$$\det(\mathbf{H}_{2}) = \left(\frac{\alpha}{k_{2}} + \frac{4}{k_{1}} - 2\right) \left(\frac{\alpha}{k_{1}} + \frac{4}{k_{2}} - 2\right) - \alpha^{2} |g_{11}|^{2} |g_{22}|^{2}$$
$$= \left(\frac{\alpha}{k_{2}} + \frac{4}{k_{1}} - 2\right) \left(\frac{\alpha}{k_{1}} + \frac{4}{k_{2}} - 2\right) - \alpha^{2} \frac{1}{k_{1}} \frac{1}{k_{2}}$$
$$= \frac{\alpha}{k_{2}} \left(\frac{4}{k_{2}} - 2\right) + \frac{\alpha}{k_{1}} \left(\frac{4}{k_{1}} - 2\right) + \left(\frac{4}{k_{1}} - 2\right) \left(\frac{4}{k_{2}} - 2\right).$$

Again, if sources are sub-Gaussian, condition $\alpha > 0$ ensures that $det(\mathbf{H}_2) > 0$ and, thus, matrix $\mathbf{H}_G J$ is positive semidefinite. This proves that the desired stationary points in J are local minima.

III. REPLY TO COMMENT 2

The second comment refers to the statement in [2] that all undesired stationary points in group 6 are not local minima. Gu *et al.* provide a counterexample to show that this statement is not true. The counterexample shows that, if $\alpha = k_1 = k_2 = 1$, the cost function presents an undesirable local minimum at the point

$$g_{11} = \frac{\sqrt{2}}{10}$$
 $g_{12} = \frac{\sqrt{582}}{30}$ $g_{21} = -g_{11}$ $g_{22} = g_{12}$ (6)

because its Hessian matrix is positive semidefinite. Thus, it is erroneous that the same conditions $\alpha > 0$, $k_1 < 2$, and $k_2 < 2$ that ensured local convergence to the desired stationary points also guarantee that the cost function J is not affected by undesired local minima. However, the counterexample does not exclude the possibility that other conditions on α , k_1 and k_2 ensure the inexistence of undesired local minima. It is important, for instance, to highlight that when $k_1 = k_2 = 1$, $g_{21} = -g_{11}$, and $g_{22} = g_{12}$, the first derivatives of the cost function take the form

$$\frac{\partial J}{\partial g_{11}} = -\frac{\partial J}{\partial g_{21}}$$

= $g_{11}^* \left(2|g_{11}|^2 + 4|g_{12}|^2 - 2 + \alpha|g_{11}|^2 \right) - \alpha|g_{12}|^2 g_{11}^*$ (7)

$$\frac{\partial J}{\partial g_{12}} = \frac{\partial J}{\partial g_{22}} = g_{12}^* (2|g_{12}|^2 + 4|g_{11}|^2 - 2 + \alpha|g_{12}|^2) - \alpha|g_{11}|^2 g_{12}^*.$$
(8)

It is straightforward to show that these two equations vanish only when $\alpha = 1$. As a consequence, the convergence to the undesirable stationary point given by (6) can be avoided by simply selecting $\alpha \neq 1$. Unfortunately, the existence of other undesirable minima in group 6 is still an open question and a more detailed analysis should be carried out in order to provide a definitive answer.

Nevertheless, we want to recall that the analysis of the stationary points in J is rather cumbersome and for this reason we only considered two sources and a two-output separating system. We did not succeed in extending the analysis to the general case of N sources. On the other hand, we investigated in [3] a more general family of cost functions for blind source separation ψ which utilizes the Shalvi and Weinstein criterion¹ for blind equalization [4] and fourth-order cross cumulants, instead of cross correlations. Both terms are related through a weighting parameter that we denoted by γ . The stationary point analysis of ψ is carried out in [3] for the general case of N complex-valued sources. If all the sources are either sub-Gaussian or super-Gaussian, we show that the condition $\gamma > 0$ is sufficient to guarantee that the desired stationary points where each output extracts a single source are local minima. In addition, the condition $\gamma > 1$ is sufficient to ensure that the cost function ψ does not contain any undesired local minima. Comment 2 does not apply to ψ since no similar counterexample can be provided to refute that it does not contain undesired local minima. The pitfall in the convergence analysis of [2] does not occur in [3].

¹Note that this criterion contains CM as a particular case.

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Unconditional Maximum Likelihood Performance at Finite Number of Samples and High Signal-to-Noise Ratio

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Abstract-This correspondence deals with the problem of estimating signal parameters using an array of sensors. In source localization, two main maximum-likelihood methods have been introduced: the conditional maximum-likelihood method which assumes the source signals nonrandom and the unconditional maximum-likelihood method which assumes the source signals random. Many theoretical investigations have been already conducted for the large samples statistical properties. This correspondence studies the behavior of unconditional maximum likelihood at high signal-to-noise ratio for finite samples. We first establish the equivalence between the unconditional and the conditional maximum-likelihood criterions at high signal-to-noise ratio. Then, thanks to this equivalence we prove the non-Gaussianity and the non-efficiency of the unconditional maximum-likelihood estimator. We also rediscover the closed-form expressions of the probability density function and of the variance of the estimates in the one source scenario and we derive a closed-form expression of this estimator variance in the two sources scenario.

Index Terms—Asymptotic performance, Cramér-Rao bound, finite number of data, high signal-to-noise ratio, unconditional maximum likelihood.

I. INTRODUCTION

Direction-of-arrival (DOA) estimation using an array of spatially distributed sensors has received a significant attention in the signal processing literature. Initial motivation was the military framework with applications such as radar and sonar. More recently, DOA estimation has also been applied to other frameworks such as friendly communication. For these numerous applications, the resolving power of the algorithm is of the utmost importance. This is why various algorithms have been proposed in the literature with a resolution which is better than

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the traditional Rayleigh beam-width [1]–[4]. An alternative to these algorithms is the maximum-likelihood (ML) method which has been extensively studied for its attractive statistical properties. When applying the ML technique to the sensor array problem, two main methods have been considered, depending on the model assumption on the signal waveforms. When the source signals are modeled as Gaussian random processes, an unconditional ML (UML) is obtained (see [5]–[7]). If, on the other hand, when the source signals are modeled as unknown deterministic quantities, the resulting estimator is referred as the conditional ML (CML) estimator (see [7]–[9]).

This correspondence deals with the asymptotic performance of the UML method. The term "asymptotic" can be understood in two different ways: in the number T of samples and in the signal-to-noise ratio (SNR). Asymptotic performance in the number T of samples (for finite SNR) have been extensively investigated [7], [10]-[12]. Concerning the asymptotic performance when the SNR tends to infinity (for finite T), few works are available. Under the deterministic model, the CML is Gaussian and efficient (it achieves the conditional Cramér-Rao bound) [13], [14]. The present work is devoted to the analysis of the UML behavior, under the stochastic signals model, when the SNR tends to infinity (for finite T): this is the meaning of asymptotic in this correspondence. Note that in [15], Athley has observed, with the help of simulation results, that the UML estimates are nonefficient at high SNR. The proposed paper aims to soundly establish the asymptotic non-Gaussianity and the asymptotic nonefficiency (in comparison with the unconditional Cramér-Rao bound) of the UML estimator in the multiple parameters case.

We have already investigated the UML asymptotic behavior for a single source [16]. The proposed paper generalizes these preliminary results to multiple sources case, providing an extended and detailed version of works reported in conference papers [17] and [18]. We first show that, at high SNR, unconditional and conditional maximum-like-lihood criterions (UMLC and CMLC) are equivalent in the sense that, with the same observations, they give the same estimates. This preliminary result is the key point for proving that the UML estimates are non-Gaussian and non-efficient when the SNR tends to infinity for any number of sources contrary to the large number of observations case. Finally, we establish a closed-form of the UML estimator variance in the case of two uncorrelated sources for centro-symmetric arrays.

In the sequel, a sample of a random vector \mathbf{y} is denoted $\mathbf{y}(\omega)$, where ω belongs to the event space Ω .

II. PROBLEM SETUP

Let us consider the classical problem of localizing N narrow-band sources impinging on an array of M sensors. The vector $\mathbf{x}_t(\omega)$ of sensors outputs is given by the following equation [9]:

$$\mathbf{x}_{t}(\omega) = \mathbf{A}(\boldsymbol{\theta}_{0})\mathbf{s}_{t}(\omega) + \mathbf{n}_{t}(\omega)$$
(1)

where t = 1, 2, ..., T and where T is the number of snapshots. $\boldsymbol{\theta} = [\theta_1, \theta_2, ..., \theta_N]^T$ denotes the candidate vector of the N DOAs whose exact value is $\boldsymbol{\theta}_0$. $\mathbf{A}(\boldsymbol{\theta}) = [\mathbf{a}(\theta_1), \mathbf{a}(\theta_2), ..., \mathbf{a}(\theta_N)]$ is the $M \times N$ steering matrix. $\mathbf{s}_t(\omega)$ is the $N \times 1$ vector of the N source signals. $\mathbf{n}_t(\omega)$ is the $M \times 1$ vector of the noise.

In the sequel $\mathbf{N}(\omega) = [\mathbf{n}_1(\omega), \mathbf{n}_2(\omega), \dots, \mathbf{n}_T(\omega)],$ and $\mathbf{S}(\omega) = [\mathbf{s}_1(\omega), \mathbf{s}_2(\omega), \dots, \mathbf{s}_T(\omega)].$

The following assumptions will be used:

A1) The signal $\mathbf{s}_t(\omega)$ is the sample of the random vector \mathbf{s}_t which is complex, circular, Gaussian, temporally white with zero mean, and covariance matrix $\Sigma_{\mathbf{s}} = \mathbb{E}[\mathbf{ss}^H]$ where \mathbb{E} denotes the expectation operator.

A2) The noise $\mathbf{n}_t(\omega)$ is the sample of the random vector \mathbf{n}_t which is complex, circular, Gaussian, spatially and temporally white with zero mean, and covariance matrix $\boldsymbol{\Sigma}_{\mathbf{n}} = \mathbb{E}[\mathbf{n}\mathbf{n}^H] = \sigma^2 \mathbf{I}_M$ where \mathbf{I}_M is the $M \times M$ identity matrix.

A3)
$$\|\mathbf{a}(\theta)\| = \sqrt{M}.$$

A4) The number of sources is less than the number of sensors, M > N.

Note that the model used in A1 differs from the conditional model, for which the signal s_t is deterministic.

III. HIGH SNR EQUIVALENCE OF THE CONDITIONAL CRITERION AND UNCONDITIONAL CRITERION

In this section, we recall the definition of the CMLC and of the UMLC and we prove the equivalence of these two criterions at high SNR in the sense where, with the same observations, they lead to the same estimates.

A. Conditional and Unconditional Maximum-Likelihood Criterion

In the conditional model case, the DOAs are obtained by minimization of the concentrated criterion [9]:

$$C_{\text{CML}}(\boldsymbol{\theta}) = \frac{1}{M-N} Tr\left\{ \boldsymbol{\Pi}_{\mathbf{A}}^{\perp}(\boldsymbol{\theta}) \, \widehat{\boldsymbol{\Sigma}}_{\mathbf{x}} \right\}$$
(2)

where $Tr \{.\}$ is the trace operator, where $\widehat{\Sigma}_{\mathbf{x}} = (1/T) \sum_{t=1}^{T} \mathbf{x}_t(\omega) \mathbf{x}_t^H(\omega)$ is the observations sample covariance matrix, and $\mathbf{\Pi}_{\mathbf{A}}^{\perp}(\boldsymbol{\theta}) = \mathbf{I}_M - \mathbf{A}(\boldsymbol{\theta}) (\mathbf{A}^H(\boldsymbol{\theta}) \mathbf{A}(\boldsymbol{\theta}))^{-1} \mathbf{A}^H(\boldsymbol{\theta})$ denotes the orthogonal projector onto the noise subspace. In the sequel, the Moore–Penrose inverse $(\mathbf{A}^H(\boldsymbol{\theta}) \mathbf{A}(\boldsymbol{\theta}))^{-1} \mathbf{A}^H(\boldsymbol{\theta})$, where $\mathbf{A}(\boldsymbol{\theta})$ is a full-column rank matrix, will be denoted $\mathbf{A}^{\dagger}(\boldsymbol{\theta})$.

In the stochastic model case, the DOAs are obtained by minimization of the concentrated criterion [9]:

$$C_{\text{UML}}(\boldsymbol{\theta}) = \left| \mathbf{A}(\boldsymbol{\theta}) \, \widehat{\mathbf{R}}_{s} \mathbf{A}^{H}(\boldsymbol{\theta}) + \hat{\sigma}^{2} \mathbf{I}_{M} \right|$$
(3)

with

$$\begin{cases} \widehat{\mathbf{R}}_{\mathbf{s}}(\boldsymbol{\theta}) = \mathbf{A}^{\dagger}(\boldsymbol{\theta}) \left(\widehat{\mathbf{\Sigma}}_{\mathbf{x}} - \hat{\sigma}^{2}(\boldsymbol{\theta}) \mathbf{I}_{M} \right) \mathbf{A}^{\dagger H}(\boldsymbol{\theta}), \\ \hat{\sigma}^{2}(\boldsymbol{\theta}) = \frac{1}{M-N} Tr \left\{ \mathbf{\Pi}_{\mathbf{A}}^{\perp}(\boldsymbol{\theta}) \, \widehat{\mathbf{\Sigma}}_{\mathbf{x}} \right\}, \end{cases}$$
(4)

where |.| denotes the determinant.

By substituting (2) and (4) into (3), we straightforwardly obtain

$$C_{\text{UML}}(\boldsymbol{\theta}) = \left| \boldsymbol{\Pi}_{\mathbf{A}}(\boldsymbol{\theta}) \, \widehat{\boldsymbol{\Sigma}}_{\mathbf{x}} \boldsymbol{\Pi}_{\mathbf{A}}(\boldsymbol{\theta}) + C_{\text{CML}}(\boldsymbol{\theta}) \, \boldsymbol{\Pi}_{\mathbf{A}}^{\perp}(\boldsymbol{\theta}) \right| \qquad (5)$$

where $\Pi_{\mathbf{A}}(\boldsymbol{\theta}) = \mathbf{A}(\boldsymbol{\theta})\mathbf{A}^{\dagger}(\boldsymbol{\theta})$ denotes the orthogonal projector onto the signal subspace.

B. Equivalence

Proposition 1: At high SNR, the UMLC and the CMLC are equivalent in the sense where the difference of DOAs obtained by minimization of $C_{\rm UML}$ (θ) and $C_{\rm CML}$ (θ) tends to zero in probability when SNR tends to infinity.

Elements of Proof: Let $\mathbf{E}_{\mathbf{s}}(\boldsymbol{\theta})$ and $\mathbf{E}_{\mathbf{n}}(\boldsymbol{\theta})$ be the $M \times N$ and $M \times (M - N)$ matrices built with the orthonormal bases of signal and noise subspaces and set $\mathbf{E}(\boldsymbol{\theta})$ be the $M \times M$ matrix such that $\mathbf{E}(\boldsymbol{\theta}) = [\mathbf{E}_{\mathbf{s}}(\boldsymbol{\theta}), \mathbf{E}_{\mathbf{n}}(\boldsymbol{\theta})]$. Equation (5) becomes

$$C_{\text{UML}}(\boldsymbol{\theta}) = \left| \mathbf{E}^{H}(\boldsymbol{\theta}) \left(\mathbf{\Pi}_{\mathbf{A}}(\boldsymbol{\theta}) \, \widehat{\mathbf{\Sigma}}_{\mathbf{x}} \mathbf{\Pi}_{\mathbf{A}}(\boldsymbol{\theta}) + C_{\text{CML}}(\boldsymbol{\theta}) \, \mathbf{\Pi}_{\mathbf{A}}^{\perp}(\boldsymbol{\theta}) \right) \mathbf{E}(\boldsymbol{\theta}) \right|$$
$$= \left| \begin{array}{c} \mathbf{E}_{\mathbf{s}}^{H}(\boldsymbol{\theta}) \, \widehat{\mathbf{\Sigma}}_{\mathbf{x}} \mathbf{E}_{\mathbf{s}}(\boldsymbol{\theta}) & \mathbf{0} \\ \mathbf{0} & C_{\text{CML}}(\boldsymbol{\theta}) \, \mathbf{I}_{M-N} \end{array} \right|.$$
(6)

The matrix involved in the determinant (6) is block diagonal so that $C_{\text{UML}}(\boldsymbol{\theta})$ can also be written as follows by writing down explicitly the dependance of each terms on the noise and $\boldsymbol{\theta}$

$$C_{\text{UML}}\left(\boldsymbol{\theta}, \mathbf{N}\left(\boldsymbol{\omega}\right)\right) = \left|\mathbf{E}_{\mathbf{s}}^{H}(\boldsymbol{\theta}) \, \widehat{\mathbf{\Sigma}}_{\mathbf{x}}\left(\boldsymbol{\theta}, \mathbf{N}(\boldsymbol{\omega})\right) \mathbf{E}_{\mathbf{s}}(\boldsymbol{\theta})\right| C_{\text{CML}}(\boldsymbol{\theta}, \mathbf{N}\left(\boldsymbol{\omega}\right))^{M-N}.$$
 (7)

Note that the minimization of $C_{\text{UML}}(\boldsymbol{\theta}, \mathbf{N}(\omega))$ is equivalent to the minimization of $\tilde{C}_{\text{UML}}(\boldsymbol{\theta}, \mathbf{N}(\omega)) = (C_{\text{UML}}(\boldsymbol{\theta}, \mathbf{N}(\omega)))^{1/(M-N)}$, consequently we will study

$$\hat{C}_{\text{UML}}\left(\boldsymbol{\theta}, \mathbf{N}\left(\boldsymbol{\omega}\right)\right) = \left|\mathbf{E}_{\mathbf{s}}^{H}(\boldsymbol{\theta}) \,\widehat{\boldsymbol{\Sigma}}_{\mathbf{x}}(\boldsymbol{\theta}, \mathbf{N}(\boldsymbol{\omega})) \,\mathbf{E}_{\mathbf{s}}(\boldsymbol{\theta})\right|^{1/(M-N)} C_{\text{CML}}(\boldsymbol{\theta}, \mathbf{N}\left(\boldsymbol{\omega}\right)). \quad (8)$$

The right-hand side of (8) is the product of two terms. Let us set

$$\alpha\left(\boldsymbol{\theta},\mathbf{N}\left(\omega\right)\right) = \left|\mathbf{E}_{\mathbf{s}}^{H}\left(\boldsymbol{\theta}\right)\widehat{\mathbf{\Sigma}}_{\mathbf{x}}\left(\boldsymbol{\theta},\mathbf{N}\left(\omega\right)\right)\mathbf{E}_{\mathbf{s}}\left(\boldsymbol{\theta}\right)\right|^{1/(M-N)}.$$
 (9)

A Taylor expansion at order zero around $(\boldsymbol{\theta}_{0}, \mathbf{0})$ of $\alpha(\boldsymbol{\theta}, \mathbf{N}(\omega))$ leads to

$$\alpha\left(\boldsymbol{\theta}, \mathbf{N}\left(\omega\right)\right) = \alpha\left(\boldsymbol{\theta}_{0}, \mathbf{0}\right) + o\left(1\right) \tag{10}$$

where o denotes the small oh notation and where

$$\alpha\left(\boldsymbol{\theta}_{0},\boldsymbol{0}\right)=\left|\mathbf{E}_{\mathbf{s}}^{H}\left(\boldsymbol{\theta}_{0}\right)\mathbf{A}\left(\boldsymbol{\theta}_{0}\right)\widehat{\boldsymbol{\Sigma}}_{\mathbf{s}}\mathbf{A}^{H}\left(\boldsymbol{\theta}_{0}\right)\mathbf{E}_{\mathbf{s}}\left(\boldsymbol{\theta}_{0}\right)\right|\neq0$$
(11)

where $\widehat{\Sigma}_{s} = (1/T) \sum_{t=1}^{T} \mathbf{s}_{t}(\omega) \mathbf{s}_{t}^{H}(\omega)$. Consequently, the first nonnull term of a Taylor expansion of $\alpha(\boldsymbol{\theta}, \mathbf{n})$ around $\boldsymbol{\theta} = \boldsymbol{\theta}_{0}$ and $\mathbf{N}(\omega) = \mathbf{0}$ is $\alpha(\boldsymbol{\theta}_{0}, \mathbf{0})$. Concerning the term $C_{\text{CML}}(\boldsymbol{\theta}, \mathbf{N}(\omega))$, a Taylor expansion at order two around $(\boldsymbol{\theta}_{0}, \mathbf{0})$ leads to

$$C_{\text{CML}}\left(\boldsymbol{\theta}, \mathbf{N}\left(\omega\right)\right) = C_{\text{CML}}\left(\boldsymbol{\theta}_{0}, \mathbf{0}\right) + \boldsymbol{\Delta}^{T}\mathbf{G} + \frac{1}{2}\boldsymbol{\Delta}^{T}\ddot{\mathbf{H}}\boldsymbol{\Delta} + o\left(\left\|\boldsymbol{\Delta}\right\|^{2}\right) \quad (12)$$

where ||.|| denotes the norm, where

$$\boldsymbol{\Delta} = \left[\left(\boldsymbol{\theta} - \boldsymbol{\theta}_0\right)^T, vec \left(\operatorname{Re} \left\{ \mathbf{N} \left(\omega \right) \right\} \right)^T, vec \left(\operatorname{Im} \left\{ \mathbf{N} \left(\omega \right) \right\} \right)^T \right]^T$$
(13)

where Re {} and Im {} denotes the real and imaginary part, respectively, and where *vec* denotes the vec operator. **G** is the gradient of $C_{\text{CML}}(\boldsymbol{\theta}, \mathbf{N}(\omega))$ at $(\boldsymbol{\theta}_0, \mathbf{0})$

$$\mathbf{G} = \left[\left(\frac{\partial C_{\text{CML}} \left(\boldsymbol{\theta}, \mathbf{N} \left(\boldsymbol{\omega} \right) \right)}{\partial \boldsymbol{\theta}} \Big|_{\boldsymbol{\theta}_{0}, \mathbf{0}} \right)^{T} \\ \left(\frac{\partial C_{\text{CML}} \left(\boldsymbol{\theta}, \mathbf{N} \left(\boldsymbol{\omega} \right) \right)}{\partial vec \left(\text{Re} \left\{ \mathbf{N} \left(\boldsymbol{\omega} \right) \right\} \right)} \Big|_{\boldsymbol{\theta}_{0}, \mathbf{0}} \right)^{T} \\ \left(\frac{\partial C_{\text{CML}} \left(\boldsymbol{\theta}, \mathbf{N} \left(\boldsymbol{\omega} \right) \right)}{\partial vec \left(\text{Im} \left\{ \mathbf{N} \left(\boldsymbol{\omega} \right) \right\} \right)} \Big|_{\boldsymbol{\theta}_{0}, \mathbf{0}} \right)^{T} \right]^{T}$$
(14)

and $\mathbf{\hat{H}}$ is the Hessian matrix of $C_{\text{CML}}(\boldsymbol{\theta}, \mathbf{N}(\omega))$ at $(\boldsymbol{\theta}_0, \mathbf{0})$, as shown by (15) at the bottom of the page.

For $(\boldsymbol{\theta}, \mathbf{n}) = (\boldsymbol{\theta}_0, \mathbf{0}), C_{\text{CML}}(\boldsymbol{\theta}, \mathbf{N}(\omega))$ is minimal and null. Consequently

$$C_{\text{CML}}\left(\boldsymbol{\theta}_{0}, \mathbf{0}\right) = 0 \text{ and } \mathbf{G} = \mathbf{0}$$
(16)

and

$$C_{\text{CML}}\left(\boldsymbol{\theta}, \mathbf{N}\left(\omega\right)\right) = \frac{1}{2} \boldsymbol{\Delta}^{T} \ddot{\mathbf{H}} \boldsymbol{\Delta} + o\left(\left\|\boldsymbol{\Delta}\right\|^{2}\right).$$
(17)

Therefore, the first non-null term of its Taylor expansion around $\boldsymbol{\theta} = \boldsymbol{\theta}_0$ and $\mathbf{N}(\omega) = \mathbf{0}$ is $(1/2)\Delta^T \mathbf{H} \boldsymbol{\Delta}$. Consequently

$$\tilde{C}_{\text{UML}}\left(\boldsymbol{\theta}, \mathbf{N}\left(\omega\right)\right) = \frac{1}{2} \alpha\left(\boldsymbol{\theta}_{0}, \mathbf{0}\right) \boldsymbol{\Delta}^{T} \ddot{\mathbf{H}} \boldsymbol{\Delta} + o\left(\|\boldsymbol{\Delta}\|^{2}\right)$$
$$= \alpha\left(\boldsymbol{\theta}_{0}, \mathbf{0}\right) C_{\text{CML}}\left(\boldsymbol{\theta}, \mathbf{N}\left(\omega\right)\right) + o\left(\|\boldsymbol{\Delta}\|^{2}\right).$$
(18)

Consequently, at high SNR, since $\tilde{C}_{\text{UML}}(\boldsymbol{\theta}, \mathbf{N}(\omega))$ is the product of $C_{\text{CML}}(\boldsymbol{\theta}, \mathbf{N}(\omega))$ by a non-null constant, both criterions provide the same estimates, concluding the proof.

IV. NON-GAUSSIANITY OF THE UML

In the sequel, concerning source signals, we are in the stochastic model framework of assumption A1 and we note $\hat{\theta}$ = arg min $C_{\text{UML}}(\theta)$ the UML estimator. The next theorem establishes

the asymptotic distribution of $\hat{\theta}$ and its non-Gaussianity for any number of sources (the single source case has already been reported in [16])

Theorem 1: Let $\tilde{\boldsymbol{\theta}} = (1/\sigma) \left(\hat{\boldsymbol{\theta}} - \boldsymbol{\theta}_0 \right)$. When SNR tends to infinity, $\tilde{\boldsymbol{\theta}}$ is non-Gaussian and converges in distribution to $\mathbf{C}(\boldsymbol{\theta}_0)\mathbf{y}$, where \mathbf{y} is a $N \times 1$ Gaussian vector with zero mean and covariance matrix \mathbf{I}_N and $\mathbf{C}(\boldsymbol{\theta}_0)$ is any $N \times N$ random matrix independent of vector \mathbf{y} , satisfying

$$\mathbf{C}(\boldsymbol{\theta}_{0})\mathbf{C}^{T}(\boldsymbol{\theta}_{0}) = \frac{1}{2T} \left(\operatorname{Re}\left\{ \mathbf{H}(\boldsymbol{\theta}_{0}) \odot \mathbf{\hat{\Sigma}}_{\mathbf{s}}^{T} \right\} \right)^{-1}$$
(19)

where \odot denotes the Hadamard product (element by element product), and where $\mathbf{H}(\boldsymbol{\theta}_0)$ is a $N \times N$ deterministic matrix which contains the information about the DOA's and about the array structure

$$\mathbf{H}(\boldsymbol{\theta}_{0}) = \mathbf{D}^{H}(\boldsymbol{\theta}_{0}) \mathbf{\Pi}_{\mathbf{A}}^{\perp}(\boldsymbol{\theta}) \mathbf{D}(\boldsymbol{\theta}_{0})$$
(20)

with

$$\mathbf{D}\left(\boldsymbol{\theta}_{0}\right) = \left[\left.\frac{d\mathbf{a}\left(\boldsymbol{\theta}\right)}{d\boldsymbol{\theta}}\right|_{\boldsymbol{\theta}_{1}}, \left.\frac{d\mathbf{a}\left(\boldsymbol{\theta}\right)}{d\boldsymbol{\theta}}\right|_{\boldsymbol{\theta}_{2}}, \dots, \left.\frac{d\mathbf{a}\left(\boldsymbol{\theta}\right)}{d\boldsymbol{\theta}}\right|_{\boldsymbol{\theta}_{N}}\right].$$
(21)

Note that in (19), $T\hat{\Sigma}_{s}$ is a $N \times N$ random matrix which follows a complex Wishart distribution with T degrees of freedom and parameter matrix the covariance Σ_{s} of source signals s_{t} .

Proof: From proposition 1, we consider that $\hat{\boldsymbol{\theta}}$ is obtained by minimization of $C_{\text{CML}}(\boldsymbol{\theta})$ given by (2). Thanks to [14], at high SNR, the conditional distribution $f\left(\tilde{\boldsymbol{\theta}} \mid \mathbf{S}\right)$ is Gaussian with asymptotic covariance given by the conditional Cramér–Rao bound, see [11]

$$\mathbf{B}_{\text{COND}}\left(\boldsymbol{\theta}_{0}\right) = \frac{1}{2T} \left(\operatorname{Re}\left\{ \mathbf{H}\left(\boldsymbol{\theta}_{0}\right) \odot \widehat{\boldsymbol{\Sigma}}_{\mathbf{s}}^{T} \right\} \right)^{-1}.$$
(22)

Let us set $\mathbf{B}_{\text{COND}}(\boldsymbol{\theta}_0) = \mathbf{C}(\boldsymbol{\theta}_0)\mathbf{C}^T(\boldsymbol{\theta}_0)$. Therefore, the asymptotic (in SNR) conditional distribution $f\left(\widetilde{\boldsymbol{\theta}} \mid \mathbf{S}\right)$ is the same as the distribution of $\mathbf{C}(\boldsymbol{\theta}_0)\mathbf{y}$, where \mathbf{y} is a Gaussian random vector with zero



mean and covariance matrix \mathbf{I}_N and where $\mathbf{C}(\boldsymbol{\theta}_0)$ is a deterministic matrix. Consequently, the asymptotic (in SNR) marginal distribution $f\left(\widetilde{\boldsymbol{\theta}}\right)$ is the same as that of $\mathbf{C}(\boldsymbol{\theta}_0)\mathbf{y}$ where \mathbf{y} is a Gaussian random vector with zero mean and covariance matrix \mathbf{I}_N and where $\mathbf{C}(\boldsymbol{\theta}_0)$ becomes a random matrix since, in (19), $T\widehat{\boldsymbol{\Sigma}}_s$ is complex Wishart distributed with T degrees of freedom, and parameter matrix $\boldsymbol{\Sigma}_s$ the source signals covariance. Since $\mathbf{C}(\boldsymbol{\theta}_0)$ becomes a random matrix, the product $\mathbf{C}(\boldsymbol{\theta}_0)\mathbf{y}$ cannot be Gaussian which completes the proof.

V. NONEFFICIENCY OF THE UML ESTIMATOR

In order to proof the nonefficiency of the UML estimator, the comparison between the asymptotic covariance of the UML estimator and the Unconditional Cramér–Rao Bound (UCRB) is provided in this section.

A. Asymptotic Covariance of $\hat{\theta}$

Corollary 1: Let $cov\left(\tilde{\boldsymbol{\theta}}\right) = \mathbb{E}\left[\tilde{\boldsymbol{\theta}}\tilde{\boldsymbol{\theta}}^{T}\right]$ be the covariance of $\tilde{\boldsymbol{\theta}}$. Then, from the above section, we have straightforwardly

$$\lim_{\sigma \to 0} cov\left(\widetilde{\boldsymbol{\theta}}\right) = \mathbb{E}\left[\mathbf{C}\mathbf{y}\mathbf{y}^{T}\mathbf{C}^{T}\right] = \mathbb{E}[\mathbf{C}\mathbf{C}^{T}],$$
$$= \frac{1}{2T} \mathbb{E}\left[\left(\operatorname{Re}\left\{\mathbf{H}\left(\boldsymbol{\theta}_{0}\right) \odot \widehat{\boldsymbol{\Sigma}}_{\mathbf{s}}^{T}\right\}\right)^{-1}\right].$$
(23)

B. Performance Bound

According to [11], the UCRB can be written as follows:

$$\mathbf{B}_{\text{UCOND}}\left(\boldsymbol{\theta}_{0}\right) = \frac{\sigma^{2}}{2T} \left(\operatorname{Re}\left\{ \mathbf{H}\left(\boldsymbol{\theta}_{0}\right) \odot \left(\boldsymbol{\Sigma}_{\mathbf{s}} \mathbf{A}^{H}\left(\boldsymbol{\theta}_{0}\right) \boldsymbol{\Sigma}_{\mathbf{x}}^{-1} \mathbf{A}\left(\boldsymbol{\theta}_{0}\right) \boldsymbol{\Sigma}_{\mathbf{s}}\right)^{T} \right\} \right)^{-1}$$
(24)

where $\Sigma_{\mathbf{x}}$ is the covariance matrix of the observations.

By using relation [11, (3.20)], it is shown that in (24)

$$\mathbf{A}^{H}(\boldsymbol{\theta}_{0})\boldsymbol{\Sigma}_{\mathbf{x}}^{-1}\mathbf{A}(\boldsymbol{\theta}_{0}) = \left(\boldsymbol{\Sigma}_{\mathbf{s}} + \sigma^{2}\left(\mathbf{A}^{H}(\boldsymbol{\theta}_{0})\mathbf{A}(\boldsymbol{\theta}_{0})\right)^{-1}\right)^{-1}$$
(25)

which tends to Σ_s^{-1} when σ tends to 0. It follows

$$\lim_{\sigma \to 0} \frac{\mathbf{B}_{\text{UCOND}}\left(\boldsymbol{\theta}_{0}\right)}{\sigma^{2}} = \frac{1}{2T} \left(\operatorname{Re}\left\{ \mathbf{H}\left(\boldsymbol{\theta}_{0}\right) \odot \boldsymbol{\Sigma}_{\mathbf{s}}^{T} \right\} \right)^{-1}.$$
(26)

C. Nonefficiency of the UML Estimator

In order to prove the nonefficiency of the UML estimator for any number of sources, the following theorem will be of interest. Note that this is a matrix extension of the well-known Jensen's inequality. This theorem has been proved in [19] without the equality condition which will be of particular interest here.

Theorem 2: Let Θ be a $N \times N$ real positive definite random matrix. Then

$$\mathbb{E}\left[\boldsymbol{\Theta}^{-1}\right] \ge \left(\mathbb{E}\left[\boldsymbol{\Theta}\right]\right)^{-1} \tag{27}$$

with equality if and only if Θ is a constant matrix with probability one. Appendix A details the proof.

Corollary 2: Let us set $\boldsymbol{\Theta} = 2T \operatorname{Re} \left\{ \mathbf{H}(\boldsymbol{\theta}_0) \odot \hat{\boldsymbol{\Sigma}}_{\mathbf{s}}^T \right\}$ in (26) and (23). Equation (27) becomes the Cramér–Rao inequality. Since $\boldsymbol{\Theta}$ is not a constant matrix with probability one, the inequality is strict and the UML estimator is nonefficient for any number of sources.

VI. SPECIFIC PERFORMANCE STUDY OF THE UML ESTIMATOR FOR THE TWO SOURCES SCENARIO

This section is devoted to a deeper statistical investigation of two specific cases frequently met in array processing: the single– and twosources case. We remind the probability density function (pdf) and the variance closed form of the UML estimates in the single-source case tediously obtained in [16]. For two uncorrelated sources and centro-symmetric arrays, we give a closed-form expression of the UML estimates covariance.

A. Distribution and Theoretical Variance in the Single Source Case

In the single-source case, $\Sigma_s = \Sigma_1$ and $\mathbf{H} = h_1$. Then $\tilde{\theta}$ is asymptotically distributed as $\sqrt{k}S_{2T}$, where S_{2T} is a Student random variable with 2*T* degrees of freedom and *k* is given by

$$k = \lim_{\sigma \to 0} \frac{\mathbf{B}_{\text{UCOND}}\left(\boldsymbol{\theta}_{0}\right)}{\sigma^{2}} = \frac{1}{2Th_{1}\Sigma_{1}}.$$
 (28)

The asymptotic variance of $\tilde{\theta}$ is then given by

$$var(\tilde{\theta}) = \frac{T}{T-1}k.$$
(29)

As established in Theorem 2, for finite T, the UML estimator is not asymptotically efficient since T/(T-1) > 1.

B. Theoretical Variance in the Two Sources Case for Uncorrelated Sources and Centro-Symmetric Array

Most arrays met in practice possess a center of symmetry (this is for instance the case of the ULA). Under this condition which will be assumed in the following, the matrix \mathbf{H} of (20) is real and symmetric (see Appendix B):

$$\mathbf{H}\left(\boldsymbol{\theta}_{0}\right) = \begin{pmatrix} h_{1} & h_{3} \\ h_{3} & h_{2} \end{pmatrix}.$$
(30)

For two uncorrelated sources, $\Sigma_s = Diag \{\Sigma_1, \Sigma_2\}$ and the asymptotic covariance of $\tilde{\theta}$ is given by

$$\lim_{\sigma \to 0} \cos\left(\tilde{\boldsymbol{\theta}}\right) = \frac{1}{2} \frac{{}^{2}F_{1}\left(1, 1; 2T; \frac{h_{3}^{2}}{h_{1}h_{2}}\right)}{T-1} \operatorname{Diag}\left\{\frac{1}{h_{1}\Sigma_{1}}, \frac{1}{h_{2}\Sigma_{2}}\right\},\\ = \frac{T}{T-1} {}_{2}F_{1}\left(1, 1; 2T; \frac{h_{3}^{2}}{h_{1}h_{2}}\right) \mathbf{K}$$
(31)

where $\mathbf{K} = \lim_{\sigma \to 0} (\mathbf{B}_{\text{UCOND}} (\boldsymbol{\theta}_0) / \sigma^2)$ and where ${}_2F_1(a, b; c; \omega)$ is the Gauss hypergeometric function defined by its integral representation [20, pp. 558]

$${}_{2}F_{1}(a,b;c;\omega) = \frac{\Gamma(c)}{\Gamma(b)\Gamma(c-b)} \int_{0}^{1} z^{b-1} (1-z)^{c-b-1} (1-z\omega)^{-a} dz \quad (32)$$

where $\Gamma(z)$ denotes the Gamma function $\Gamma(z) = \int_0^\infty t^{z-1} e^{-t} dt$.

The derivation of (31) is given in Appendix C. As expected, the UML estimator is not asymptotically efficient since T/(T-1) > 1 and $_2F_1(1,1;2T;h_3^2/h_1h_2) \geq 1$.

VII. SIMULATION EXAMPLES

In this section, results of some Monte Carlo simulations concerning the UML estimator are presented. The purpose is to illustrate the applicability of the derived expressions of the pdf and of the variance. In all simulations, the array is an ULA of M = 10 sensors with half-wavelength spacing (the beamwidth of the array is equal to 10°). The UML DOA estimation is conducted with T = 2 snapshots. We consider the case of two uncorrelated sources with equal power located at 0° and 5° .



Fig. 1. Asymptotic variance of the UML estimator in the two-sources case. $\theta_0 = [0^\circ, 5^\circ], T = 2$ snapshots and M = 10 sensors.



Fig. 2. Histogram of UML estimates in the two-sources case. $\theta_0 = [0^\circ, 5^\circ]$, M = 10 sensors, T = 2 snapshots, and SNR = 30 dB.

DOA are given with respect to the broadside. The ML DOA estimation is performed with a Gauss–Newton algorithm thanks to a global search over a grid.

We have reported in Fig. 1 the evolution of the UML empirical variance, of the theoretical variance (31), and of the UCRB versus SNR. Monte Carlo simulations have been performed with r = 1000 independent realizations. Here, $(T/T - 1) {}_2F_1(1, 1; 2T; h_3^2/h_1h_2) = 2.9$. In this asymptotic region, one can notice the good match between theoretical results and simulations. The nonefficiency of UML at high SNR is observed. We also observe the well-known threshold effect [21] of the estimator variance when the SNR becomes weak (approximatively 20 dB in this case). This phenomena due to outliers gives the validity domain this asymptotic analysis (see [15] for more details concerning the UML threshold prediction).

Fig. 2 gives the histograms of the estimated DOA corresponding to the previous case with Monte Carlo simulations performed with r = 10000 independent realizations and a SNR of 30 dB. We also reported the pdf of a Gaussian distribution with the same variance. The non-Gaussianity of the UML estimates is observed. To confirm this

TABLE ICHI-SQUARE TEST IN THE TWO SOURCES CASE. k = 15 BINS, r = 10000REALIZATIONS. $\boldsymbol{\theta}_0 = [0^\circ, 5^\circ], M = 10$ SENSORS, T = 2 SNAPSHOTS,
AND SNR = 30 dB

Two sources case	Gaussian pdf: source 1	Gaussian pdf: source 2
Δ	$1.32 \ 10^{22}$	$4 \ 10^{24}$
$\Pr\left(X \ge \Delta\right)$	0%	0%
Hypothesis	rejected	rejected

"visual" result, we have used the classical Chi-square test which tests a distribution observed against another theoretical distribution. For the Chi-square fit computation, the data are divided into k = 15 bins and the statistical test requires the computation of

$$\Delta = \sum_{i=1}^{k} \frac{(O_i - rp_i)^2}{rp_i}$$
(33)

where O_i is the observed frequency for bin *i*, and p_i is the candidate probability for bin *i*. The hypothesis that the data are from a population following the candidate distribution is rejected if

$$\Pr\left(X \ge \Delta\right) = \frac{\Gamma\left(\frac{k-1}{2}, \frac{\Delta}{2}\right)}{\Gamma\left(\frac{k-1}{2}\right)} \le 5\%$$
(34)

where X follows a Chi-square distribution with k - 1 degrees of freedom. Table I shows that the pdf of the estimates is not Gaussian for a SNR of 30 dB.

VIII. CONCLUSION

The statistical properties of the UML estimator have been investigated. We have shown that, for any number of sources, this estimator is non-Gaussian and nonefficient at high signal-to-noise ratio for a finite number of samples. The key point of the analysis is the equivalence between the UML and the CML method at high SNR. Moreover, we have provided the UML estimator covariance closed-form expression for two uncorrelated sources and centro-symmetric array.

APPENDIX A PROOF OF THE UML NONEFFICIENCY

Lemma 1: Let Ω be a $N \times N$ real symmetric positive semidefinite matrix. Then $\forall \mathbf{q}$

$$\mathbf{q}^{T} \mathbf{\Omega} \mathbf{q} + \mathbf{q}^{T} \mathbf{\Omega}^{-1} \mathbf{q} - 2 \mathbf{q}^{T} \mathbf{q} \ge 0$$
(35)

with equality if and only if q is an eigenvector of Ω with eigenvalue one.

Proof: Let us set $\Omega = \sum_{i=1}^{N} \lambda_i \mathbf{r}_i \mathbf{r}_i^T$ the eigendecomposition of Ω on an orthonormal basis $\{\mathbf{r}_i\}_{i=1...N}$ with associated eigenvalues λ_i . Equation (35) can be written as

$$\sum_{i=1}^{N} \left(\lambda_{i} - 2 + \frac{1}{\lambda_{i}} \right) \left(\mathbf{q}^{T} \mathbf{r}_{i} \right)^{2} \ge 0.$$
(36)

Noticing that $\lambda_i - 2 + 1/\lambda_i \ge 0$ for $\lambda_i > 0$, and that $\lambda_i - 2 + 1/\lambda_i = 0$ for $\lambda_i = 1$ the proof is straightforward.

Lemma 2: Let Ω a $N \times N$ random real symmetric positive semidefinite matrix with probability one such that $\mathbb{E}[\Omega] = \mathbf{I}_N$. Then there is a vector \mathbf{q} such that

$$\mathbf{q}^{T} \mathsf{E} \left[\mathbf{\Omega}^{-1} \right] \mathbf{q} - \mathbf{q}^{T} \mathbf{q} > 0 \tag{37}$$

if and only if $\Pr[\mathbf{\Omega} = \mathbf{I}_N] \neq 1$.

Proof: Let us set $\zeta_{\mathbf{q}} = \mathbf{q}^T \Omega \mathbf{q} + \mathbf{q}^T \Omega^{-1} \mathbf{q} - 2\mathbf{q}^T \mathbf{q}$. Since $\mathbb{E}[\Omega] = \mathbf{I}_N$, we have $\mathbb{E}[\zeta_{\mathbf{q}}] = \mathbf{q}^T \mathbb{E}[\Omega^{-1}] \mathbf{q} - \mathbf{q}^T \mathbf{q}$. Consequently, proving Lemma 2 amounts to prove that $\exists \mathbf{q}$ such that $\mathbb{E}[\zeta_{\mathbf{q}}] > 0$ if and only

if $\Pr[\Omega = \mathbf{I}_N] \neq 1$. From (35) $\zeta_{\mathbf{q}}$ is a nonnegative random variable. Thus, $\mathbb{E}[\zeta_{\mathbf{q}}] > 0$ if and only if $\Pr[\zeta_{\mathbf{q}} = 0] \neq 1$. From lemma 1, $\Pr \left[\zeta_{\mathbf{q}} = 0 \right] = \Pr \left[\mathbf{\Omega} \mathbf{q} = \mathbf{q} \right].$ Consequently

$$\forall \mathbf{q} \Pr[\zeta_{\mathbf{q}} = 0] = 1 \iff \forall \mathbf{q} \Pr[\mathbf{\Omega} \mathbf{q} = \mathbf{q}] = 1 \iff \Pr[\mathbf{\Omega} = \mathbf{I}_N] = 1.$$

This completes the proof.

Finally, with the notations of Theorem 2, let us set Ω $\mathbb{E}[\Theta]^{-1/2} \Theta \mathbb{E}[\Theta]^{-1/2}$. Theorem 2 follows from Lemma 2.

APPENDIX B STUDY OF MATRIX H CENTRO-SYMMETRIC SENSOR ARRAYS

We prove in this appendix that \mathbf{H} [see equation (20)] is a real symmetric matrix. It is obvious that **H** is an hermitian matrix. Therefore, we must prove that **H** is a real matrix under the assumption that the array has a center of symmetry

$$\mathbf{H} \left(\boldsymbol{\theta}_{0} \right) = \mathbf{D}^{H} \left(\boldsymbol{\theta}_{0} \right) \mathbf{\Pi}_{\mathbf{A}}^{\perp} \left(\boldsymbol{\theta} \right) \mathbf{D} \left(\boldsymbol{\theta}_{0} \right)$$
$$= \mathbf{D}^{H} \left(\boldsymbol{\theta}_{0} \right) \mathbf{D} \left(\boldsymbol{\theta}_{0} \right) - \left(\mathbf{A}^{H} \left(\boldsymbol{\theta}_{0} \right) \mathbf{D} \left(\boldsymbol{\theta}_{0} \right) \right)^{H}$$
$$\cdot \left(\mathbf{A}^{H} \left(\boldsymbol{\theta}_{0} \right) \mathbf{A} \left(\boldsymbol{\theta}_{0} \right) \right)^{-1} \left(\mathbf{A}^{H} \left(\boldsymbol{\theta}_{0} \right) \mathbf{D} \left(\boldsymbol{\theta}_{0} \right) \right). \quad (38)$$

The *i*th element of the steering vector is^1

$$a_i(\theta_k) = e^{j(2\pi/\lambda)\mathbf{v}_i^T \mathbf{u}(\theta_k)}$$
(39)

where \mathbf{v}_i is the coordinate vector of the *i*th sensor, and $\mathbf{u}(\theta_k)$ is the unit vector pointing towards the kth source. Therefore, the ith element of $\mathbf{d}(\theta_k) = |d\mathbf{a}(\theta)/d\theta|_{\theta_k}$ [see (21)] is

$$d_{i}(\theta_{k}) = j \frac{2\pi}{\lambda} \mathbf{v}_{i}^{T} \underbrace{\frac{d\mathbf{u}}{d\theta}}_{\mathbf{u}(\theta_{k})} e^{j(2\pi/\lambda)\mathbf{v}_{i}^{T}\mathbf{u}(\theta_{k})}.$$
(40)

If the array has a center of symmetry, the sensors can be labeled so that $\mathbf{v}_i = -\mathbf{v}_{M-i+1}$. The *m*th row and *n*th column element of each term in (38) is as shown by (41) at the bottom of the page.

 $\mathbf{A}^{H}(\boldsymbol{\theta}_{0}) \mathbf{A}(\boldsymbol{\theta}_{0})\big|_{m,n} \text{ is a sum of two by two complex conjugates}$ with the same magnitude.² Therefore, $\mathbf{A}^{H}(\boldsymbol{\theta}_{0}) \mathbf{A}(\boldsymbol{\theta}_{0})\big|_{m,n} \in \mathbb{R}$. Sim-ilarly, $\mathbf{D}^{H}(\boldsymbol{\theta}_{0}) \mathbf{D}(\boldsymbol{\theta}_{0})\big|_{m,n} \in \mathbb{R}$ since $(\mathbf{v}_{i}^{T} \dot{\mathbf{u}}(\boldsymbol{\theta}_{m}))(\mathbf{v}_{i}^{T} \dot{\mathbf{u}}(\boldsymbol{\theta}_{n})) =$ $\left(\mathbf{v}_{M-i+1}^{T}\dot{\mathbf{u}}\left(\theta_{m}\right)\right)\left(\mathbf{v}_{M-i+1}^{T}\dot{\mathbf{u}}\left(\theta_{n}\right)\right)$, and $\left.\mathbf{A}^{H}\left(\boldsymbol{\theta}_{0}\right)\mathbf{D}\left(\boldsymbol{\theta}_{0}\right)\right|_{m,n}\in\mathbb{R}$, as shown by (42) at the bottom of the page.

Therefore, H is a real symmetric matrix.

 $^{1}\lambda$ is the wavelength of emitted signals.

²If the number of sensors is odd, $\mathbf{v}_{M/2+1} = 0$.

APPENDIX C

THEORETICAL VARIANCE IN THE TWO-SOURCES CASE

According to (23), and with the assumption that the array has a center of symmetry, i.e., $\mathbf{H}(\boldsymbol{\theta}_0)$ becomes a real symmetric matrix

$$\lim_{\sigma \to 0} cov\left(\widetilde{\boldsymbol{\theta}}\right) = \frac{1}{2T} \mathbb{E}\left[\left(\mathbf{H}\left(\boldsymbol{\theta}_{0}\right) \odot \operatorname{Re}\left(\widehat{\boldsymbol{\Sigma}}_{\mathbf{s}}^{T}\right)\right)^{-1}\right]$$
$$= \frac{1}{2} \mathbb{E}\left[\left(\mathbf{H}\left(\boldsymbol{\theta}_{0}\right) \odot \operatorname{Re}\left(\mathbf{W}\right)\right)^{-1}\right]$$
(43)

where \mathbf{W} is a $N \times N$ random matrix which follows a complex Wishart distribution with T degrees of freedom and parameter matrix the covariance, $\Sigma_{s} = \text{Diag} \{\Sigma_{1}, \Sigma_{2}\}$, of source signals s_{t} .

Under assumptions A1 and uncorrelated sources, $\mathbf{W}_R = \operatorname{Re} \{\mathbf{W}\}\$ is a $N \times N$ symmetric positive definite random matrix which follows a real Wishart distribution with 2T degrees of freedom and parameter matrix the covariance $(1/2)\Sigma_s$. From the Cholesky factorization, $\mathbf{W}_R = \mathbf{D}\mathbf{D}^T$, with

$$\mathbf{D} = \begin{pmatrix} \rho_1 & 0\\ \alpha & \rho_2 \end{pmatrix}. \tag{44}$$

The elements of \mathbf{D} are independent and satisfy [22]

$$\begin{cases} \rho_1 \sim \sqrt{\frac{\Sigma_1}{2} \chi^2 (2T)} \\ \rho_2 \sim \sqrt{\frac{\Sigma_2}{2} \chi^2 (2T-1)} \\ \alpha \sim \mathcal{N}\left(0, \frac{\Sigma_2}{2}\right) \end{cases}$$
(45)

where $\mathcal{N}(0,\varepsilon)$ is a Gaussian distribution with mean value 0 and variance ε and where $\chi^2(P)$ is a Chi-square distribution with P degrees of freedom.

The covariance of $\tilde{\boldsymbol{\theta}}$ is given by

 $\overline{\sigma}$

$$\lim_{\sigma \to 0} cov\left(\widetilde{\boldsymbol{\theta}}\right) = \begin{pmatrix} var(\theta_1) & \Psi \\ \Psi & var(\widetilde{\theta}_2) \end{pmatrix}$$
(46)
$$= \frac{1}{2} \mathbb{E} \left[(\mathbf{H} \left(\boldsymbol{\theta}_0\right) \odot \mathbf{W}_R \right)^{-1} \right]$$
$$= \frac{1}{2} \mathbb{E} \left[\frac{1}{\Phi} \begin{pmatrix} h_2 \left(\rho_2^2 + \alpha^2\right) & -h_3 \rho_1 \alpha \\ -h_3 \rho_1 \alpha & h_1 \rho_1^2 \end{pmatrix} \right]$$
(47)

where $var(\hat{\theta}_1)$ (respectively, $var(\hat{\theta}_2)$) is the variance of the first source (respectively, the second source), Ψ is the cross-correlation and Φ = $h_1 h_2 \rho_1^2 \left(\rho_2^2 + \alpha^2 \right) - (h_3 \rho_1 \alpha)^2.$ From (47)

$$var(\tilde{\theta}_{1}) = \frac{1}{2} \mathbb{E} \left[\frac{h_{2} \left(\rho_{2}^{2} + \alpha^{2}\right)}{h_{1}h_{2}\rho_{1}^{2} \left(\rho_{2}^{2} + \alpha^{2}\right) - \left(h_{3}\rho_{1}\alpha\right)^{2}} \right]$$
$$= \frac{1}{2h_{1}} \mathbb{E} \left[\frac{1}{\rho_{1}^{2}} \frac{1}{1 - \frac{h_{3}^{2}}{h_{1}h_{2}} \frac{\alpha^{2}}{\alpha^{2} + \rho_{2}^{2}}} \right]$$
(48)

$$\begin{cases} \mathbf{A}^{H}\left(\boldsymbol{\theta}_{0}\right)\mathbf{A}\left(\boldsymbol{\theta}_{0}\right)\Big|_{m,n} = \sum_{i=1}^{M} e^{j\left(2\pi/\lambda\right)\mathbf{v}_{i}^{T}\left(\mathbf{u}\left(\boldsymbol{\theta}_{n}\right)-\mathbf{u}\left(\boldsymbol{\theta}_{m}\right)\right)} \\ \mathbf{D}^{H}\left(\boldsymbol{\theta}_{0}\right)\mathbf{D}\left(\boldsymbol{\theta}_{0}\right)\Big|_{m,n} = \sum_{i=1}^{M} \left(\frac{2\pi}{\lambda}\right)^{2} \left(\mathbf{v}_{i}^{T}\dot{\mathbf{u}}\left(\boldsymbol{\theta}_{m}\right)\right) \left(\mathbf{v}_{i}^{T}\dot{\mathbf{u}}\left(\boldsymbol{\theta}_{n}\right)\right) e^{j\left(2\pi/\lambda\right)\mathbf{v}_{i}^{T}\left(\mathbf{u}\left(\boldsymbol{\theta}_{n}\right)-\mathbf{u}\left(\boldsymbol{\theta}_{m}\right)\right)} \\ \mathbf{A}^{H}\left(\boldsymbol{\theta}_{0}\right)\mathbf{D}\left(\boldsymbol{\theta}_{0}\right)\Big|_{m,n} = \sum_{i=1}^{M} j\frac{2\pi}{\lambda} \left(\mathbf{v}_{i}^{T}\dot{\mathbf{u}}\left(\boldsymbol{\theta}_{n}\right)\right) e^{j\left(2\pi/\lambda\right)\mathbf{v}_{i}^{T}\left(\mathbf{u}\left(\boldsymbol{\theta}_{n}\right)-\mathbf{u}\left(\boldsymbol{\theta}_{m}\right)\right)} \end{cases}$$
(41)

$$\mathbf{A}^{H}\left(\boldsymbol{\theta}_{0}\right)\mathbf{D}\left(\boldsymbol{\theta}_{0}\right)\Big|_{m,n} = j\frac{2\pi}{\lambda}\sum_{i=1}^{M/2} \left(\mathbf{v}_{i}^{T}\dot{\mathbf{u}}\left(\boldsymbol{\theta}_{n}\right)\right)\underbrace{\left(e^{j(2\pi/\lambda)\mathbf{v}_{i}^{T}\left(\mathbf{u}\left(\boldsymbol{\theta}_{n}\right)-\mathbf{u}\left(\boldsymbol{\theta}_{m}\right)\right)}-e^{-j(2\pi/\lambda)\mathbf{v}_{i}^{T}\left(\mathbf{u}\left(\boldsymbol{\theta}_{n}\right)-\mathbf{u}\left(\boldsymbol{\theta}_{m}\right)\right)}\right)}_{\text{imaginary number}}.$$
(42)

where $\alpha^2 \sim (\Sigma_2/2)\chi^2(1)$ and the ratio $\alpha^2/(\alpha^2 + \rho_2^2) = Z$ follows a beta distribution with 1 and 2T - 1 degrees of freedom which is independent of $Y = \rho_1^2$. Therefore, (48) becomes

$$var(\tilde{\theta}_1) = \frac{1}{2h_1} \mathbb{E}\left[\frac{1}{Y}\right] \mathbb{E}\left[\frac{1}{1 - \frac{h_3^2}{h_1 h_2}Z}\right] = \frac{I_1 I_2}{2h_1}.$$
 (49)

 $I_1 = \mathbb{E}[1/Y]$ and $I_2 = \mathbb{E}\left[1/(1 - (h_3^2/h_1h_2)Z)\right]$ satisfy

$$\begin{cases} I_1 = \int_0^\infty \frac{1}{y} \prod_{\chi^2} (y) \, dy \\ I_2 = \int_0^1 \frac{1}{1 - \frac{h_2^2}{h_1 h_2} z} \prod_\beta (z) \, dz \end{cases}$$
(50)

where $\Pi_{\chi^2}(y)$ and $\Pi_{\beta}(z)$ are, respectively, the probability density functions of a chi-square random variable $(\Sigma_1/2)\chi^2(2T)$ and of a beta random variable $\beta(1, 2T - 1)$

$$\begin{cases} \Pi_{\chi^2}(y) = \frac{1}{2^{T-1}\Gamma(T)\Sigma_1^T} y^{T-1} e^{-y/\Sigma_1} \\ \Pi_{\beta}(z) = (2T-1)(1-z)^{2(T-1)}. \end{cases}$$
(51)

When $T \ge 2$, I_1 converges, it is a Gamma function. I_2 is the integral representation of a Gauss hypergeometric function [20, pp. 556–565]

$${}_{2}F_{1}(a,b;c;\omega) = \frac{\Gamma(c)}{\Gamma(b)\Gamma(c-b)} \int_{0}^{1} z^{b-1} (1-z)^{c-b-1} (1-z\omega)^{-a} dz \quad (52)$$

where a = 1, b = 1, c = 2T and $\omega = h_3^2/h_1h_2$. Note that (52) is finite

$$\begin{cases} \text{for all } (a, b, c), & \text{if } -1 < \omega < 1\\ \text{for } c > a + b, & \text{if } \omega = \pm 1. \end{cases}$$
(53)

In our case, **H** is a semi positive definite matrix, then $|\mathbf{H}| \ge 0 \Leftrightarrow \omega = h_3^2/h_1h_2 \le 1$. It signifies that I_2 is finite for $T \ge 2$. Finally

$$\begin{cases} I_1 = \frac{1}{(T-1)\Sigma_1}, \\ I_2 = {}_2F_1\left(1, 1; 2T; \frac{h_3^2}{h_1 h_2}\right). \end{cases}$$
(54)

and

$$var(\tilde{\theta}_{1}) = \frac{{}_{2}F_{1}\left(1,1;2T;\frac{\hbar_{3}^{2}}{\hbar_{1}\hbar_{2}}\right)}{2\left(T-1\right)h_{1}\Sigma_{1}}.$$
(55)

Similarly

$$var(\tilde{\theta}_2) = \frac{{}_2F_1\left(1,1;2T;\frac{h_3^2}{h_1h_2}\right)}{2\left(T-1\right)h_2\Sigma_2}.$$
(56)

It can be easily shown that $\Psi = 0$ [see (46)]: it is the integral from minus infinity to plus infinity of an odd function of the variable α . According to (26), the UCRB in the two sources case is

$$\lim_{\sigma \to 0} \frac{1}{\sigma^2} \mathbf{B}_{\text{UCOND}} \left(\boldsymbol{\theta}_0 \right) = \frac{1}{2T} \text{Diag} \left\{ \frac{1}{h_1 \Sigma_1}, \frac{1}{h_2 \Sigma_2} \right\}.$$
(57)

Therefore, using (46), (55)–(57), one obtains (31).

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