mathbf{s}(n)}_{D \times 1} + \underbrace{\mathbf{n}(n)}_{M \times 1}$$

5.1 Signal model

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Such that the output signal can be written as $y(n) = \mathbf{h}^H \mathbf{x}(n) = \mathbf{h}^H [\mathbf{As}(n) + \mathbf{n}(n)]$ If we now assume one single signal, s(n), coming from direction θ , then $\mathbf{x}(n) = s(n)\mathbf{a}(\theta) + \mathbf{n}(n)$

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- If we make $\mathbf{h}^{H}\mathbf{a}(\theta) = 1$, the output signal would correspond to $y(n) = s(n) + \underbrace{\mathbf{h}^{H}\mathbf{n}(n)}_{\mathbf{u}}$

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- If we make $\mathbf{h}^{H}\mathbf{a}(\theta) = 1$, the output signal would correspond to $y(n) = s(n) + \underbrace{\mathbf{h}^{H}\mathbf{n}(n)}_{\text{noise}}$
- Also note that $E[|y(n)|^2] = \mathbf{h}^H \mathbf{R}_x \mathbf{h}$, $\mathbf{R}_x = E[\mathbf{x}(n)\mathbf{x}^H(n)]$

5.2 Non-parametric methods: BF (beamforming a.k.a. Delay & Sum) and Capon

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- Minimizing $E[|y(n)|^2] = \mathbf{h}^H \mathbf{h}$ s.t. $\mathbf{h}^H \mathbf{a}(\theta) = 1$, the result,

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 obtain E[|y(n)|²] = h^Hh
- Minimizing $E[|y(n)|^2] = \mathbf{h}^H \mathbf{h}$ s.t. $\mathbf{h}^H \mathbf{a}(\theta) = 1$, the result, after using Lagrange multiplier, taking the gradient, and equating to zero, is $\mathbf{h} = \mathbf{a}(\theta)/M$ which leads to $E[|y(n)|^2] = \frac{\mathbf{a}^H(\theta)\mathbf{R}_x\mathbf{a}(\theta)}{M^2}$

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- Omitting factor $\frac{1}{M^2}$, we estimate the autocorrelation matrix as $\hat{\mathbf{R}}_x = \frac{1}{N} \sum_{n=1}^{N} \mathbf{x}(n) \mathbf{x}^H(n)$ and find the direction of interest by varying θ and obtaining the

eak in
$$P_{DS}(\theta) = \mathbf{a}^{H}(\theta) \hat{\mathbf{R}}_{x} \mathbf{a}(\theta)$$

р



In the method known as Capon, we minimize $E[|y(n)|^2] = \mathbf{h}^H \mathbf{R}_x \mathbf{h}$ subject to $\mathbf{h}^H \mathbf{a}(\theta) = 1$



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- Using Lagrange multiplier, we write $\xi = \mathbf{h}^H \mathbf{R}_x \mathbf{h} + \lambda (\mathbf{h}^H \mathbf{a}(\theta) - 1)$, and make $\nabla_{\mathbf{h}} \xi = \mathbf{0}$ such that $\mathbf{h} = \frac{\mathbf{R}_x^{-1} \mathbf{a}(\theta)}{\mathbf{a}^H(\theta) \mathbf{R}_x^{-1} \mathbf{a}(\theta)}$



- In the method known as Capon, we minimize $E[|y(n)|^2] = \mathbf{h}^H \mathbf{R}_x \mathbf{h}$ subject to $\mathbf{h}^H \mathbf{a}(\theta) = 1$
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- Replacing the above coefficient vector in $E[|y(n)|^2]$, we obtain $E[|y(n)|^2] = \frac{1}{\mathbf{a}^H(\theta)\mathbf{R}_x^{-1}\mathbf{a}(\theta)}$



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- Using Lagrange multiplier, we write $\xi = \mathbf{h}^H \mathbf{R}_x \mathbf{h} + \lambda (\mathbf{h}^H \mathbf{a}(\theta) - 1)$, and make $\nabla_{\mathbf{h}} \xi = \mathbf{0}$ such that $\mathbf{h} = \frac{\mathbf{R}_x^{-1} \mathbf{a}(\theta)}{\mathbf{a}^H(\theta) \mathbf{R}_x^{-1} \mathbf{a}(\theta)}$
- Replacing the above coefficient vector in $E[|y(n)|^2]$, we obtain $E[|y(n)|^2] = \frac{1}{\mathbf{a}^H(\theta)\mathbf{R}_x^{-1}\mathbf{a}(\theta)}$
- Therefore, in the Capon DoA, we estimate $\hat{\mathbf{R}}_x = \frac{1}{N} \sum_{n=1}^{N} \mathbf{x}(n) \mathbf{x}^H(n)$ and find the direction of interest by varying θ and obtaining the peak in

$$P_{CAPON}(\theta) = \frac{1}{\mathbf{a}^{H}(\theta)\hat{\mathbf{R}}_{x}^{-1}\mathbf{a}(\theta)}$$

5.3 Eigenvalue-Based DoA

MUSIC

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- Also note that A is $M \times D$, s is $D \times 1$, and $\mathbf{n}(n)$ is $M \times 1$
- We then write $\mathbf{R}_x = E\left[\mathbf{x}(n)\mathbf{x}^H(n)\right] = \mathbf{A}\mathbf{R}_s\mathbf{A}^H + \mathbf{R}_n$, this last matrix becoming $\mathbf{R}_n = \sigma_n^2\mathbf{I}$ when assuming spatially white noise; \mathbf{R}_s is the $D \times D$ autocorrelation matrix of the signal vector, i.e., $E\left[\mathbf{s}(n)\mathbf{s}^H(n)\right]$
• $\mathbf{R}_x = \mathbf{A}\mathbf{R}_s\mathbf{A}^H + \mathbf{R}_n$ with D < M implies that $\mathbf{A}\mathbf{R}_s\mathbf{A}^H$ is singular (rank D), its determinant is equal to zero and, therefore, det $[\mathbf{R}_x - \sigma_n^2 \mathbf{I}] = 0$ and σ_n^2 is a (minimum) eigenvalue with multiplicity M - D

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- Spectral decomposition of matrix R_x: vector e_m being an eigenvector of R_x means that R_xe_m = λ_me_m. Collecting all eigenvectors in matrix E, we may write R_xE = EΛ = [e₁ ··· e_M] diag {[λ₁ ··· λ_M]}
 ⇒ R_x = EΛE^H

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 ⇒ R_x = EΛE^H
- Dividing matrix E in two parts, the first D columns and the last N = M - D columns, we have: $E = [\underbrace{e_1 \cdots e_D}_{E_S} \underbrace{e_{D+1} \cdots e_M}_{E_N}] = [E_S \ E_N]$

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- The columns of \mathbf{E}_S span the *D*-dimensional signal subspace while the columns of \mathbf{E}_N span the *N*-dimensional noise subspace
- A vector in the signal subspace is a linear combination of the columns of \mathbf{E}_S . An example: $\sum_{d=1}^{D} x_d \mathbf{e}_d = \mathbf{E}_S \mathbf{x}, \mathbf{x} = [x_1 \cdots x_D]^T$

- Noting that $\mathbf{E}\mathbf{E}^H = \mathbf{I}$, we can write $\mathbf{E}_S\mathbf{E}_S^H + \mathbf{E}_N\mathbf{E}_N^H = \mathbf{I}$
- The columns of E_S span the D-dimensional signal subspace while the columns of E_N span the N-dimensional noise subspace
- A vector in the signal subspace is a linear combination of the columns of \mathbf{E}_S . An example: $\sum_{d=1}^{D} x_d \mathbf{e}_d = \mathbf{E}_S \mathbf{x}, \mathbf{x} = [x_1 \cdots x_D]^T$
- We can find the distance d from a vector v to the signal subspace \mathbf{E}_S by obtaining x that minimizes $d = |\mathbf{v} \mathbf{E}_S \mathbf{x}|$; the result is $d^2 = \mathbf{v}^H \mathbf{E}_N \mathbf{E}_N^H \mathbf{v}$

• The squared distance from vector $\mathbf{a}(\theta)$ to the signal subspace (spanned by \mathbf{E}_S) is $d^2 = \mathbf{a}^H(\theta)\mathbf{E}_N\mathbf{E}_N^H\mathbf{a}(\theta)$

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- When θ belongs to $\{\theta_1 \cdots \theta_D\}$, this distance should be close to zero
- Its inverse will present peaks. In algorithm MUSIC, we estimate D from the eigenvalues of $\hat{\mathbf{R}}_x$; from its eigenvectors, we form \mathbf{E}_S and \mathbf{E}_N , and by varying θ , we shall find peaks in the directions of θ_1 to θ_D in

$$P_{MUSIC}(\theta) = \frac{1}{d_{\mathbf{a}(\theta)}^2} = \frac{1}{\mathbf{a}^H(\theta)\mathbf{E}_N\mathbf{E}_N^H\mathbf{a}(\theta)}$$

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$$P_{MUSIC}(\theta) = \frac{1}{d_{\mathbf{a}(\theta)}^2} = \frac{1}{\mathbf{a}^H(\theta)\mathbf{E}_N\mathbf{E}_N^H\mathbf{a}(\theta)}$$

• If \mathbf{R}_{S} is required, we compute $\mathbf{R}_{S} = (\mathbf{A}^{H}\mathbf{A})^{-1}\mathbf{A}^{H}(\mathbf{R}_{x} - \sigma_{n}^{2}\mathbf{I})\mathbf{A}(\mathbf{A}^{H}\mathbf{A})^{-1}$

5.4 GCC-Based DoA

M microphones of an array are in positions p_1 to p_M : 6



4

- \mathcal{Z} Mic 1 positioned at \mathbf{p}_1 6 u θ $\mathbf{2}$ 5 \mathcal{Y} M = 7 \mathcal{X} 3
- \blacksquare M microphones of an array are in positions p_1 to p_M :
- \mathbf{P} -u: unit vector in the direction of propagation



- Mic 1 positioned at \mathbf{p}_1 6 u 4 Ĥ 25 \mathcal{Y} M = 7 \mathcal{X} 3
- M microphones of an array are in positions p_1 to p_M :
- $-\mathbf{u}$: unit vector in the direction of propagation
- \bullet θ : grazing angle $(\frac{\pi}{2}$ - elevation angle)
- ϕ : horizontal angle (azimuth)



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•
$$r_{x_m x_l}(\tau) = E[x_m(n)x_l(n-\tau)]$$

• When the sound frontwave first hits microphone m ($\tau_{ml} < 0$):



• When the sound frontwave first hits microphone m ($\tau_{ml} < 0$):



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GCC

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 $\hat{R}_{x_m x_l}(e^{j\omega}) \approx |S(e^{j\omega})|^2 H_m(e^{j\omega}) H_l^*(e^{j\omega}) \text{ and}$ $\hat{r}_{x_m x_l}(\tau) \approx \frac{1}{2\pi} \int_{-\pi}^{\pi} H_m(e^{j\omega}) H_l^*(e^{j\omega}) \hat{R}_s(e^{j\omega}) e^{j\omega\tau} d\omega$

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- Which motivates the GCC:

$$r_{x_m x_l}^G(\tau) = \frac{1}{2\pi} \int_{-\pi}^{\pi} \psi(\omega) \hat{R}_{x_m x_l}(e^{j\omega}) e^{j\omega\tau} d\omega$$

Types of $\psi(\omega)$



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$$\hat{R}_{x_m}(e^{j\omega}) = |X_m(e^{j\omega})|^2$$

$$\hat{R}_{n_m}(e^{j\omega}) = |N_m(e^{j\omega})|^2$$
(estimated during silence interval)

$$\hat{R}_{n_l}(e^{j\omega}) = |N_l(e^{j\omega})|^2$$
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GCC

• PHAT (Phase Transform): $\psi(\omega) = \frac{1}{|\hat{R}_{x_m x_l}(e^{j\omega})|}$

• Replacing this function in the expression of $r_{x_m x_l}^G(\tau)$: $r_{x_m x_l}^{PHAT}(\tau) = \frac{1}{2\pi} \int_{-\pi}^{\pi} \frac{\hat{R}_{x_m x_l}(e^{j\omega})}{|\hat{R}_{x_m x_l}(e^{j\omega})|} e^{j\omega\tau} d\omega$ in which, after making $\hat{R}_{x_m x_l}(e^{j\omega}) = |S(e^{j\omega})|^2 H_m(e^{j\omega}) H_l^*(e^{j\omega})$, we have $r_{x_m x_l}^{PHAT}(\tau) = \frac{1}{2\pi} \int_{-\pi}^{\pi} e^{j(\triangleleft H_m - \triangleleft H_l + \omega\pi)} d\omega$

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- For the PHAT, in case of having $h_m(n) = \alpha_m \delta(n)$ and $h_l(n) = \alpha_l \delta(n - \Delta \tau)$, the cross-correlation would be $r_{x_m x_l}^{PHAT}(\tau) = \delta(\tau + \Delta \tau) \Rightarrow$ peak in $\tau_{ml} = -\Delta \tau$ (a perfect indication of a temporal delay!)

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$$\xi = \left(\bar{\tau}_{12} - \Delta \bar{\mathbf{p}}_{12}^T \mathbf{u}\right)^2 + \dots + \left(\bar{\tau}_{(M-1)M} - \Delta \bar{\mathbf{p}}_{(M-1)M}^T \mathbf{u}\right)^2$$

with $\bar{\tau}_{ml} = \tau_{ml}/f_s$ and $\Delta \bar{\mathbf{p}}_{ml} = (\mathbf{p}_m - \mathbf{p}_l)/v_{sound}$

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- We then find u that minimizes ξ by making $\nabla_{\mathbf{u}}\xi = \mathbf{0}$:

Au = b

where $\mathbf{A} = \Delta \bar{\mathbf{p}}_{12} \Delta \bar{\mathbf{p}}_{12}^T + \dots + \Delta \bar{\mathbf{p}}_{(M-1)M} \Delta \bar{\mathbf{p}}_{(M-1)M}^T$ and $\mathbf{b} = \bar{\tau}_{12} \Delta \bar{\mathbf{p}}_{12}^T + \dots + \bar{\tau}_{(M-1)M} \Delta \bar{\mathbf{p}}_{(M-1)M}$

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• And this unit vector is given as $\mathbf{u} = \begin{bmatrix} u_x \\ u_y \\ u_z \end{bmatrix} = \mathbf{A}^{-1}\mathbf{b}$

Azimuth and elevation

Moving u and also the fact that it corresponds to

$$\begin{bmatrix} \sin\theta\cos\phi\\ \sin\theta\sin\phi\\ \cos\theta \end{bmatrix}, \cdots$$

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And the elevation:

elevation =
$$90^{\circ} - \theta = 90^{\circ} - \arccos u_z$$

Last slide ©

Thank you!