



IPSA TOULOUSE

Signal processing for Radar array

Project: Study and implementation of beamforming and DoA estimation techniques

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Introduction

In this report, we consider a network of N ULA antennas that receives P sources with directions $\boldsymbol{\theta}_S = [\theta_1, \dots, \theta_P]^\top \in \mathbb{R}^P$. First, we will implement the beamforming to identify the source of interest, and in the second part, estimate the directions of arrival of every source by using the DoA estimation methods.

1 Beamforming techniques

In this first section, the objective is to implement beamforming techniques by considering a specific scenario. For that, we consider that only one source, at each instant $k \in \llbracket 1, K \rrbracket$, $s(k) \in \mathbb{C}$ with power P_S and direction, $\mathbf{a}(\theta_S)$ is received by the ULA and corrupted by an interference $i(k) \in \mathbb{C}$, with power P_I and direction $\mathbf{a}_I(\theta_I)$. The signal received by the ULA can be expressed through the expression

$$\mathbf{y}(k) = \mathbf{a}(\theta_S) s(k) + \mathbf{y}_I(k) + \mathbf{n}(k) \quad \mathbf{n}(k) \sim \mathcal{CN}(0, \sigma^2 \mathbf{I}) \quad (1)$$

with $\mathbf{y}(k) \in \mathbb{C}^N$

$$s(k) \sim \mathcal{N}(0, P_S) \quad \mathbf{y}_I(k) = \mathbf{a}_I(\theta_I) i(k) \quad i(k) \sim \mathcal{N}(0, P_I) \quad (2)$$

We consider a ULA, reception is done over $\theta \in [0, 2\pi]$ and the expression of a source direction for such a network is :

$$\mathbf{a}(\theta) = \begin{pmatrix} 1 \\ \vdots \\ \exp^{-2\pi j(n-1)\frac{d}{\lambda} \sin(\theta)} \\ \vdots \\ \exp^{-2\pi j(N-1)\frac{d}{\lambda} \sin(\theta)} \end{pmatrix} \quad (3)$$

The study is carried out by considering two cases depending on the knowledge of the interference-noise covariance matrix \mathbf{C} and of the signal-interference-noise covariance matrix \mathbf{R} .

1.1 Known covariance matrices

In the first case, the covariance matrices involved in the problem are assumed to be completely known. Consequently, the matrices \mathbf{C} and \mathbf{R} can be computed theoretically and implemented without any approximation.

1.1.1 Covariance matrices \mathbf{R} and \mathbf{C} expression

By definition, we have the interference-noise covariance matrix \mathbf{C} and the signal-interference-noise covariance matrix \mathbf{R} :

$$\begin{aligned} \mathbf{C} &= \sigma^2 \mathbf{I} + P_I \mathbf{a}_I(\theta_I) \mathbf{a}_I(\theta_I)^\dagger = \mathbb{E} \left((\mathbf{y}_I(k) + \mathbf{n}(k)) (\mathbf{y}_I(k) + \mathbf{n}(k))^\dagger \right) \\ \mathbf{R} &= \mathbf{C} + P_S \mathbf{a}(\theta_S) \mathbf{a}(\theta_S)^\dagger = \mathbb{E} \left((\mathbf{a}(\theta_S) s(k) + \mathbf{y}_I(k) + \mathbf{n}(k)) (\mathbf{a}(\theta_S) s(k) + \mathbf{y}_I(k) + \mathbf{n}(k))^\dagger \right) \end{aligned}$$

First for \mathbf{C} :

$$\begin{aligned} \mathbf{C} &= \mathbb{E} \left((\mathbf{y}_I(k) + \mathbf{n}(k)) (\mathbf{y}_I(k) + \mathbf{n}(k))^\dagger \right) \\ &= \mathbb{E} \left(\mathbf{y}_I(k) \mathbf{y}_I(k)^\dagger + \mathbf{y}_I(k) \mathbf{n}(k)^\dagger + \mathbf{n}(k) \mathbf{y}_I(k)^\dagger + \mathbf{n}(k) \mathbf{n}(k)^\dagger \right) \\ &= \mathbb{E} \left(\mathbf{y}_I(k) \mathbf{y}_I(k)^\dagger \right) + \mathbb{E} \left(\mathbf{y}_I(k) \mathbf{n}(k)^\dagger \right) + \mathbb{E} \left(\mathbf{n}(k) \mathbf{y}_I(k)^\dagger \right) + \mathbb{E} \left(\mathbf{n}(k) \mathbf{n}(k)^\dagger \right) \end{aligned} \quad (4)$$

By assuming independence between the noise and the interference, the crossed terms are 0 :

$$\mathbb{E} \left(\mathbf{y}_I(k) \mathbf{n}(k)^\dagger \right) = \mathbb{E} \left(\mathbf{n}(k) \mathbf{y}_I(k)^\dagger \right) = 0$$

We also know from Equation 1 that :

$$\mathbb{E} \left(\mathbf{y}_I(k) \mathbf{y}_I(k)^\dagger \right) = \mathbb{E} \left(\left(\mathbf{a}_I(\theta_I) i(k) \right) \left(\mathbf{a}_I(\theta_I) i(k) \right)^\dagger \right) \quad i(k) \sim \mathcal{N}(0, P_I)$$

As $\mathbf{a}_I(\theta_I)$ is not random and $\mathbb{E} \left(i(k) i(k)^\dagger \right) \in \mathbb{R}$:

$$\begin{aligned} \mathbb{E} \left(\mathbf{y}_I(k) \mathbf{y}_I(k)^\dagger \right) &= \mathbf{a}_I(\theta_I) \mathbb{E} \left(i(k) i(k)^\dagger \right) \mathbf{a}_I(\theta_I)^\dagger & i(k) \sim \mathcal{N}(0, P_I) \\ &= \mathbb{E} \left(i(k) i(k)^\dagger \right) \mathbf{a}_I(\theta_I) \mathbf{a}_I(\theta_I)^\dagger & i(k) \sim \mathcal{N}(0, P_I) \\ &= P_I \mathbf{a}_I(\theta_I) \mathbf{a}_I(\theta_I)^\dagger \end{aligned}$$

From Equation 1 directly with $\mathbf{n}(k) \sim \mathcal{CN}(0, \sigma^2 \mathbf{I})$: $\mathbb{E} \left(\mathbf{n}(k) \mathbf{n}(k)^\dagger \right) = \sigma^2 \mathbf{I}$ Replacing terms in Equation 4 :

$$\mathbf{C} = \sigma^2 \mathbf{I} + P_I \mathbf{a}_I(\theta_I) \mathbf{a}_I(\theta_I)^\dagger$$

Likewise for \mathbf{R} , we identify the same terms from \mathbf{C} plus the additional ones :

$$\mathbf{R} = \mathbb{E} \left(\left(\mathbf{a}(\theta_S) s(k) + \mathbf{y}_I(k) + \mathbf{n}(k) \right) \left(\mathbf{a}(\theta_S) s(k) + \mathbf{y}_I(k) + \mathbf{n}(k) \right)^\dagger \right)$$

$$\begin{aligned} \mathbf{R} &= \mathbf{C} + \mathbb{E} \left(\left(\mathbf{a}(\theta_S) s(k) \right) \left(\mathbf{a}(\theta_S) s(k) \right)^\dagger \right) + \mathbb{E} \left(\left(\mathbf{a}(\theta_S) s(k) \right) \left(\mathbf{y}_I(k) + \mathbf{n}(k) \right)^\dagger \right) + \\ &\quad \mathbb{E} \left(\left(\mathbf{y}_I(k) + \mathbf{n}(k) \right) \left(\mathbf{a}(\theta_S) s(k) \right)^\dagger \right) \end{aligned}$$

By assuming independence between the noise and the signal and independence between the interference and the signal, crossed terms are 0 and with $\mathbf{a}(\theta_S)$ not random and $\mathbb{E} \left(s(k) s(k)^\dagger \right) \in \mathbb{R}$:

$$\begin{aligned} \mathbf{R} &= \mathbf{C} + \mathbb{E} \left(s(k) s(k)^\dagger \right) \mathbf{a}(\theta_S) \mathbf{a}(\theta_S)^\dagger \quad s(k) \sim \mathcal{N}(0, P_S) \\ &= \mathbf{C} + P_S \mathbf{a}(\theta_S) \mathbf{a}(\theta_S)^\dagger \end{aligned}$$

With these equations, we can implement the beamforming techniques. We can note that beamforming only depends on the antenna network geometry, including the positions of antennas, wavelength, and DoA, so we do not need to implement the received signal $\mathbf{y}(t)$ for these techniques.

The goal of this part is to compare the different methods. We start by implementing the conventional beamforming \mathbf{w}_{CBF} and verifying the relation between the Signal-Interference-Noise ratios ($SINR_{out} = N SINR_{in}$). Once it is done, we implement the optimal adaptive beamforming \mathbf{w}_{opt} and its $SINR$. We obtain :

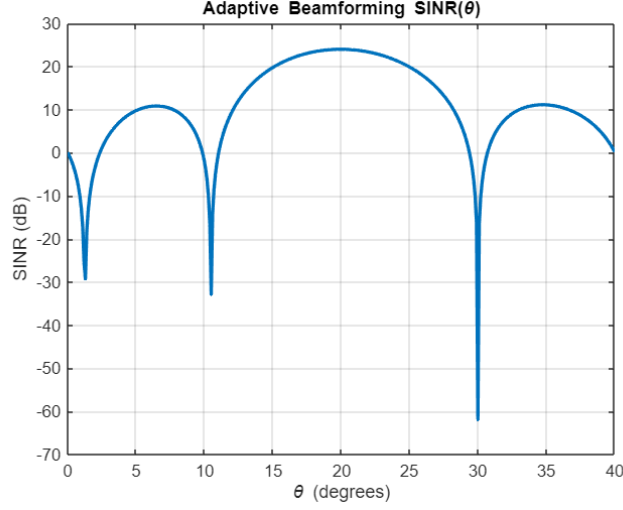


Figure 1: Adaptive Beamforming

Now, we assume that the angle θ_S is not very well-known and is provided by an estimator $\widehat{\theta}_S$. To handle this, we need to implement both the Minimum Power Distortionless Response (*MVDR*) filter \mathbf{w}_{MVDR} and the Minimum Variance Distortionless Response (*MPDR*) filter \mathbf{w}_{MPDR} , considering that only the matrix \mathbf{R} is known.

1.1.2 Signal-to-Interference and Noise Ratio (SINR) expression

We filter the signal at reception with the filter \mathbf{w} . Meaning the received signal is:

$$\mathbf{y}_R(k) = \mathbf{w}^\dagger \mathbf{y}(k) = \underbrace{\mathbf{w}^\dagger \mathbf{a}(\theta) s(k)}_{\text{signal of interest}} + \underbrace{\mathbf{w}^\dagger (\mathbf{y}_I(k) + \mathbf{n}(k))}_{\text{interf. and noise}}$$

The Signal to Interference and Noise Ratio (SINR) is, by definition, a ratio of the *useful* power part (associated with the signal) over the *lost* part (associated with the noise and the interference). At reception, it can be expressed as :

$$\begin{aligned} SINR_{out}(\mathbf{w}) &= \frac{\mathbb{E}(|\mathbf{w}^\dagger \mathbf{a}(\theta) s(k)|^2)}{\mathbb{E}(|\mathbf{w}^\dagger (\mathbf{y}_I(k) + \mathbf{n}(k))|^2)} \\ &= \frac{\mathbb{E}\left((\mathbf{w}^\dagger \mathbf{a}(\theta) s(k)) (\mathbf{w}^\dagger \mathbf{a}(\theta) s(k))^\dagger\right)}{\mathbb{E}\left((\mathbf{w}^\dagger (\mathbf{y}_I(k) + \mathbf{n}(k))) (\mathbf{w}^\dagger (\mathbf{y}_I(k) + \mathbf{n}(k)))^\dagger\right)} \\ &= \frac{\mathbb{E}\left(\mathbf{w}^\dagger \mathbf{a}(\theta) s(k) s(k)^\dagger \mathbf{a}(\theta)^\dagger \mathbf{w}\right)}{\mathbb{E}\left(\mathbf{w}^\dagger (\mathbf{y}_I(k) + \mathbf{n}(k)) (\mathbf{y}_I(k) + \mathbf{n}(k))^\dagger \mathbf{w}\right)} \end{aligned}$$

As \mathbf{w} and $\mathbf{a}(\theta_S)$ not random, the expectation is $\in \mathbb{R}$ and from Equation 1 and Equation 2

$$\begin{aligned} SINR_{out}(\mathbf{w}) &= \frac{\mathbf{w}^\dagger \mathbf{a}(\theta) \mathbb{E}\left(s(k) s(k)^\dagger\right) \mathbf{a}(\theta)^\dagger \mathbf{w}}{\mathbf{w}^\dagger \mathbb{E}\left((\mathbf{y}_I(k) + \mathbf{n}(k)) (\mathbf{y}_I(k) + \mathbf{n}(k))^\dagger\right) \mathbf{w}} \\ &= \frac{P_S \mathbf{w}^\dagger \mathbf{a}(\theta) \mathbf{a}(\theta)^\dagger \mathbf{w}}{\mathbf{w}^\dagger \mathbf{C} \mathbf{w}} = \frac{P_S |\mathbf{w}^\dagger \mathbf{a}(\theta)|^2}{\mathbf{w}^\dagger \mathbf{C} \mathbf{w}} \end{aligned}$$

From this expression, the maximization of the $SINR(\mathbf{w})$ is equivalent to the minimization of $\mathbf{w}^\dagger \mathbf{C} \mathbf{w}$.

1.1.3 Minimum Power Distortionless Response (MPDR) filter expression

We want to maximize $SINR(\mathbf{w})$. This is equivalent to solving the problem :

$$\arg \max_{\mathbf{w}} SINR(\mathbf{w}) = \arg \min_{\mathbf{w}} \mathbf{w}^\dagger \mathbf{C} \mathbf{w}$$

In practice, the matrix \mathbf{C} is not available and cannot even be reliably estimated. Therefore, we must work with \mathbf{R} , and our objective becomes to maximize the $SSINR$:

$$\arg \max_{\mathbf{w}} SSINR_{out}(\mathbf{w}) = \frac{\mathbb{E}(|\mathbf{w}^\dagger \mathbf{a}(\theta) s(k)|^2)}{\mathbb{E}(|\mathbf{w}^\dagger (\mathbf{a}(\theta_S) s(k) + \mathbf{y}_I(k) + \mathbf{n}(k))|^2)} = \frac{P_S |\mathbf{w}^\dagger \mathbf{a}(\theta)|^2}{\mathbf{w}^\dagger \mathbf{R} \mathbf{w}}$$

The Minimum Power Distortionless Response filter $MPDR$ satisfies the following expression:

$$\mathbf{w}_{MPDR} = \arg \min_{\mathbf{w}} \mathbf{w}^\dagger \mathbf{R} \mathbf{w} \quad (5)$$

We consider, moreover, that we want to verify the unit constraint $|\mathbf{w}^\dagger \mathbf{a}(\theta_S)| = 1$. We end up with a constrained optimization problem, which can be solve by using the Lagrangian of the problem. Let us define the Lagrangian with respect to $\mathbf{w} \in \mathbb{R}^N$ and $\lambda \in \mathbb{C}$ and its associated problem :

$$\{\tilde{\mathbf{w}}, \tilde{\lambda}\} = \arg \min_{\mathbf{w}, \lambda} \mathcal{L}(\lambda, \mathbf{w}) = \arg \min_{\mathbf{w}, \lambda} \mathbf{w}^\dagger \mathbf{R} \mathbf{w} + \lambda \left(\mathbf{w}^\dagger \mathbf{a}(\theta_S) - 1 \right) + \lambda^* \left(\mathbf{a}(\theta_S)^\dagger \mathbf{w} - 1 \right)$$

To solve the optimization problem, we compute the derivatives of the Lagrangian:

$$\mathcal{L}(\lambda, \mathbf{w}) = \mathbf{w}^\dagger \mathbf{R} \mathbf{w} + \lambda \left(\mathbf{w}^\dagger \mathbf{a}(\theta_S) - 1 \right) + \lambda^* \left(\mathbf{a}(\theta_S)^\dagger \mathbf{w} - 1 \right)$$

Step 1: Directional derivative

For any perturbation \mathbf{h} in \mathbf{w} , the first-order (Taylor) expansion gives us:

$$f(\mathbf{w} + \epsilon \mathbf{h}) = f(\mathbf{w}) + \epsilon \left\langle \mathbf{h}, \frac{\partial f}{\partial \mathbf{w}} \right\rangle + o(\epsilon),$$

where $\langle \mathbf{h}, \mathbf{v} \rangle = \mathbf{h}^\dagger \mathbf{v}$ denotes the Hermitian scalar product. We apply this to the two relevant terms:

1. Quadratic term:

$$f_1(\mathbf{w}) = \mathbf{w}^\dagger \mathbf{R} \mathbf{w}$$

Expanding in the direction \mathbf{h} :

$$f_1(\mathbf{w} + \epsilon \mathbf{h}) = (\mathbf{w} + \epsilon \mathbf{h})^\dagger \mathbf{R} (\mathbf{w} + \epsilon \mathbf{h}) = \mathbf{w}^\dagger \mathbf{R} \mathbf{w} + \epsilon \mathbf{h}^\dagger \mathbf{R} \mathbf{w} + \epsilon \mathbf{w}^\dagger \mathbf{R} \mathbf{h} + \mathcal{O}(\epsilon^2)$$

The derivative with respect to \mathbf{w} is:

$$\frac{\partial f_1}{\partial \mathbf{w}} = \mathbf{R} \mathbf{w}$$

2. Linear constraint term:

$$f_2(\mathbf{w}) = \lambda^* (\mathbf{a}(\theta_S)^\dagger \mathbf{w} - 1)$$

Directional derivative:

$$f_2(\mathbf{w} + \epsilon \mathbf{h}) = \lambda^* (\mathbf{a}(\theta_S)^\dagger (\mathbf{w} + \epsilon \mathbf{h}) - 1) = f_2(\mathbf{w}) + \epsilon \lambda^* \mathbf{a}(\theta_S)^\dagger \mathbf{h} + \mathcal{O}(\epsilon^2)$$

so that,

$$\frac{\partial f_2}{\partial \mathbf{w}} = \lambda^* \mathbf{a}(\theta_S)$$

Step 2: Derivative of the Lagrangian

Combining the terms, we obtain:

$$\frac{\partial \mathcal{L}(\lambda, \mathbf{w})}{\partial \mathbf{w}} = \mathbf{R}\mathbf{w} + \lambda^* \mathbf{a}(\theta_S).$$

The derivative with respect to the Lagrange multiplier λ is straightforward:

$$\frac{\partial \mathcal{L}(\lambda, \mathbf{w})}{\partial \lambda} = \mathbf{w}^\dagger \mathbf{a}(\theta_S) - 1.$$

Lagrange conditions give us that the problem solutions are $\{\tilde{\mathbf{w}}, \tilde{\lambda}\}$ such that :

$$\begin{aligned} & \begin{cases} \left. \frac{\partial \mathcal{L}(\lambda, \mathbf{w})}{\partial \lambda} \right|_{\lambda=\tilde{\lambda}, \mathbf{w}=\tilde{\mathbf{w}}} = 0 \\ \left. \frac{\partial \mathcal{L}(\lambda, \mathbf{w})}{\partial \mathbf{w}} \right|_{\lambda=\tilde{\lambda}, \mathbf{w}=\tilde{\mathbf{w}}} = 0 \end{cases} \\ & \Leftrightarrow \begin{cases} \mathbf{R}\tilde{\mathbf{w}} + \tilde{\lambda}^* \mathbf{a}(\theta_S) = 0 \\ \tilde{\mathbf{w}}^\dagger \mathbf{a}(\theta_S) - 1 = 0 \end{cases} \\ & \Leftrightarrow \begin{cases} \tilde{\mathbf{w}} = -\tilde{\lambda}^* \mathbf{R}^{-1} \mathbf{a}(\theta_S) \\ -\tilde{\lambda}^* = \frac{1}{\mathbf{a}(\theta_S)^\dagger \mathbf{R}^{-1} \mathbf{a}(\theta_S)} \end{cases} \end{aligned}$$

By substituting the second equation in the first one, we obtain the *MPDR* filter solution of Equation 5

$$\mathbf{w}_{MPDR} = \tilde{\mathbf{w}} = \frac{\mathbf{R}^{-1} \mathbf{a}(\theta_S)}{\mathbf{a}(\theta_S)^\dagger \mathbf{R}^{-1} \mathbf{a}(\theta_S)}$$

Now that the expressions of the filters are computed, we can implement the two approaches by assuming the error between the true angle and the estimated angle is equal to 2° , and compare them to the Conventional and Adaptive beamforming :

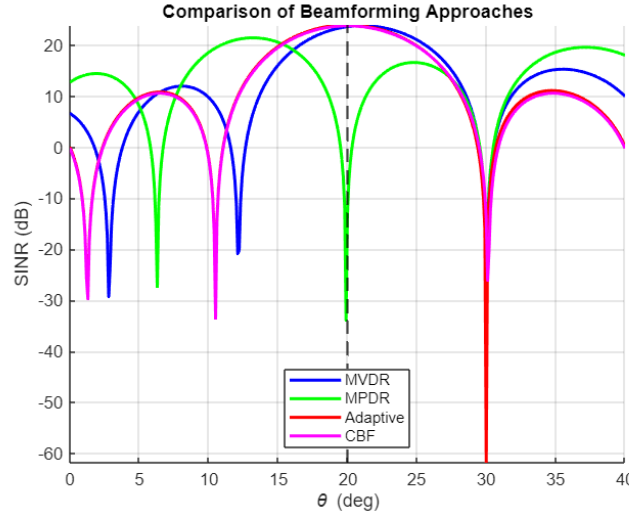


Figure 2: Comparison of Beamforming Approaches

We notice that around $\theta_s = 20^\circ$ the four methods exhibit a strong peak, which corresponds to the direction of the desired signal. This peak represents the main beam of the array response and indicates that each beamforming successfully steers its gain toward the target source. A sharp drop appears around $\theta_i = 30$ which corresponds to the direction of the interfering source. The deeper the drop, the more effective the beamforming is in canceling the interference. Adaptive, MVDR, and MPDR achieve strong interference suppression, whereas the Conventional Beamformer (CBF) produces a much shallower null.

1.2 Unknown covariance matrices

In the second case, we explore a more realistic scenario where the matrices \mathbf{C} and \mathbf{R} are unknown and must be estimated. To accomplish this, we need to collect multiple snapshots of the received signal, denoted as $\mathbf{Y} = \{\mathbf{y}(k)\}_{k \in \llbracket 1, K \rrbracket}$, at different time instances. Typically, the covariance matrices are estimated using:

$$\begin{aligned}\hat{\mathbf{C}} &= \frac{1}{K} \sum_{k=1}^K ((\mathbf{y}_I(k) + \mathbf{n}(k))(\mathbf{y}_I(k) + \mathbf{n}(k))^\dagger) \\ \hat{\mathbf{R}} &= \frac{1}{K} \sum_{k=1}^K \mathbf{y}(k) \mathbf{y}(k)^\dagger\end{aligned}$$

We are now interested in implementing and analyzing beamforming approaches based on these estimations.

1.2.1 Noise probability density reformulation

We have $\mathbf{n}(k) \sim \mathcal{CN}(0, \sigma^2 \mathbf{I})$ with \mathcal{CN} the complex normal distribution:

$$\frac{1}{\pi^N |\boldsymbol{\Sigma}|} \exp^{-(\mathbf{n}(k) - \mathbf{m})^\dagger \boldsymbol{\Sigma}^{-1} (\mathbf{n}(k) - \mathbf{m})} \quad (6)$$

From there, we can apply it specifically for the studied case with $\mathbf{m} = 0$ and $\boldsymbol{\Sigma} = \sigma^2 \mathbf{I}$

$$p(\mathbf{n}(k)) = \frac{1}{(\pi \sigma^2)^N} \exp^{-\frac{1}{\sigma^2} \mathbf{n}(k)^\dagger \mathbf{n}(k)}$$

We want to separate the real and imaginary parts of this signal to simplify the problem with $\mathbf{n}(k) = \mathbf{n}_{re}(k) + j \mathbf{n}_{im}(k)$ where $\mathbf{n}_{re}(k) \in \mathbb{R}^N$ and $\mathbf{n}_{im}(k) \in \mathbb{R}^N$. Under these assumptions, $\mathbf{n}(k)$ hermitian norm becomes:

$$\mathbf{n}(k)^\dagger \mathbf{n}(k) = (\mathbf{n}_{re}(k) - j \mathbf{n}_{im}(k)) (\mathbf{n}_{re}(k) + j \mathbf{n}_{im}(k)) = \|\mathbf{n}_{re}(k)\|^2 + \|\mathbf{n}_{im}(k)\|^2$$

We replace :

$$\begin{aligned}p(\mathbf{n}(k)) &= \frac{1}{(\pi \sigma^2)^{N/2}} \frac{1}{(\pi \sigma^2)^{N/2}} \exp^{-\frac{1}{\sigma^2} (\|\mathbf{n}_{re}(k)\|^2 + \|\mathbf{n}_{im}(k)\|^2)} \\ &= \frac{1}{(2\pi \frac{\sigma^2}{2})^{N/2}} \exp^{-\frac{1}{2\frac{\sigma^2}{2}} \|\mathbf{n}_{re}(k)\|^2} \times \frac{1}{(2\pi \frac{\sigma^2}{2})^{N/2}} \exp^{-\frac{1}{2\frac{\sigma^2}{2}} \|\mathbf{n}_{im}(k)\|^2}\end{aligned}$$

By the real vectorial Gaussian for $\mathcal{N}(\mathbf{m}_0, \sigma_0^2 \mathbf{I})$ with the following distribution:

$$p(\mathbf{x}) = \frac{1}{(2\pi \sigma_0^2)^{N/2}} \exp^{-\frac{1}{2\sigma_0^2} \|\mathbf{x} - \mathbf{m}_0\|^2}$$

We arrive at the separation in a product of two real gaussians:

$$p(\mathbf{n}(k)) = \mathbf{n}_{re}(k) \mathbf{n}_{im}(k) = \mathcal{N}(0, \frac{\sigma^2}{2} \mathbf{I}) \mathcal{N}(0, \frac{\sigma^2}{2} \mathbf{I})$$

We implement the received signal $y(k) \quad \forall k \in \llbracket 1, K \rrbracket$:

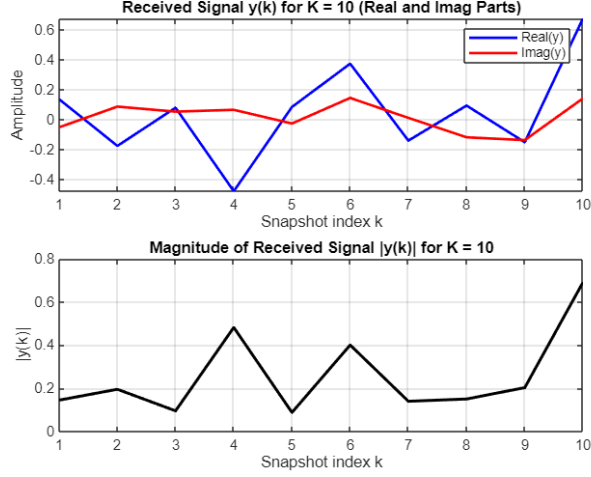


Figure 3: Received signal

This allows us to implement the new *MVDR* and *MPDR* filters with the matrices R and C unknown. By computing the errors, we obtain :

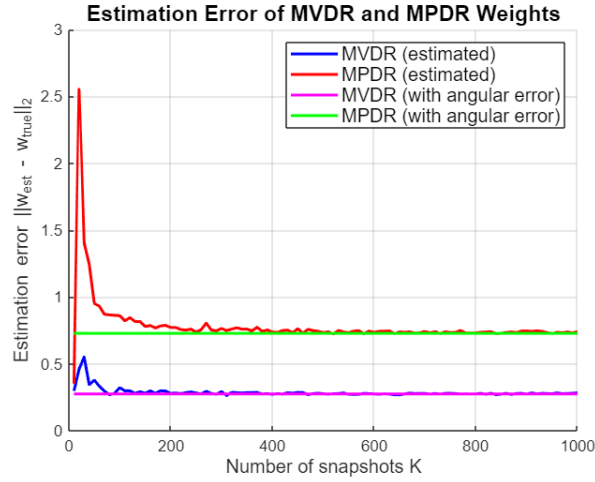


Figure 4: Comparison error of *MVDR* and *MPDR*

We see that the errors of the new filters converged to the filter's values when R and C were known. Overall, the *MPDR* filter has higher error and converges more slowly than the *MVDR* filter. To compare them further, we calculate the *SINR* of the two methods over 50 Monte-Carlo realizations. For this, we consider R and C known. We obtain:

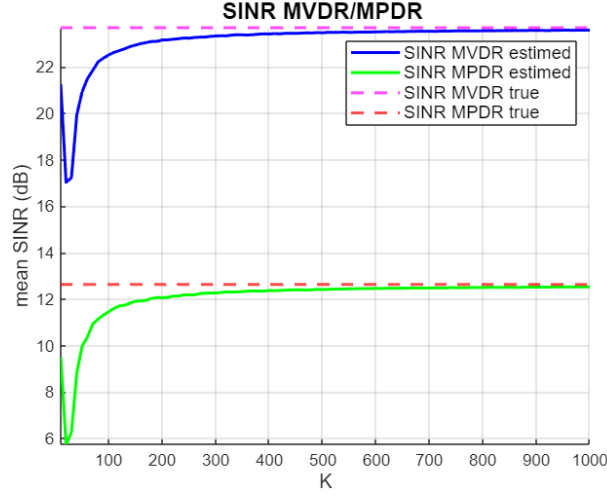


Figure 5: Comparison SINR for MVDR and MPDR

The mean $SINR$ over the 50 realizations for the two methods converge to the real values of $SINR$: as K increases, noise in the sample covariance estimate decreases, so the estimated beamformer weights approach the optimal ones. $MVDR$ reaches a much higher $SINR$ than $MPDR$ and converges faster, confirming its superior interference suppression capability.

1.2.2 Robust MPDR

We implement the Robust MPDR filter, which is a form of MPDR where a factor μ is added to R , for $\mu = [0, 1e^{-5}, 1e^{-4}, 1e^{-3}, 1e^{-2}, 1e^{-1}, 1]$.

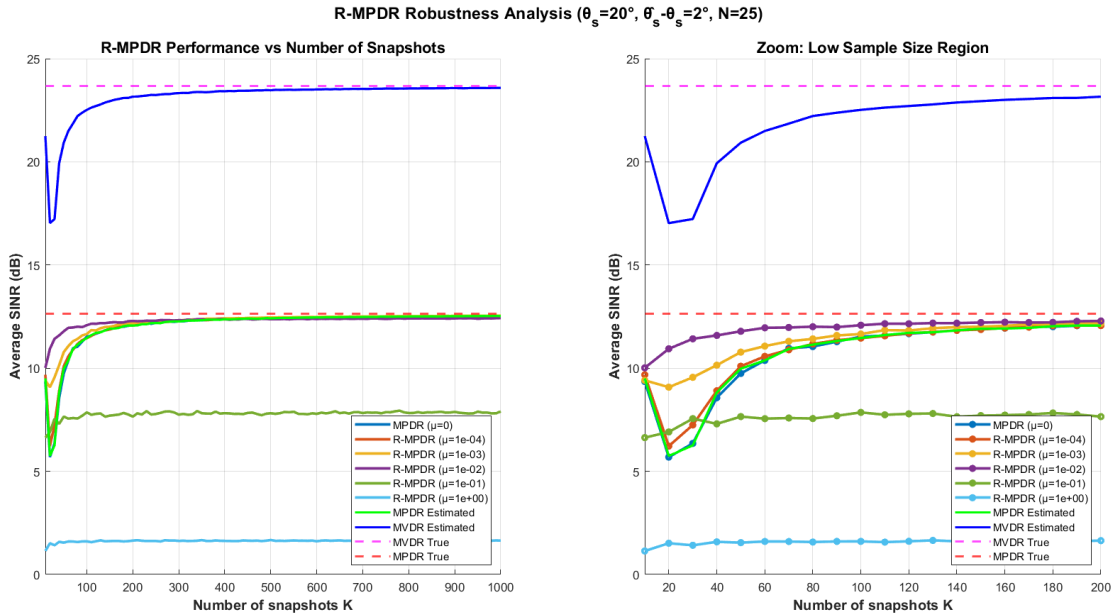


Figure 6: Comparison MVDR, MPDR, and Robust MPDR

When we examine the standard deviation along the diagonal of R , we find values around 0.0089 on average across all Monte Carlo runs. Then, if we introduce the μ factor and set it too large (around 1), its contribution becomes dominant—around 11,000%. In this case, all diagonal terms become nearly identical, which degrades performance since the algorithm can no longer distinguish differences in the covariance values.

Conversely, if μ is too small (around 10^{-5}), its contribution represents only about 0.11%. Its effect is therefore negligible, and no improvement is observed in the beamformer's performance.

With more moderate values of μ (e.g., 10^{-3} in this case), the contribution is around 11 %. It increases the overall diagonal, improves matrix conditioning and stability during inversion, while still preserving enough variation to distinguish between entries. We also note that for $\mu = 0$, the SINR of the robust MPDR collapses to the classic MPDR value, and if μ is too large, the SINR drops.

Therefore, to improve MPDR performance, μ values between 10^{-4} and 10^{-2} are recommended. Overall, this method helps accelerate convergence and improve the accuracy of the MPDR algorithm.

2 Implementation of DoA estimation techniques

In this section, we will implement the various DoA estimation techniques. To achieve this, we assume a new scenario where the ULA network receives three unknown sources with directions $\boldsymbol{\theta}_S = [\theta_1, \theta_2, \theta_3]^\top$. The simulation parameters remain unchanged. The primary modification is that we now consider $P = 3$ sources with $\theta_1 = -25^\circ$, $\theta_2 = 20^\circ$, $\theta_3 = 25^\circ$, and $\mathbf{P}_S = 10$. As previously, the matrices \mathbf{C} and \mathbf{R} are estimated and K is fixed to 500. Additionally, we assume the presence of interference. Consequently, the signal is modeled by:

$$\mathbf{y}(k) = \mathbf{A}(\boldsymbol{\theta}_S) \mathbf{s}(k) + \mathbf{n}(k) \quad \mathbf{n}(k) \sim \mathcal{CN}(0, \sigma^2 \mathbf{I}) \quad (7)$$

with $\mathbf{A}(\boldsymbol{\theta}_S) = [\mathbf{a}(\theta_1), \mathbf{a}(\theta_2), \mathbf{a}(\theta_3)]$ and for each k $\mathbf{s}(k) \in \mathbb{C}^3$.

2.1 Implementation of non-parametric approaches

Non-parametric methods do not assume a specific model for the source signals: they work directly with the output power received after beamforming. We start by implementing the two methods based on the conventional \mathbf{w}_{CBF} and adaptive beamforming with Capon's method $\hat{\mathbf{w}}_{MPDR}$, with interferences. We obtain:

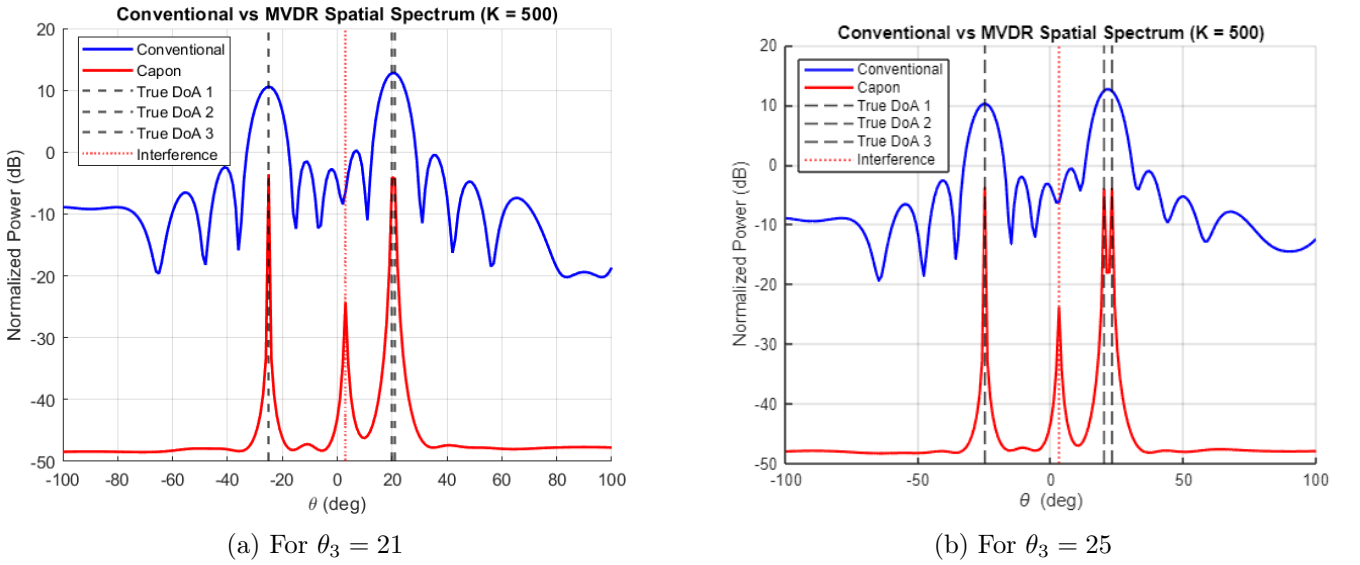


Figure 7: Conventional and Capon comparison

We can see that both methods show spectral peaks at the true angles of arrival, indicated by the dashed vertical lines. However, the Capon exhibits much sharper and narrower peaks than the conventional method. This means *MVDR* is more effective at separating closely spaced sources. The conventional beamformer has large side lobes while the *MVDR* suppresses them, dropping more than 40 dB below the main peaks. The *MVDR* spectrum includes deep nulls at the positions of interferers, unlike the conventional method. This is a signature of adaptive beamforming: *MVDR* actively cancels interference while preserving the desired signal. When the third source is $\theta_3 = 21^\circ$, we see that the second and third DoA are not accurately separated. The methods are not precise enough.

2.2 Implementation of parametric approaches

2.2.1 MUSIC method

In the MUSIC method, we assume the structure of the signal covariance matrix $\mathbf{R} = \mathbf{A}(\theta)\mathbf{R}_S\mathbf{A}(\theta)^\dagger + \sigma^2\mathbf{I}$. This is considered prior information, so that this method can be considered parametric.

Power expression with the MUSIC approach

We first assume: $\mathbf{R} = \mathbf{A}(\theta)\mathbf{R}_S\mathbf{A}(\theta)^\dagger + \sigma^2\mathbf{I}$

We decompose $\mathbf{A}(\theta)\mathbf{R}_S\mathbf{A}(\theta)^\dagger$ into an orthonormal basis of eigenvectors $\{\mathbf{u}_i\}_{i \in [1;N]}$. \mathbf{U} composed of all \mathbf{u}_i , then gives the following equation :

$$\mathbf{A}(\theta)\mathbf{R}_S\mathbf{A}(\theta)^\dagger = \mathbf{U}\mathbf{\Lambda}\mathbf{U}^\dagger$$

where $\mathbf{\Lambda}$ is the diagonal matrix of the eigenvalues λ_i associated with the eigenvectors \mathbf{u}_i . The matrix $\mathbf{A}(\theta)\mathbf{R}_S\mathbf{A}(\theta)^\dagger$ is of rank P , which means that (in descending order) the P first eigenvalues are non-zero, and the $N - P$ others are null. By separating \mathbf{U} into two separate parts $\mathbf{U}_S = [\mathbf{u}_1, \dots, \mathbf{u}_P]$ and $\mathbf{U}_N = [\mathbf{u}_{P+1}, \dots, \mathbf{u}_N]$, we have :

$$\mathbf{A}(\theta)\mathbf{R}_S\mathbf{A}(\theta)^\dagger = \sum_{p=1}^P \lambda_p \mathbf{u}_p \mathbf{u}_p^\dagger = \mathbf{U}_S \mathbf{\Lambda}_S \mathbf{U}_S^\dagger$$

And then, the orthogonality of the matrix \mathbf{U}_N (from the SVD definition) yields to :

$$\mathbf{R} = \mathbf{U}_S \mathbf{\Lambda}_S \mathbf{U}_S^\dagger + \sigma^2 \mathbf{U}_N \mathbf{U}_N^\dagger \quad (8)$$

To compute the estimated power, subsection 2.2.1 gives us that both $\mathbf{a}(\theta_S)$ and \mathbf{U}_S generate the same space, meaning they are different bases of the same space. Then, as we know $\mathbf{U}_S \perp \mathbf{U}_N$, we obtain $\mathbf{a}(\theta_S) \perp \mathbf{U}_N$. We also have $\mathbf{a}(\theta_S)^\dagger \cdot \mathbf{U}_N = 0$. Numerically, a power is the sum of the energies of the studied signal. Here, in the considered subspace, it can be written as $\mathbf{a}(\theta_S)^\dagger \mathbf{U}_N \mathbf{U}_N^\dagger \mathbf{a}(\theta_S)$. We want :

$$P_{MUSIC}(\theta) = \frac{1}{\mathbf{a}(\theta_S)^\dagger \mathbf{U}_S \mathbf{U}_S^\dagger \mathbf{a}(\theta_S)}$$

From our assumptions, we determine $\hat{\mathbf{U}}_S$ from \mathbf{R} . In reality, we do not have access to \mathbf{R} so we have to estimate it from a distribution. Meaning \mathbf{R} becomes $\hat{\mathbf{R}}$ and \mathbf{U}_S becomes $\hat{\mathbf{U}}_S$:

$$P_{MUSIC}(\theta) = \frac{1}{\mathbf{a}(\theta_S)^\dagger \hat{\mathbf{U}}_S \hat{\mathbf{U}}_S^\dagger \mathbf{a}(\theta_S)}$$

We implement this method with interferences. We set θ between $[-100^\circ, +100^\circ]$.

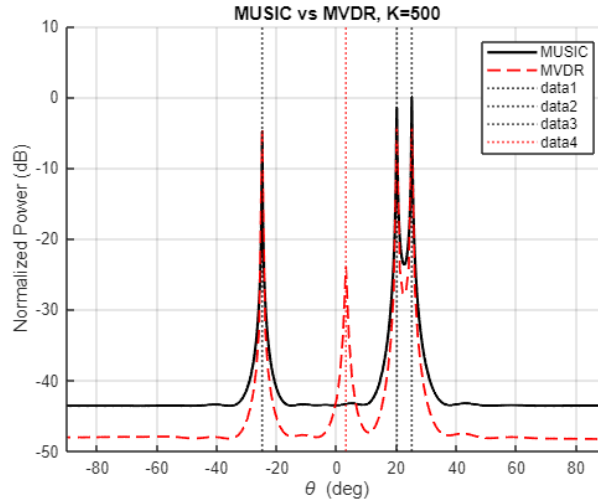


Figure 8: Comparison Music and Capon

MUSIC shows very sharp, narrow peaks aligned with the true DoAs. This is expected because *MUSIC* exploits the orthogonality between the noise subspace and the array steering vectors. As a result, *MUSIC* can separate closely spaced sources even when conventional or *MVDR* beamforming struggles. On the other hand, *MVDR* also shows peaks at the correct DoAs, but the peaks are wider, less pronounced, and sometimes slightly shifted depending on estimation noise. So, overall, *MUSIC* exhibits a stronger discrimination capability.

We added interferences, and *MVDR* shows a suppressed response near the interference DoA, because it tries to place a null while maintaining gain toward the desired direction, but *MUSIC* does not. Which leads us to say that *MUSIC* is a detection or estimation technique, while *MVDR* is also more of a beamforming or filtering technique.

2.2.2 Maximum Likelihood method

Now, we propose to implement the parametric approach based on the maximization of the likelihood so that:

$$\boldsymbol{\theta}_S = \arg \max_{\boldsymbol{\theta}_S} p(\mathbf{Y}|\boldsymbol{\theta}_S) \quad \mathbf{Y} = \{\mathbf{y}(1), \dots, \mathbf{y}(K)\} \quad (9)$$

Particularly, we propose to study the deterministic likelihood approach, which assumes that $s(k)$ is not random.

1. When s and σ^2 are known

Theoretical estimator expression

From Equation 7, when $\mathbf{s}(k)$ modeling is known :

$$\mathbf{y}(k) = \mathbf{A}(\boldsymbol{\theta}_S) \mathbf{s}(k) + \mathbf{n}(k) \quad \mathbf{n}(k) \sim \mathcal{CN}(0, \sigma^2 \mathbf{I}) \quad \Leftrightarrow \quad \mathbf{y}(k) \sim \mathcal{CN}(\mathbf{A}(\boldsymbol{\theta}_S) \mathbf{s}(k), \sigma^2 \mathbf{I})$$

The equation Equation 9 becomes :

$$\begin{aligned} \hat{\boldsymbol{\theta}}_{S,ML1} &= \arg \max_{\boldsymbol{\theta}_S} p(\mathbf{Y}|\boldsymbol{\theta}_S) = \arg \min_{\boldsymbol{\theta}_S} -\ln p(\mathbf{Y}|\boldsymbol{\theta}_S) \\ &= \arg \min_{\boldsymbol{\theta}_S} -\ln \left(\prod_{k=1}^K p(\mathbf{y}(k)|\boldsymbol{\theta}_S) \right) \end{aligned}$$

Supposing independence and replacing with the complex Gaussian from Equation 6 :

$$\hat{\boldsymbol{\theta}}_{S,ML1} = \arg \min_{\boldsymbol{\theta}_S} -\ln \prod_{k=1}^K \left(\frac{1}{\pi^N |\boldsymbol{\Sigma}|} \exp^{-\frac{1}{\sigma^2} (\mathbf{y}(k) - \mathbf{A}(\boldsymbol{\theta}_S) \mathbf{s}(k))^{\dagger} (\mathbf{y}(k) - \mathbf{A}(\boldsymbol{\theta}_S) \mathbf{s}(k))} \right)$$

All terms independent in $\boldsymbol{\theta}_S$ do not influence the arg min, they are only a constant with respect to this problem. Let's develop by simplifying them :

$$\begin{aligned} \hat{\boldsymbol{\theta}}_{S,ML1} &= \arg \min_{\boldsymbol{\theta}_S} -\sum_{k=1}^K \left(-\frac{1}{\sigma^2} (\mathbf{y}(k) - \mathbf{A}(\boldsymbol{\theta}_S) \mathbf{s}(k))^{\dagger} (\mathbf{y}(k) - \mathbf{A}(\boldsymbol{\theta}_S) \mathbf{s}(k)) \right) \\ &= \arg \min_{\boldsymbol{\theta}_S} \sum_{k=1}^K \left((\mathbf{y}(k) - \mathbf{A}(\boldsymbol{\theta}_S) \mathbf{s}(k))^{\dagger} (\mathbf{y}(k) - \mathbf{A}(\boldsymbol{\theta}_S) \mathbf{s}(k)) \right) \\ &= \arg \min_{\boldsymbol{\theta}_S} \sum_{k=1}^K \|\mathbf{y}(k) - \mathbf{A}(\boldsymbol{\theta}_S) \mathbf{s}(k)\|^2 \end{aligned}$$

When s is known, this expression can be computed directly because $\mathbf{y}(k)$ is the received signal, $\mathbf{A}(\theta)$ is by design. A Gauss-Newton algorithm can solve for our only unknown $\boldsymbol{\theta}_S$.

Cramér-Rao Bound expression

To find the expression of the *CRB*, we use the Slepian Bangs formula. The general expression of the Fisher information matrix terms $\forall l \in \llbracket 1, P \rrbracket, \forall l' \in \llbracket 1, P \rrbracket$ considering $\mathbf{y}(k) \sim \mathcal{CN}(\boldsymbol{\mu}_k(x), \boldsymbol{\Sigma}(x))$ is :

$$[\mathbf{I}(x)]_{l,l'} = 2 \sum_{k=1}^K \operatorname{Re} \left(\frac{\partial \boldsymbol{\mu}_k(x)}{\partial (x)_l}^{\dagger} \boldsymbol{\Sigma}(x)^{-1} \frac{\partial \boldsymbol{\mu}_k(x)}{\partial (x)_{l'}} \right) + K \operatorname{tr} \left(\boldsymbol{\Sigma}(x)^{-1} \frac{\partial \boldsymbol{\Sigma}(x)}{\partial (x)_l} \boldsymbol{\Sigma}(x)^{-1} \frac{\partial \boldsymbol{\Sigma}(x)}{\partial (x)_{l'}} \right)$$

In our case, we have $\boldsymbol{\mu}_k(\boldsymbol{\theta}) = \mathbf{A}(\boldsymbol{\theta}_S) \mathbf{s}(k)$ and $\boldsymbol{\Sigma}(\boldsymbol{\theta}) = \sigma^2 \mathbf{I}$ independent from $\boldsymbol{\theta}$ (derivatives are 0), so that it simplifies to the first (mean) term only:

$$[\mathbf{I}(\boldsymbol{\theta}_S)]_{l,l'} = \frac{2}{\sigma^2} \sum_{k=1}^K \operatorname{Re} \left(\frac{\partial \boldsymbol{\mu}_k(\boldsymbol{\theta}_S)}{\partial (\boldsymbol{\theta}_S)_l}^{\dagger} \frac{\partial \boldsymbol{\mu}_k(\boldsymbol{\theta}_S)}{\partial (\boldsymbol{\theta}_S)_{l'}} \right)$$

Using $\boldsymbol{\mu}_k(\boldsymbol{\theta}_S) = \mathbf{A}(\boldsymbol{\theta}_S)\mathbf{s}(k)$, we obtain

$$\frac{\partial \boldsymbol{\mu}_k(\boldsymbol{\theta}_S)}{\partial (\boldsymbol{\theta}_S)_l} = \frac{\partial \mathbf{A}(\boldsymbol{\theta}_S)}{\partial (\boldsymbol{\theta}_S)_l} \mathbf{s}(k) = \mathbf{D}\mathbf{A}_l \mathbf{s}(k)$$

so that the Fisher Information Matrix entries become:

$$[\mathbf{I}(\boldsymbol{\theta}_S)]_{l,l'} = \frac{2}{\sigma^2} \sum_{k=1}^K \operatorname{Re} \left(\mathbf{s}(k)^\dagger \mathbf{D}\mathbf{A}_l^\dagger \mathbf{D}\mathbf{A}_{l'} \mathbf{s}(k) \right).$$

Finally, the Cramér–Rao bound for this case is given by:

$$\text{CRB}(\boldsymbol{\theta}_S) = \mathbf{I}(\boldsymbol{\theta}_S)^{-1}.$$

2. When s and σ^2 are unknown

Theoretical estimator expression

When neither s nor σ^2 are known, the terms with σ^2 that we integrated in the constant for the development above now have to be considered. By developing from Equation 2.2.2 like above, we end up with:

$$\hat{\boldsymbol{\theta}}_{S,ML2} = \arg \min_{\boldsymbol{\theta}_S, \sigma^2, \mathbf{s}(k)} -K \ln \left(\frac{1}{\pi^N |\boldsymbol{\Sigma}|} \right) + \frac{1}{\sigma^2} \sum_{k=1}^K \|\mathbf{y}(k) - \mathbf{A}(\boldsymbol{\theta}_S)\mathbf{s}(k)\|^2$$

With $\boldsymbol{\Sigma} = \sigma^2 \mathbf{I}$, we have $|\boldsymbol{\Sigma}| = \sigma^{2N}$, and the term in π is still a constant for the minimization:

$$\hat{\boldsymbol{\theta}}_{S,ML2} = \arg \min_{\boldsymbol{\theta}_S, \sigma^2, \mathbf{s}(k)} KN \ln(\sigma^2) + \frac{1}{\sigma^2} \sum_{k=1}^K \|\mathbf{y}(k) - \mathbf{A}(\boldsymbol{\theta}_S)\mathbf{s}(k)\|^2 \quad (10)$$

First, we minimize with respect to $\mathbf{s}(k)$ is a classical linear least-squares problem, so it can be expressed easily:

$$\hat{\mathbf{s}}(k) = \arg \min_{\mathbf{s}(k)} \|\mathbf{y}(k) - \mathbf{A}(\boldsymbol{\theta}_S)\mathbf{s}(k)\|^2 = \arg \min_{\mathbf{s}(k)} (\mathbf{y}(k) - \mathbf{A}(\boldsymbol{\theta}_S)\mathbf{s}(k))^\dagger (\mathbf{y}(k) - \mathbf{A}(\boldsymbol{\theta}_S)\mathbf{s}(k))$$

The normal equation directly gives us:

$$\hat{\mathbf{s}}(k) = (\mathbf{A}(\boldsymbol{\theta}_S)^\dagger \mathbf{A}(\boldsymbol{\theta}_S))^{-1} \mathbf{A}(\boldsymbol{\theta}_S)^\dagger \mathbf{y}(k)$$

This estimator is the unique vector whose image through $\mathbf{A}(\boldsymbol{\theta}_S)$ is the orthogonal projection of $\mathbf{y}(k)$ onto the column space of $\mathbf{A}(\boldsymbol{\theta}_S)$. This defines the orthogonal projector and complement:

$$\mathbf{P}_\mathbf{A}(\boldsymbol{\theta}_S) = \mathbf{A}(\boldsymbol{\theta}_S) (\mathbf{A}(\boldsymbol{\theta}_S)^\dagger \mathbf{A}(\boldsymbol{\theta}_S))^{-1} \mathbf{A}(\boldsymbol{\theta}_S)^\dagger \quad \text{and} \quad \mathbf{P}_\mathbf{A}^\perp(\boldsymbol{\theta}_S) = \mathbf{I} - \mathbf{P}_\mathbf{A}(\boldsymbol{\theta}_S)$$

With these defined, the vector fitted to this particular subspace becomes $\hat{\mathbf{y}}(k) = \mathbf{A}(\boldsymbol{\theta}_S)\mathbf{s}(k) = \mathbf{P}_\mathbf{A}(\boldsymbol{\theta}_S)\mathbf{y}(k)$. Thus, the solution becomes:

$$\arg \min_{\mathbf{s}(k)} \|\mathbf{y}(k) - \mathbf{A}(\boldsymbol{\theta}_S)\mathbf{s}(k)\|^2 = \mathbf{y}(k)^\dagger \mathbf{P}_\mathbf{A}^\perp(\boldsymbol{\theta}_S) \mathbf{y}(k)$$

After projecting out the signal subspace, the log-likelihood from Equation 10 becomes independent from $\mathbf{s}(k)$:

$$\hat{\boldsymbol{\theta}}_{S,ML2} = \arg \min_{\boldsymbol{\theta}_S, \sigma^2} KN \ln(\sigma^2) + \frac{1}{\sigma^2} \sum_{k=1}^K \mathbf{y}(k)^\dagger \mathbf{P}_\mathbf{A}^\perp(\boldsymbol{\theta}_S) \mathbf{y}(k)$$

With respect to σ^2 it takes the form:

$$f(\sigma^2) = A \ln(\sigma^2) + \frac{B}{\sigma^2} \quad \text{with} \quad A = KN, \quad B = \sum_{k=1}^K \mathbf{y}(k)^\dagger \mathbf{P}_\mathbf{A}^\perp(\boldsymbol{\theta}_S) \mathbf{y}(k)$$

Differentiating with respect to σ^2 and setting the derivative to zero gives

$$\frac{\partial f}{\partial \sigma^2} = \frac{A}{\sigma^2} - \frac{B}{(\sigma^2)^2} = 0 \implies \hat{\sigma}^2 = \frac{B}{A}$$

Substituting this expression back into the likelihood removes all terms independent of $\boldsymbol{\theta}_S$. Therefore, the maximum likelihood estimator of $\boldsymbol{\theta}_S$ is obtained by minimizing:

$$\hat{\boldsymbol{\theta}}_{S, ML2} = \arg \min_{\boldsymbol{\theta}_S} \frac{1}{K} \sum_{k=1}^K \mathbf{y}(k)^\dagger \mathbf{P}_A^\perp(\boldsymbol{\theta}_S) \mathbf{y}(k)$$

We use that $\mathbf{y}(k)^\dagger \mathbf{P}_A^\perp(\boldsymbol{\theta}_S) \mathbf{y}(k) \in \mathbb{C}$ to replace with $\mathbf{y}(k)^\dagger \mathbf{P}_A^\perp(\boldsymbol{\theta}_S) \mathbf{y}(k) = \text{tr}(\mathbf{y}(k)^\dagger \mathbf{P}_A^\perp(\boldsymbol{\theta}_S) \mathbf{y}(k))$, and as tr is a circular operator, the equality extends to $= \text{tr}(\mathbf{P}_A^\perp(\boldsymbol{\theta}_S) \mathbf{y}(k) \mathbf{y}(k)^\dagger)$. We then have by linearity of both sum and tr :

$$\begin{aligned} \hat{\boldsymbol{\theta}}_{S, ML2} &= \arg \min_{\boldsymbol{\theta}_S} \frac{1}{K} \sum_{k=1}^K \text{tr}(\mathbf{P}_A^\perp(\boldsymbol{\theta}_S) \mathbf{y}(k) \mathbf{y}(k)^\dagger) \\ &= \arg \min_{\boldsymbol{\theta}_S} \text{tr}(\mathbf{P}_A^\perp(\boldsymbol{\theta}_S) \underbrace{\frac{1}{K} \sum_{k=1}^K \mathbf{y}(k) \mathbf{y}(k)^\dagger}_{\hat{\mathbf{R}}}) \\ \hat{\boldsymbol{\theta}}_{S, ML2} &= \arg \min_{\boldsymbol{\theta}_S} \text{tr}(\mathbf{P}_A^\perp(\boldsymbol{\theta}_S) \hat{\mathbf{R}}) \end{aligned}$$

with $\mathbf{P}_A^\perp(\boldsymbol{\theta}_S) = \mathbf{I} - \mathbf{A}(\boldsymbol{\theta}_S)(\mathbf{A}(\boldsymbol{\theta}_S)^\dagger \mathbf{A}(\boldsymbol{\theta}_S))^{-1} \mathbf{A}(\boldsymbol{\theta}_S)^\dagger$

Cramér-Rao Bound expression

To find the expression of the *CRB*, we use the Slepian Bangs formula. The general expression of the Fisher information matrix terms $\forall l \in \llbracket 1, P \rrbracket, \forall l' \in \llbracket 1, P \rrbracket$ considering $\mathbf{y}(k) \sim \mathcal{CN}(\boldsymbol{\mu}_k(x), \boldsymbol{\Sigma}(x))$ is :

$$[\mathbf{I}(x)]_{l, l'} = 2 \sum_{k=1}^K \text{Re} \left(\frac{\partial \boldsymbol{\mu}_k(x)^\dagger}{\partial (x)_l} \boldsymbol{\Sigma}(x)^{-1} \frac{\partial \boldsymbol{\mu}_k(x)}{\partial (x)_{l'}} \right) + K \text{tr} \left(\boldsymbol{\Sigma}(x)^{-1} \frac{\partial \boldsymbol{\Sigma}(x)}{\partial (x)_l} \boldsymbol{\Sigma}(x)^{-1} \frac{\partial \boldsymbol{\Sigma}(x)}{\partial (x)_{l'}} \right)$$

In our case, we have $\boldsymbol{\mu}_k(\boldsymbol{\theta}) = 0$ and $\boldsymbol{\Sigma}(\boldsymbol{\theta}) = \mathbf{R}$, so that it simplifies to:

$$[\mathbf{I}(\boldsymbol{\theta})]_{l, l'} = K \text{tr} \left(\mathbf{R}^{-1} \frac{\partial \mathbf{R}}{\partial (\boldsymbol{\theta})_l} \mathbf{R}^{-1} \frac{\partial \mathbf{R}}{\partial (\boldsymbol{\theta})_{l'}} \right)$$

We need to compute \mathbf{R} derivative for each $l \in \llbracket 1, P \rrbracket$ and $l' \in \llbracket 1, P \rrbracket$. Let l be in $\llbracket 1, P \rrbracket$:

$$\frac{\partial \mathbf{R}}{\partial (\boldsymbol{\theta})_l} = \frac{\partial \mathbf{A}(\boldsymbol{\theta}) \mathbf{R}_S \mathbf{A}(\boldsymbol{\theta})^\dagger + \sigma^2 \mathbf{I}}{\partial (\boldsymbol{\theta})_l} = \frac{\partial \mathbf{A}(\boldsymbol{\theta})}{\partial (\boldsymbol{\theta})_l} \mathbf{R}_S \mathbf{A}(\boldsymbol{\theta})^\dagger + \mathbf{A}(\boldsymbol{\theta}) \mathbf{R}_S \frac{\partial \mathbf{A}(\boldsymbol{\theta})^\dagger}{\partial (\boldsymbol{\theta})_l} \quad (11)$$

Now, let's decompose the matrix derivative of $\mathbf{A}(\boldsymbol{\theta}) \in \mathcal{M}_{N, P}(\mathbb{R})$ for our l index :

$$\frac{\partial \mathbf{A}(\boldsymbol{\theta})}{\partial (\boldsymbol{\theta})_l} = \left[\frac{\partial \mathbf{a}(\theta_1)}{\partial (\boldsymbol{\theta})_l}, \frac{\partial \mathbf{a}(\theta_2)}{\partial (\boldsymbol{\theta})_l}, \dots, \frac{\partial \mathbf{a}(\theta_P)}{\partial (\boldsymbol{\theta})_l} \right] \quad (12)$$

Let k be in $\llbracket 1; P \rrbracket$. If $k \neq l$, $\frac{\partial \mathbf{a}(\theta_k)}{\partial (\boldsymbol{\theta})_l} = 0$ and if $k = l$, $\frac{\partial \mathbf{a}(\theta_k)}{\partial (\boldsymbol{\theta})_l} = \frac{\partial \mathbf{a}(\theta_l)}{\partial (\boldsymbol{\theta})_l}$. Now, let's decompose the derivative of $\mathbf{a}(\theta_l) \in \mathbb{R}^N$, with respect to $(\boldsymbol{\theta})_l$. From Equation 3 as it's a ULA network with N antennas:

$$\mathbf{a}(\theta_l) = \{\exp^{-2\pi j(n-1)\frac{d}{\lambda} \sin(\theta_l)}\}_{n \in \llbracket 1; N \rrbracket} = \begin{pmatrix} 1 \\ \vdots \\ \exp^{-2\pi j(n-1)\frac{d}{\lambda} \sin(\theta_l)} \\ \vdots \\ \exp^{-2\pi j(N-1)\frac{d}{\lambda} \sin(\theta_l)} \end{pmatrix}$$

Meaning that the vectorial derivative of $\mathbf{a}(\theta_l)$ is:

$$\frac{\partial \mathbf{a}(\theta_l)}{\partial (\boldsymbol{\theta})_l} = \frac{\partial}{\partial (\boldsymbol{\theta})_l} \begin{pmatrix} 1 \\ \vdots \\ \exp^{-2\pi j(n-1)\frac{d}{\lambda} \sin(\theta_l)} \\ \vdots \\ \exp^{-2\pi j(N-1)\frac{d}{\lambda} \sin(\theta_l)} \end{pmatrix} = \begin{pmatrix} 0 \\ \vdots \\ 2\pi j(n-1)\frac{d}{\lambda} \cos(\theta_l) \exp^{-2\pi j(n-1)\frac{d}{\lambda} \sin(\theta_l)} \\ \vdots \\ 2\pi j(N-1)\frac{d}{\lambda} \cos(\theta_l) \exp^{-2\pi j(N-1)\frac{d}{\lambda} \sin(\theta_l)} \end{pmatrix}$$

Put that back in Equation 12. We define $[\mathbf{DA}_l] \in \mathcal{M}_{N,P}(\mathbb{R})$ such as:

$$[\mathbf{DA}_l] := \frac{\partial \mathbf{A}(\boldsymbol{\theta})}{\partial (\boldsymbol{\theta})_l} = \begin{pmatrix} 0 & \cdots & 0 & \cdots & 0 \\ \vdots & \vdots & \vdots & \vdots & \vdots \\ 0 & \cdots & 2\pi j(n-1)\frac{d}{\lambda} \cos(\theta_l) \exp^{-2\pi j(n-1)\frac{d}{\lambda} \sin(\theta_l)} & \cdots & 0 \\ \vdots & \vdots & \vdots & \vdots & \vdots \\ 0 & \cdots & \underbrace{2\pi j(N-1)\frac{d}{\lambda} \cos(\theta_l) \exp^{-2\pi j(N-1)\frac{d}{\lambda} \sin(\theta_l)}}_{\text{column } l} & \cdots & 0 \end{pmatrix}$$

Then once this is computed, we can express $[\mathbf{DR}_l] \in \mathcal{M}_{N,P}(\mathbb{R})$ when substituting in Equation 11:

$$[\mathbf{DR}_l] := \mathbf{DA}_l \mathbf{R}_S \mathbf{A}(\boldsymbol{\theta})^\dagger + \mathbf{A}(\boldsymbol{\theta}) \mathbf{R}_S \mathbf{DA}_l^\dagger$$

With the same expression for l' the derivative. We finally have $\forall l \in \llbracket 1, P \rrbracket, \forall l' \in \llbracket 1, P \rrbracket$:

$$[\mathbf{I}(\boldsymbol{\theta})]_{l,l'} = K \text{Tr}(\mathbf{R}^{-1} [\mathbf{DR}_l] \mathbf{R}^{-1} [\mathbf{DR}_{l'}])$$

Let $\boldsymbol{\theta} = [\theta_1, \dots, \theta_P]^T$ be the parameter vector. The Fisher Information Matrix is computed element-wise as follows.

For all $(l, l') \in \llbracket 1, P \rrbracket^2$:

- Computation of the Jacobian matrices for \mathbf{A} :

$$[\mathbf{DA}_l] = \frac{\partial \mathbf{A}(\boldsymbol{\theta})}{\partial \theta_l}, \quad [\mathbf{DA}_{l'}] = \frac{\partial \mathbf{A}(\boldsymbol{\theta})}{\partial \theta_{l'}}.$$

- Computation of the Jacobian matrices for \mathbf{R} :

$$[\mathbf{DR}_l] = \mathbf{DA}_l \mathbf{R}_S \mathbf{A}(\boldsymbol{\theta})^\dagger + \mathbf{A}(\boldsymbol{\theta}) \mathbf{R}_S \mathbf{DA}_l^\dagger$$

$$[\mathbf{DR}_{l'}] = [\mathbf{DA}_{l'}] \mathbf{R}_S \mathbf{A}(\boldsymbol{\theta})^\dagger + \mathbf{A}(\boldsymbol{\theta}) \mathbf{R}_S [\mathbf{DA}_{l'}]^\dagger$$

- From the (l, l') entry of the Fisher matrix:

$$[\mathbf{I}(\boldsymbol{\theta})]_{l,l'} = K \text{Tr}(\mathbf{R}^{-1} [\mathbf{DR}_l] \mathbf{R}^{-1} [\mathbf{DR}_{l'}])$$

Once the Fisher Information Matrix $\mathbf{I}(\boldsymbol{\theta})$ is completed, the Cramér–Rao Bound is obtained from its inverse:

$$CRB(\boldsymbol{\theta}) = \mathbf{I}(\boldsymbol{\theta})^{-1}$$

Implementation

We implemented the Maximum Likelihood (ML) estimation approach for both cases: when the source signals s are known and when they are unknown. In the case where s is known, the CRB is identical for all sources because the signals are assumed perfectly known and uncorrelated, so the only limiting factor is the additive noise. When s is unknown, the CRB is significantly higher, particularly for closely spaced sources, reflecting the increased difficulty in separating them due to the correlation between their steering vectors.

The Monte-Carlo simulations show that the ML estimator is sensitive to initialization and noise. For known signals, the algorithm generally converges even with moderate random perturbations of the initial

angles, and the estimation error closely follows the CRB when convergence is achieved. When both s and σ^2 are unknown, convergence is much more difficult. Small deviations in the initial guess can prevent convergence, and only very good initializations or low-noise conditions allow the algorithm to reach the global minimum.

Overall, knowing the source signals simplifies the DoA estimation problem and leads to much lower CRB and more robust convergence. When the signals are unknown, the estimation is more challenging, especially for closely spaced sources, highlighting the practical need for careful initialization or the use of subspace-based methods to improve performance.

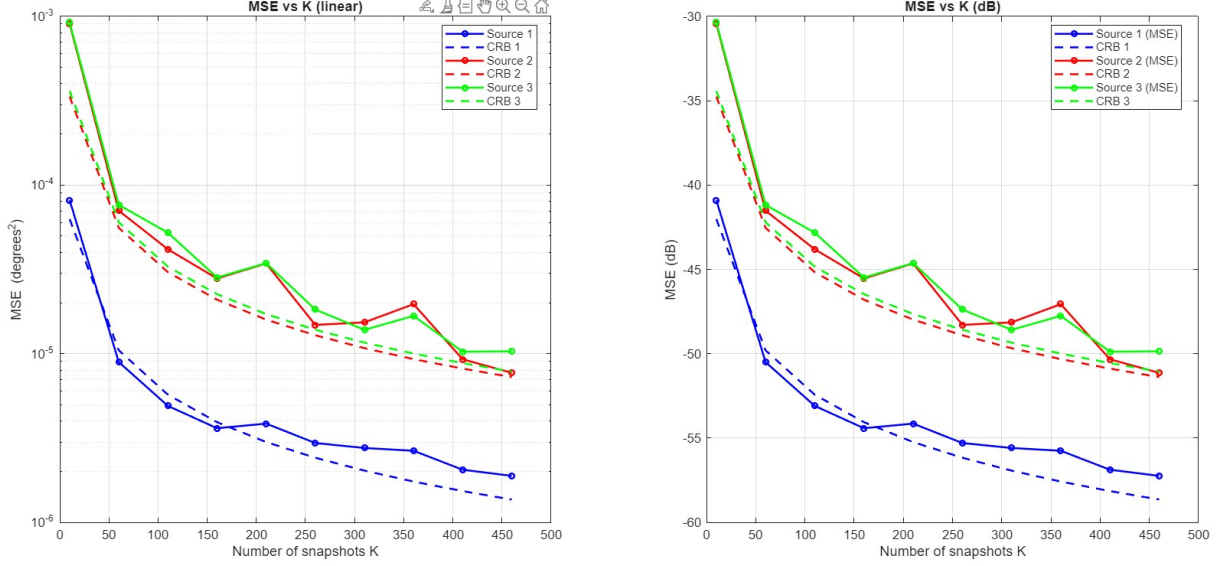


Figure 9: MSE for MLE of the three sources

Although maximum likelihood estimators are statistically optimal, their practical implementation is challenging due to the highly non-convex nature of the likelihood function. When the source signals are unknown, the estimation problem becomes significantly more difficult, exhibiting strong sensitivity to initialization and noise. In practice, ML is often used as a benchmark or combined with subspace-based methods for improved robustness.