

Super-Exponential Blind Adaptive Beamforming

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Abstract—The objective of the beamforming with the exploitation of a sensor array is to enhance the signals of the sources from desired directions, suppress the noises and the interfering signals from other directions, and/or simultaneously provide the localization of the associated sources. In this paper, we present a higher order cumulant-based beamforming algorithm, namely, the super-exponential blind adaptive beamforming algorithm, which is extended from the super-exponential algorithm (SEA) and the inverse filter criteria (IFC). While both SEA and IFC assume noise-free conditions, this requirement is no longer needed, and all the noise components are taken into account in the proposed algorithm. Two special conditions are derived under which the proposed blind beamforming algorithm achieves the performance of the corresponding optimal nonblind beamformer in the sense of minimum mean square error (MMSE). Simulation results show that the proposed algorithm is effective and robust to diverse initial weight vectors; its performance with the use of the fourth-order cumulants is close to that of the nonblind optimal MMSE beamformer.

Index Terms—Adaptive array processing, blind beamforming, high-order cumulants, inverse filter criteria (IFC), super-exponential algorithm (SEA).

I. INTRODUCTION

THE concept of the adaptive antenna [1] was first proposed more than three decades ago. Since then, the adaptive array or, in general, adaptive beamforming techniques, have been widely investigated for various purposes [2], [3], [20], [26]. In a variety of applications, the objective of the beamforming is to enhance the signals of the sources from desired directions, suppress the noises and the interfering signals from other directions, and/or simultaneously provide the localization of the associated source [20]. The optimal beamformer (or spatial filter), in the sense of minimum mean square error (MMSE) criteria, is specified by the Wiener–Hopf equation [2], [3], provided that the reference signal waveform or, alternatively, the

spatial signature (spatial response of the associated antenna array) of the desired source is available. However, it is often encountered in practical situations that the duration of the reference signal is not sufficiently long or that the reference signal is not available at all. The former situations happen, for example, in many mobile communication applications where the length of pilot signals is limited to maintain a high communication capacity. The later situations arise in communication countermeasure and various passive radar and sonar applications. In such situations, adaptive algorithms based on the MMSE criterion cannot work properly. As a result, it becomes increasingly important to develop blind adaptive beamforming algorithms in the absence of a reference signal.

The constant modulus algorithm (CMA) is a well-known algorithm originally developed for blind equalization [15], [16] and has recently been used for blind beamforming [17]. Instead of using training signals, the CMA uses the constant amplitude property of the desired signals. Analyses reported in [18] and [19] have shown that, in certain cases where, for instance, the source signals are of high signal-to-noise ratio (SNR) and are not closely spaced, the steady-state and convergence performances of the CMA-based blind beamforming algorithm are very close to those obtained from the corresponding nonblind ones. However, the CMA-based blind beamformers usually incur considerable performance degradation when these conditions are not satisfied.

Applying higher order statistics, which are usually described in terms of cumulants [7], or cyclostationary properties of signals to blind beamforming is considered to be a new trend to effectively make use of the inherent properties of the signals [11]–[14]. However, the methods reported in the open literature so far are only applicable to certain specific signal environments. For example, the kurtosis (fourth-order cumulant) maximization algorithm [11] globally converges only in noise-free situations. The fourth-order cumulant-based blind beamforming approach [12], [13] assumes that the interfering signals are Gaussian distributed, that incoming signals are coherent, or that the incoming signals are independent and with the identical fourth-order cumulants. On the other hand, to apply cyclostationarity in blind adaptive beamforming, *a priori* knowledge of the exact cyclic frequency of the desired signal is required [14]. Adaptive beamforming algorithms based on cyclic correlation do not work properly when the cyclic frequency of the desired source signal is not accurately estimated [14].

The recent development of the cumulants-based blind deconvolution algorithms, namely, the super-exponential algorithm (SEA) [4], [6] and the inverse filter criteria (IFC) [5], [8], [9], opened a new avenue for blind deconvolution and equalization

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[4]. The extensions of SEA [6] to multi-input multi-output (MIMO) blind deconvolution problems proved to be effective for blind equalization in multiuser wireless communications [10], [22]–[24]. All the formulations of these single-input single-output (SISO) or MIMO blind deconvolution problems make the following assumptions:

- 1) infinite length of the overall channel (convolution of the channel and the equalizer);
- 2) finite length of the FIR filters;
- 3) the absence of noise.

The noise-free assumption is basically justified by using a high-input SNR. According to [6], [10], and [22]–[24], an implicit condition required for formulating the SEA is that the number of independent sources (which also include independent intersymbol interference components, if there are any) must exceed the degrees of freedom (DOFs) of the system [3]. Violation of this condition will result in failure to apply SEA because the corresponding correlation matrix is rank deficient and, therefore, noninvertible [6], [22], [23]. This is a strict constraint and is in conflict with the well-known requirement that the system DOFs must exceed the number of independent sources for effective interference suppression. Situations with the number of the independent sources being less than either the number of array sensors or the number of the total taps of the space-time FIR filters are often encountered in various communications and radar applications. In spite of such importance, applications of SEA to such practical scenarios have not been well investigated thus far.

Another very important factor to be considered is the system noise, which highly limits the system performance. In practice, noise is always present, and the input SNR could be moderate or even low (examples of such situations include countermeasure in communications, radar, and sonar, etc.). This fact underscores the importance of developing new formulations of the SEA with the noises taken into account.

The purpose of this paper is to formulate an SEA-based blind beamforming algorithm with the consideration of array noise and interference signals. The proposed algorithm works robustly, independent of whether the number of independent sources is more or less than the DOFs of the array system. Without loss of generality, we confine our formulation to the classical adaptive beamforming for independent non-Gaussian source signals.

The main contributions of this paper are three fold. First, this paper presents a super-exponential blind beamforming algorithm with the consideration of both noise and interference signals. Second, two special conditions are derived for the algorithm to converge to the optimal beamformer specified by the Wiener–Hopf solution. Third, this paper proves that the inverse filter criteria [8], [9] can be used as the objective function for the proposed beamforming algorithm, and a complete blind beamforming algorithm is subsequently obtained.

The paper is organized as follows. In Section II, we formulate the super-exponential blind beamforming algorithm and prove the two special conditions. Section III reveals the relationship between the proposed algorithm and the existing IFC, and the complete blind beamforming algorithm is then derived. Sec-

tion IV presents simulation results of the conventional applications. Section V concludes this paper.

II. FORMULATION OF BLIND ADAPTIVE BEAMFORMING

A. Formulation of Nonblind Adaptive Beamforming

Let us consider an antenna array of $N(N > 1)$ elements with $P(P \geq 1)$ incident non-Gaussian sources. The signal vector $\mathbf{x}(t)$ at the antenna array is expressed as

$$\mathbf{x}(t) = \sum_{p=1}^P \Gamma \mathbf{a}(\theta_p) b_p(t) + \vec{\omega}(t) \quad (1)$$

where $\mathbf{a}(\theta)$ and $\vec{\omega}(t) = [\varpi_1(t), \dots, \varpi_N(t)]^T$ (the superscript “ T ” denotes transpose) are the array steering vector and the array noise vector, respectively, and $\Gamma = \text{diag}\{\zeta_i e^{j\varphi_i}, i = 1, \dots, N\}$ is a diagonal matrix denoting array channels’ complex gains (gains and phases). The steering vector, for example, of a half-wavelength spaced linear array can be expressed as

$$\mathbf{a}(\theta) = \left[e^{-j(2\pi/\lambda)((N-1)/2)d\sin(\theta)}, e^{-j(2\pi/\lambda)((N-3)/2)d\sin(\theta)}, \dots, e^{j(2\pi/\lambda)((N-3)/2)d\sin(\theta)}, e^{j(2\pi/\lambda)((N-1)/2)d\sin(\theta)} \right]^T$$

with the phase reference point taken at the center of the array. To simplify the notation without loss of generality, hereafter, let us use $\mathbf{a}(\theta)$ instead of $\Gamma \mathbf{a}(\theta)$.

The following assumptions are made for the sources’ signals and the employed antenna array.

- A1) The source signals $b_p(t)$, $p = 1, \dots, P$ are independent, stationary, and non-Gaussian processes with zero mean and variance $E|b_p(t)|^2 = \sigma_{b_p}^2$.
- A2) The noise vector $\vec{\omega}(t)$ is zero-mean, spatially white (but not limited) with $E\{\vec{\omega}(t)\vec{\omega}^T(t)\} = 0$ and $E\{\vec{\omega}(t)\vec{\omega}^H(t)\} = \sigma_{\vec{\omega}}^2 \mathbf{I}_{N \times N}$, where the superscript “ H ” denotes conjugate transpose, $\sigma_{\vec{\omega}}^2$ expresses the noise power at each element, and $\mathbf{I}_{N \times N}$ is the $N \times N$ identity matrix. $\vec{\omega}(t)$ is uncorrelated with all the source signals.
- A3) The steering vectors of all the associated sources are linearly independent when $P \leq N$.

For the purpose of later use, we denote $b_p(t) = \sigma_{b_p} \bar{b}_p(t)$, $p = 1, \dots, P$, $\varpi_l(t) = \sigma_{\varpi} \bar{\varpi}_l(t)$, $l = 1, \dots, N$, $\mathbf{I}_{N \times N} = [\mathbf{e}_{N \times 1}^{(1)}, \dots, \mathbf{e}_{N \times 1}^{(N)}]$, and $\bar{\mathbf{A}}(\Theta) = [\sigma_{b_1} \mathbf{a}(\theta_1), \dots, \sigma_{b_P} \mathbf{a}(\theta_P)]$, and we define

$$\mathbf{H} = \left[\sigma_{b_1} \mathbf{a}(\theta_1), \dots, \sigma_{b_P} \mathbf{a}(\theta_P), \sigma_{\varpi} \mathbf{e}_{N \times 1}^{(1)}, \dots, \sigma_{\varpi} \mathbf{e}_{N \times 1}^{(N)} \right] \\ = \left[\bar{\mathbf{A}}(\Theta)_{N \times P} \sigma_{\varpi} \mathbf{I}_{N \times N} \right]_{N \times (N+P)} \quad (2)$$

and

$$\vec{\alpha}(t) = [\alpha_1(t), \dots, \alpha_{P+N}(t)]^T \\ = [\bar{b}_1(t), \dots, \bar{b}_P(t), \bar{\varpi}_1(t), \dots, \bar{\varpi}_N(t)]^T. \quad (3)$$

Then

$$E[\vec{\alpha}(t)\vec{\alpha}^H(t)] = \mathbf{I}_{(P+N) \times (P+N)}. \quad (4)$$

With straightforward manipulations, the signal vector of (1) can be expressed as the following linear vector model:

$$\begin{aligned} \mathbf{x}(t) &= \left[\sigma_{b_1} \mathbf{a}(\theta_1), \dots, \sigma_{b_P} \mathbf{a}(\theta_P), \sigma_{\varpi} \mathbf{e}_{N \times 1}^{(1)}, \dots, \sigma_{\varpi} \mathbf{e}_{N \times 1}^{(N)} \right] \\ &\quad \times \begin{bmatrix} \bar{b}_1(t) \\ \vdots \\ \bar{b}_P(t) \\ \bar{\varpi}_1(t) \\ \vdots \\ \bar{\varpi}_N(t) \end{bmatrix} \\ &= \mathbf{H}_{N \times (N+P)} \bar{\boldsymbol{\alpha}}(t)_{(N+P) \times 1}. \end{aligned} \quad (5)$$

In order to avoid confusion between the subsequent SEA extension and that presented in [6], [10], and [22]–[24], we highlight the following two cases for ease of understanding.

Case 1) $1 \leq P < N$. This represents the kind of typical cases in communications, radar, and sonar applications.

Case 2) $P \geq N$. This is the case where the associated system does not have enough DOFs.

According to [6], [10], and [22]–[24], an implicit condition necessary to formulate the existing SEA is that the number of independent sources (including independent intersymbol components, if there are any) is greater than the DOFs of the system [3]. It is obvious that Case 2 satisfies this requirement. The above-mentioned extension [6], [10], [22]–[24] can be directly applied to Case 2, regardless of whether the system noises are considered or not.

Remark 1: Note that in the absence of noises and for $P \leq N$ (i.e., the assumption A2 in [23]), the signal model expressed by (1) and (5) is a special case of [23, (1) and (2)] by setting $k = 0$. It is clear that \mathbf{H} in (5) is of full column rank if $1 \leq P < N$ as the noises are absent.

In this paper, we concentrate on Case 1 with the presence of noise. In (5), the dimension of the combined signal vector $\bar{\boldsymbol{\alpha}}$ is $(N + P) \times 1$, which, regardless of the value of N and P , is always greater than N .

Using assumptions A1 and A2, the correlation matrix of the array signal vector of $\mathbf{x}(t)$ can be expressed as

$$\begin{aligned} \mathbf{R} &= E[\mathbf{x}(t)\mathbf{x}^H(t)] = \mathbf{H}E[\bar{\boldsymbol{\alpha}}(t)\bar{\boldsymbol{\alpha}}^H(t)]\mathbf{H}^H = \mathbf{H}\mathbf{H}^H \\ &= \bar{\mathbf{A}}(\Theta)\bar{\mathbf{A}}^H(\Theta) + \sigma_{\varpi}^2 \mathbf{I}_{N \times N}. \end{aligned} \quad (6)$$

The output of adaptive antenna is formulated as

$$z(t) = \mathbf{c}^T \mathbf{x}(t) = \mathbf{c}^T \mathbf{H} \bar{\boldsymbol{\alpha}}(t) \quad (7)$$

where $\mathbf{c} = [c_1, c_2, \dots, c_N]^T$ is the complex weight vector. Define

$$\mathbf{s} = [s_1 \quad \cdots \quad s_{P+N}]^T = \mathbf{H}^T \mathbf{c}_{N \times 1} \quad (8)$$

as the gain vector of all the signals and the noises. Then, the output of the adaptive antenna can be rewritten as

$$z(t) = \mathbf{c}^T \mathbf{x}(t) = \mathbf{s}^T \bar{\boldsymbol{\alpha}}(t) = \sum_{n=1}^{P+N} s_n \alpha_n(t). \quad (9)$$

Note that there are only N weights that are actually controllable, and the gains corresponding to all the signals and the noises depend on the N controllable weights.

Assume that $b_{n_0}(t)$, $1 \leq n_0 \leq P$ is the source signal of interest. Using assumptions A1 and A2, the mean square error (MSE) between $b_{n_0}(t)$ and the output of the adaptive antenna can be expressed as

$$\begin{aligned} \text{MSE} &= E\left(|b_{n_0}(t) - z(t)|^2\right) = |s_{n_0} - \sigma_{b_{n_0}}|^2 + \sum_{\substack{n=1 \\ n \neq n_0}}^{P+N} |s_n|^2 \\ &= \left\| \mathbf{s} - \sigma_{b_{n_0}} \mathbf{e}_{(P+N) \times 1}^{(n_0)} \right\|^2 = \left\| \mathbf{H}^T \mathbf{c} - \sigma_{b_{n_0}} \mathbf{e}_{(P+N) \times 1}^{(n_0)} \right\|^2. \end{aligned} \quad (10)$$

Minimizing the MSE with respect to \mathbf{c} results in the optimal weight vector, represented by the following Wiener–Hopf equation, which is called the optimal adaptive beamformer or MMSE beamformer [2], [3], [26]:

$$\mathbf{c}_{opt} = \left(\mathbf{H}\mathbf{H}^H \right)^{-1} \left(\sigma_{b_{n_0}} \mathbf{H}^* \mathbf{e}_{(P+N) \times 1}^{(n_0)} \right) = \left(\mathbf{R}^T \right)^{-1} \mathbf{d} \quad (11)$$

where

$$\mathbf{d} = \sigma_{b_{n_0}} \mathbf{H}^* \mathbf{e}_{(P+N) \times 1}^{(n_0)} = \sigma_{b_{n_0}}^2 \mathbf{a}^*(\theta_{n_0}) \quad (12)$$

where the superscript “*” denotes complex conjugate.

B. Super-Exponential Iteration Principle

Without loss of generality, we present the super-exponential iteration principle in the same manner as that of [6] with the use of a finite dimensional complex gain vector \mathbf{s} . Define

$$s'_n = s_n^p (s_n^*)^q, \quad n = 1, \dots, P + N \quad (13a)$$

$$s''_n = \frac{1}{\|s'\|} s'_n, \quad n = 1, \dots, P + N \quad (13b)$$

where $\mathbf{s}' = [s'_1, \dots, s'_{P+N}]^T$, and $\|s'\| = \sqrt{\sum_{n=1}^{P+N} |s'_n|^2}$ is the norm of \mathbf{s}' . Substituting $s_n = s_n^{(k-1)}$ and $s_k = s_n$ into (13a) and (13b), we obtain the following iteration equation:

$$s_n^{(k)} = \frac{\left(s_n^{(k-1)} \right)^p \left(s_n^{(k-1)*} \right)^q}{\sqrt{\sum_{l=1}^{P+N} \left| \left(s_l^{(k-1)} \right)^p \left(s_l^{(k-1)*} \right)^q \right|^2}}, \quad n = 1, \dots, P + N. \quad (14)$$

In the case that $s_n^{(0)} = s_n = |a_n| e^{j\varphi_n}$, $a_n \neq 0$, $n = 1, \dots, P + N$, for $k \geq 1$, we have

$$s_n^{(k)} = \frac{|a_n|^{(p+q)^k} e^{j(p-q)^k \varphi_n}}{\sqrt{\sum_{l=1}^{P+N} |a_l|^{2(p+q)^k}}}, \quad n = 1, \dots, P + N. \quad (15)$$

When $p+q \geq 2$ and $p-q = 1$, if $|a_{n_0}| > |a_n|$, $\forall n \neq n_0$, then, as $k \rightarrow \infty$, $s_n^{(k)}$ goes to

$$\begin{cases} s_{n_0}^{(k)} \xrightarrow{k \rightarrow \infty} e^{j\varphi_{n_0}} \\ |s_n^{(k)}| \xrightarrow{k \rightarrow \infty} 0, & n = 1, \dots, P + N, n \neq n_0 \end{cases} \quad (16)$$

From (15), it is seen that s_n is forced to super-exponentially converge to the desired response, i.e., the magnitude of the leading component approaches 1, and the others approach 0. For more general cases, let us reconsider the iteration procedure with a finite-modulus complex weighting sequence, and then, (13a) and (13b) become

$$s'_n = \beta_n s_n^p (s_n^*)^q, \quad n = 1, \dots, P + N \quad (17a)$$

$$s''_n = \frac{1}{\|s'\|} s'_n, \quad n = 1, \dots, P + N \quad (17b)$$

where $\beta_n = |\beta_n| e^{j\psi_n}$, $|\beta_n| < \infty$, and $n = 1, \dots, P + N$. For $k \geq 1$, the equation corresponding to (14) becomes

$$s_n^{(k)} = \frac{\beta_n \left(s_n^{(k-1)} \right)^p \left(s_n^{(k-1)*} \right)^q}{\sqrt{\sum_{l=1}^{P+N} \left| \beta_l \left(s_l^{(k-1)} \right)^p \left(s_l^{(k-1)*} \right)^q \right|^2}}, \quad n = 1, \dots, P + N. \quad (18)$$

Similar to (15), in the case where $s_n^{(0)} = s_n = |a_n| e^{j\varphi_n}$, $n = 1, \dots, P + N$, for $k \geq 1$, we have

$$s_n^{(k)} = \frac{|\beta_n| \sum_{l=0}^{k-1} (p+q)^l |a_n|^{(p+q)^k} e^{j(p-q)^k \varphi_n + j\psi_n \sum_{l=0}^{k-1} (p-q)^l}}{\sqrt{\sum_{l=1}^{P+N} |\beta_l| \sum_{l=0}^{k-1} (p+q)^l |a_l|^{2(p+q)^k}}}, \quad n = 1, \dots, P + N. \quad (19)$$

When $p + q \geq 2$, (19) can be rewritten as

$$s_n^{(k)} = \frac{|\beta_n|^{((p+q)^k - 1)/(p+q-1)} |a_n|^{(p+q)^k} e^{j(p-q)^k \varphi_n + j\psi_n \sum_{l=0}^{k-1} (p-q)^l}}{\sqrt{\sum_{l=1}^{P+N} |\beta_l|^{2 \times ((p+q)^k - 1)/(p+q-1)} |a_l|^{2(p+q)^k}}}, \quad n = 1, \dots, P + N. \quad (20)$$

The convergent behavior of the above iteration sequence is described by the following theorem.

Theorem 1: For an arbitrary weighting sequence with elements of finite modulus, i.e., $\beta_n = |\beta_n| e^{j\psi_n}$, $|\beta_n| < \infty$, $n = 1, \dots, P + N$, the conditions for the iterative updated vector (20) to super-exponentially converge to a vector that the magnitude of the leading component equal to 1 and the other components equal to 0, i.e.,

$$\left\{ \begin{array}{l} |s_{n_0}^{(k)}| \xrightarrow{k \rightarrow \infty} 1 \\ |s_n^{(k)}| \xrightarrow{k \rightarrow \infty} 0, \quad n = 1, \dots, P + N, n \neq n_0 \end{array} \right. \quad (21)$$

are $p + q \geq 2$, and

$$|a_{n_0}| |\beta_{n_0}|^{1/(p+q-1)} > |a_n| |\beta_n|^{1/(p+q-1)} \quad n = 1, \dots, P + N, n \neq n_0. \quad (22)$$

The proof of Theorem 1 is given in Appendix A.

C. Formulation of Blind Adaptive Beamforming

It is seen that for $k \geq 1$, the associated gain vector is updated by (8) with a given weight vector $\mathbf{c}^{(k)}$, i.e.,

$$\mathbf{s}^{(k)} = [s_1^{(k)} \quad \dots \quad s_{P+N}^{(k)}]^T = \mathbf{H}^T \mathbf{c}^{(k)}. \quad (23)$$

On the other hand, the gain $s_n^{(k)}$ can be directly updated based on $s_n^{(k-1)}$, i.e., $s_n^{(k)} = \beta_n \left(s_n^{(k-1)} \right)^p \left(s_n^{(k-1)*} \right)^q$, $n = 1, \dots, P + N$ in accordance with (17a) without $\mathbf{c}^{(k)}$. This implies that there are evident differences between $s_n^{(k)}$ of (23) and $s_n^{(k)}$ of (17a). However, due to the attractive property of the super-exponential convergence shown in Theorem 1, it is desired that the gains $s_n^{(k)}$, $n = 1, \dots, P + N$ be updated under (17a) and (17b). The best way to solve the problem is to combine the iteration procedures of (17a) and (17b) with that of (23) and then to minimize all the differences. This leads to the following constrained minimum least square (LS) problem, i.e.,

$$\min_{\mathbf{c}^{(k)}} \sum_{n=1}^{P+N} \left| s_n^{(k)} - \beta_n \left(s_n^{(k-1)} \right)^p \left(s_n^{(k-1)*} \right)^q \right|^2$$

subject to $\sum_{n=1}^{P+N} \left| s_n^{(k)} \right|^2 = 1. \quad (24)$

Define

$$\mathbf{g}^{(k-1)} = \left[\beta_1 \left(s_1^{(k-1)} \right)^p \left(s_1^{(k-1)*} \right)^q, \dots, \beta_{P+N} \left(s_{P+N}^{(k-1)} \right)^p \left(s_{P+N}^{(k-1)*} \right)^q \right]^T. \quad (25)$$

Then, substituting (23) into (24) yields

$$\min_{\mathbf{c}^{(k)}} \left\| \mathbf{H}^T \mathbf{c}^{(k)} - \mathbf{g}^{(k-1)} \right\|^2, \quad \text{subject to } \left\| \mathbf{H}^T \mathbf{c}^{(k)} \right\| = 1. \quad (26)$$

With the use of the method of Lagrange multipliers, the solution for $\mathbf{c}^{(k)}$ can be derived straightforwardly as

$$\mathbf{c}^{(k)} = \frac{(\mathbf{H}^* \mathbf{H}^T)^{-1} \mathbf{H}^* \mathbf{g}^{(k-1)}}{\sqrt{(\mathbf{H}^* \mathbf{g}^{(k-1)})^H (\mathbf{H}^* \mathbf{H}^T)^{-1} \mathbf{H}^* \mathbf{g}^{(k-1)}}}. \quad (27)$$

Define

$$\mathbf{d}^{(k-1)} = \mathbf{H}^* \mathbf{g}^{(k-1)} \quad (28)$$

and with the definitions of (6) and (27), (28) can also be expressed as

$$\mathbf{c}^{(k)} = \frac{(\mathbf{R}^T)^{-1} \mathbf{d}^{(k-1)}}{\sqrt{(\mathbf{d}^{(k-1)})^H (\mathbf{R}^T)^{-1} \mathbf{d}^{(k-1)}}} \quad (29)$$

where \mathbf{R} can be estimated from the samples (or array snapshots) of $\mathbf{x}(t)$, and $\mathbf{d}^{(k-1)}$ will be estimated subsequently through a high-order cumulant. Regarding (26), (27), and (29), we have the following remark.

Remark 2: As mentioned in Remark 1, in the absence of noise, \mathbf{H} is of full column rank when $1 \leq P < N$. In this case, the correlation matrix \mathbf{R} is rank-deficient. This implies that in such a case, it is difficult to obtain an analytical solution for the blind beamformer $\mathbf{c}^{(k)}$ from (26). As a result, the iteration procedure given by [23, (39) and (40)] becomes inapplicable.

Therefore, it is necessary to take the noises into account, and treating the noise components as independent sources is a convenient and effective way to solve the problem for Case 1.

Next, we present the method of estimating $\mathbf{d}^{(k-1)}$. First of all, let $\text{cum}(x_1; x_2; \dots; x_L)$ denote the joint cumulant of L random variables $\{x_1, x_2, \dots, x_L\}$, and define

$$\begin{aligned} & \text{cum}\left(z^{(l)}(t) : p; z^{(l)*}(t) : q; \mathbf{x}^*(t)\right) \\ &= \text{cum}\left(\underbrace{z^{(l)}(t); \dots; z^{(l)}(t)}_{p \text{ terms}}; \underbrace{z^{(l)*}(t); \dots; z^{(l)*}(t)}_{q \text{ terms}}; \mathbf{x}^*(t)\right) \end{aligned} \quad (30)$$

where $z^{(k)}(t)$ is the output of the beamformer $\mathbf{c}^{(k)}$ at the k th iteration, i.e.,

$$z^{(k)}(t) = \mathbf{c}^{(k)T} \mathbf{x}(t) = \mathbf{s}^{(k)T} \bar{\alpha}(t) = \sum_{n=1}^{P+N} s_n^{(k)} \alpha_n(t). \quad (31)$$

Using the linearity property of cumulants [7], i.e., $\text{cum}(\sum_i b_i x_i; \dots) = \sum_i b_i \text{cum}(x_i; \dots)$, and substituting (31) and (5) into (30), (30) can be rewritten as (32), shown at the bottom of the page. Define

$$\beta_n = \text{cum}(\alpha_n(t) : p; \alpha_n^*(t) : q + 1), n = 1, \dots, P + N \quad (33)$$

as the cumulants of order $(p, q + 1)$ of all sources. From (25), (28), and (32), we conclude that

$$\mathbf{d}^{(k-1)} = \text{cum}\left(z^{(k-1)}(t) : p; z^{(k-1)*}(t) : q; \mathbf{x}^*(t)\right). \quad (34)$$

According to (31), $\mathbf{d}^{(k-1)}$ can be estimated by $\mathbf{c}^{(k-1)}$ and the data samples $\mathbf{x}(t)$.

Equations (29), (31), and (34), rewritten as

$$\mathbf{c}^{(k)} = \frac{(\mathbf{R}^T)^{-1} \mathbf{d}^{(k-1)}}{\sqrt{(\mathbf{d}^{(k-1)})^H (\mathbf{R}^T)^{-1} \mathbf{d}^{(k-1)}}} \quad (35a)$$

$$\mathbf{d}^{(k-1)} = \text{cum}\left(z^{(k-1)}(t) : p; z^{(k-1)*}(t) : q; \mathbf{x}^*(t)\right) \quad (35b)$$

$$z^{(k-1)}(t) = \left(\mathbf{c}^{(k-1)}\right)^T \mathbf{x}(t) \quad (35c)$$

define the super-exponential blind beamforming algorithm. A remark on it is given as follows.

Remark 3: It is noted that the above blind super-exponential beamforming algorithm is valid for $P \geq 1$ and can be applied regardless of the array configurations and the number of non-Gaussian (i.e., nonzero cumulants) sources, which, therefore, is very important to diverse applications, e.g., communications, radar, sonar, etc.

When the noises are Gaussian distributed, then for $p + q \geq 2$ and $n = P + 1, \dots, P + N$, β_n equals zero [7]. In this case, the beamforming weight vector given by (35a) can be written as

$$\mathbf{c}^{(k)} = \frac{\sum_{i=1}^P \sigma_{b_i} (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_i) \beta_i \left(s_i^{(k-1)}\right)^p \left(s_i^{(k-1)*}\right)^q}{\sqrt{(\mathbf{d}^{(k-1)})^H (\mathbf{R}^T)^{-1} \mathbf{d}^{(k-1)}}}. \quad (36)$$

It is evident from (36) that the iterative weight vector is a linear combination of the optimal weight vectors corresponding to all the source signals. Furthermore, substituting (36) into (8) yields

$$\mathbf{s}^{(k)} = \frac{\sum_{i=1}^P \bar{\xi}(\theta_i) \beta_i \left(s_i^{(k-1)}\right)^p \left(s_i^{(k-1)*}\right)^q}{\left\| \sum_{i=1}^P \bar{\xi}(\theta_i) \beta_i \left(s_i^{(k-1)}\right)^p \left(s_i^{(k-1)*}\right)^q \right\|} \quad (37)$$

where

$$\bar{\xi}(\theta_i) = \sigma_{b_i} \mathbf{H}^T (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_i) \quad (38)$$

is the gain vector of the optimal adaptive antenna with respect to the i th source, and

$$\left\| \bar{\xi}(\theta_i) \right\|^2 = \sigma_{b_i}^2 \mathbf{a}^T(\theta_i) (\mathbf{R}^T)^{-1} \mathbf{a}(\theta_i). \quad (39)$$

It is obvious that the optimum $\mathbf{s}^{(k)}$ is also a linear combination of $\bar{\xi}(\theta_i)$.

From (36) and (37), it is seen that the convergent beamformer $\mathbf{c}^{(k)}$ depends on the initial vector $\mathbf{c}^{(0)}$, $p + q$, input SNRs, the spatial responses of the sources, and spatial correlation between the sources.

We know that due to the constraint of (24), $|s_n^{(k)}| < 1$ always hold for $k \geq 1$ and $n = 1, \dots, P + N$. If the initial weight vector $\mathbf{c}^{(0)}$ is chosen such that

$$\left| s_{n_0}^{(0)} \right| > \max \left\{ \left| s_n^{(0)} \right|, n = 1, \dots, P + N, n \neq n_0 \right\}, 1 \leq n_0 \leq P \quad (40)$$

$$\begin{aligned} \text{cum}\left(z^{(k)}(t) : p; z^{(k)*}(t) : q; \mathbf{x}^*(t)\right) &= \text{cum}\left(z^{(k)}(t) : p; z^{(k)*}(t) : q; \mathbf{H}^* \bar{\alpha}^*(t)\right) \\ &= \mathbf{H}^* \text{cum}\left(z^{(k)}(t) : p; z^{(k)*}(t) : q; \bar{\alpha}^*(t)\right) \\ &= \mathbf{H}^* \text{cum}\left(\underbrace{\sum_{n=1}^{P+N} s_n^{(k)} \alpha_n(t); \dots; \sum_{n=1}^{P+N} s_n^{(k)} \alpha_n(t)}_{p \text{ terms}}; \underbrace{\sum_{n=1}^{P+N} s_n^{(k)*} \alpha_n^*(t); \dots; \sum_{i=1}^{P+N} s_n^{(k)*} \alpha_n^*(t)}_{q \text{ terms}}; \bar{\alpha}^*(t)\right) \\ &= \mathbf{H}^* \begin{bmatrix} \left(s_1^{(k)}\right)^p \left(s_1^{(k)*}\right)^q \text{cum}\left(\alpha_1(t) : p; \alpha_1^*(t) : q + 1\right) \\ \vdots \\ \left(s_{P+N}^{(k)}\right)^p \left(s_{P+N}^{(k)*}\right)^q \text{cum}\left(\alpha_{P+N}(t) : p; \alpha_{P+N}^*(t) : q + 1\right) \end{bmatrix}. \end{aligned} \quad (32)$$

then, from (18) or (37), it is obvious that for arbitrary $|\beta_n| < +\infty$, $n = 1, \dots, P$

$$\mathbf{s}^{(1)} \xrightarrow{p+q \rightarrow \infty} \frac{e^{j(p-q) \arg(s_{n_0}^{(0)})} \vec{\xi}(\theta_{n_0})}{\|\vec{\xi}(\theta_{n_0})\|} \vec{\xi}(\theta_{n_0}) \quad (41)$$

where $\arg(s_{n_0}^{(0)})$ denotes the phase of $s_{n_0}^{(0)}$. Note that (40) and (41) indicate that the super-exponential blind beamforming algorithm converges faster for a larger value of $p+q$. However, a larger value of $p+q$ usually gives rise to a larger variance of the estimated $\mathbf{d}^{(k-1)}$ [see (35b)] thus leading to a larger variance of the designed beamformer. Nevertheless, for an appropriate $p+q \geq 2$, as the gain of all the undesired signals are smaller than that of the desired signal under $\mathbf{c}^{(0)}$, the super-exponential beamforming algorithm will converge fast and closely to the optimal beamformer. In other words, if $\mathbf{c}^{(0)}$ is chosen such that

$$\mathbf{s}^{(0)} = \frac{\vec{\xi}(\theta_{n_0})}{\|\vec{\xi}(\theta_{n_0})\|} \quad (42)$$

then under an appropriate $p+q \geq 2$, we have

$$\mathbf{s}^{(1)} = \frac{e^{j(p-q) \arg(s_{n_0}^{(0)})} \vec{\xi}(\theta_{n_0}) + \vec{\varepsilon}}{\|e^{j(p-q) \arg(s_{n_0}^{(0)})} \vec{\xi}(\theta_{n_0}) + \vec{\varepsilon}\|} \quad (43)$$

where $\vec{\varepsilon}$ is an error vector with $\|\vec{\varepsilon}\|$ much smaller than $\|\vec{\xi}(\theta_{n_0})\|$.

In order to give more insight into the behavior and the performance of the proposed algorithm, we present two theorems in the following for two special cases, respectively. In the first case, the SNRs of all source signals approach infinity, whereas in the second, the desired source is spatially orthogonal to the other sources.

Theorem 2: When $P \leq N$, $p+q \geq 2$, and the noise components are Gaussian distributed, the algorithm specified by (35a)–(35c) performs perfect suppression of all the interfering signals if $\sigma_{b_i}^2/\sigma_{\omega}^2 \rightarrow +\infty$, $i = 1, \dots, P$, and

$$|\beta_{n_0}|^{1/(p+q-1)} \left| s_{n_0}^{(0)} \right| > \max \left\{ |\beta_n|^{1/(p+q-1)} \left| s_n^{(0)} \right| \right. \\ \left. n = 1, \dots, P, n \neq n_0 \right\}$$

where β_n is defined by (33).

Theorem 2 is proved in Appendix B.

Remark 4: Theorem 2 implies that the higher the SNRs of all the source signals, the closer the blind beamformer approaches the nonblind optimal beamformer (the MMSE beamformer). In addition, the blind beamformer approaches the nonblind one when the SNRs of all the sources approach infinity.

Theorem 3: When $P \leq N$, $p+q \geq 2$, and the noise components are Gaussian distributed, the algorithm specified by (35a)–(35c) converges to the nonblind optimal (MMSE) beamformer if the steering vector of the n_0 th ($1 \leq n_0 \leq P$) source is spatially orthogonal to those of all the other sources, i.e.,

$$\mathbf{a}^T(\theta_n) \mathbf{a}^*(\theta_{n_0}) = 0, n = 1, \dots, P, n \neq n_0 \quad (44)$$

and

$$\mu_0 |\beta_{n_0}|^{1/(p+q-1)} \left| s_{n_0}^{(0)} \right| > \tilde{\beta}_{r_0}^{1/(p+q-1)} \left| s_{m_0}^{(0)} \right| \\ n_0 \neq m_0, 1 \leq m_0, r_0 \leq P \quad (45)$$

where

$$\beta_n = \begin{cases} \text{cum}(\alpha_n(t) : p; \alpha_n^*(t) : q + 1), & n = 1, \dots, P \\ 0, & n = P + 1, \dots, P + N \end{cases} \quad (46)$$

$$\tilde{\beta}_{r_0} = \max \left\{ \sum_{\substack{i=1 \\ i \neq n_0}}^P \sigma_{b_n} \sigma_{b_i} \mathbf{a}^T(\theta_n) (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_i) \|\beta_i\|, n=1, \dots, P, n \neq n_0 \right\} \quad (47)$$

$$\left| s_{m_0}^{(0)} \right| = \max \left\{ \left| s_n^{(0)} \right|, n = 1, \dots, P \right\} \quad (48)$$

$$\mu_0 = \frac{\frac{\|\mathbf{a}(\theta_{n_0})\|^2 \sigma_{b_{n_0}}^2}{\sigma_{\omega}^2}}{1 + \frac{\|\mathbf{a}(\theta_{n_0})\|^2 \sigma_{b_{n_0}}^2}{\sigma_{\omega}^2}} < 1. \quad (49)$$

The proof is in Appendix C.

Remark 5: Note that under (45), the algorithm converges to the conventional beamformer shown by (C8) in Appendix C. This provides some insights of the performance behavior of the proposed algorithm.

III. OBJECTIVE FUNCTION OF THE BLIND ADAPTIVE BEAMFORMING

In many adaptive processing problems, we know that appropriate objective functions, or cost functions, are very important in developing adaptive algorithms and evaluating their performance. It has been shown in [5], [8], and [9] that the following inverse filter criteria

$$J_{m,n} = \frac{|C_{m,n}\{e(k)\}|}{|C_{1,1}\{e(k)\}|^{(m+n)/2}} = \frac{|C_{m,n}\{e(k)\}|}{\left| E \left\{ |e(k)|^2 \right\} \right|^{(m+n)/2}} \quad (50)$$

are suitable objective functions for blind deconvolution, where $e(k)$ represents the output of the associated blind equalizer, $m+n > 2$, and the cumulant of $e(k)$ is defined as

$$C_{m,n}\{e(k)\} = \text{cum} \{e(k) : m, e^*(k) : n\} \\ = \text{cum} \left\{ \underbrace{e(k), \dots, e(k)}_{m \text{ terms}}, \underbrace{e^*(k), \dots, e^*(k)}_{n \text{ terms}} \right\}. \quad (51)$$

In the absence of noises, the objective function is bounded by the maximum modulus of the cumulants of all the signal components [5], [8], [9]. In this section, our purpose is to show that the inverse filter criteria can also be an objective function for the proposed blind beamforming algorithm.

Replacing $e(k)$ in (50) with the $z(t)$ in (7)–(9), one can obtain

$$J_{m,n} = \frac{|C_{m,n}\{z(t)\}|}{|C_{1,1}\{z(t)\}|^{(m+n)/2}} = \frac{\left| \sum_{l=1}^{P+N} \gamma_l^{(m,n)} s_l^m (s_l^*)^n \right|}{\left| \sum_{l=1}^{P+N} \gamma_l^{(1,1)} |s_l|^2 \right|^{(m+n)/2}} \quad (52)$$

where

$$\gamma_l^{(m,n)} = \text{cum} \left\{ \underbrace{\alpha_l(t), \dots, \alpha_l(t)}_{m \text{ terms}}, \underbrace{\alpha_l^*(t), \dots, \alpha_l^*(t)}_{n \text{ terms}} \right\} \quad (53)$$

represents the (m, n) th-order cumulant of the l th source signal. According to (4), we have

$$\gamma_l^{(1,1)} = \text{cum} \{ \alpha_l(t), \alpha_l^*(t) \} = E \{ |\alpha_l(t)|^2 \} = 1. \quad (54)$$

Defining

$$\delta_s = \left| \sum_{l=1}^{P+N} |s_l|^2 \right|^{(1/2)} = (\mathbf{c}^T \mathbf{R} \mathbf{c}^*)^{(1/2)} \quad (55)$$

which is the root mean square (RMS) of the total output power, (52) can then be expressed as

$$J_{m,n}(\mathbf{c}) = \left| \sum_{l=1}^{P+N} \gamma_l^{(m,n)} \left(\frac{s_l}{\delta_s} \right)^m \left(\frac{s_l^*}{\delta_s} \right)^n \right|. \quad (56)$$

Since most of the digital signal waveforms employed in communications and radar applications are of symmetric distributions, it can be shown that all their odd-order moments and cumulants are equal to zero. In the following, the situation where $m = n > 1$ is considered. In this case, (56) becomes

$$J_{2m}(\mathbf{c}) = J_{m,m}(\mathbf{c}) = \left| \sum_{l=1}^{P+N} \gamma_l^{(m,m)} \left| \frac{s_l}{\delta_s} \right|^{2m} \right|. \quad (57)$$

From the above equation, we see that for a given n_0 , $1 \leq n_0 \leq P$, the necessary and sufficient condition for $J_{2m}(\mathbf{c}) = \left| \gamma_{n_0}^{(m,m)} \right|$ is that $|s_{n_0}| = 1$ and $s_n = 0$, $n = 1, \dots, P + N$, $n \neq n_0$, which corresponds to the perfect suppression of all the interfering signals described in Theorem 2. In practice, when the blind beamformer makes $|s_{n_0}| > |s_n|$, $n = 1, \dots, P + N$, $n \neq n_0$, $1 \leq n_0 \leq P$, then (57) can be expressed as

$$J_{2m}(\mathbf{c}) = \left| \gamma_{n_0}^{(m,m)} \left| \frac{s_{n_0}}{\delta_s} \right|^{2m} \left[1 + \sum_{\substack{l=1 \\ l \neq n_0}}^{P+N} \frac{\gamma_l^{(m,m)}}{\gamma_{n_0}^{(m,m)}} \left(\left| \frac{s_l}{s_{n_0}} \right|^2 \right)^m \right] \right| \quad (58)$$

where the ratio $(|s_l|^2/|s_{n_0}|^2)^m$ will decrease as m increases. This indicates that the larger the m , i.e., the higher the employed order of the cumulant, the more the reduced contributions of $\gamma_n^{(m,m)}$, $n = 1, \dots, P + N$, $n \neq n_0$ to $J_{2m}(\mathbf{c})$. The term $|s_{n_0}|^2/|s_l|^2$ denotes the suppression ratio of the beamformer to the l th source. The higher the suppression ratio, the smaller the contribution of the l th source to the objective function.

The gradient of the objective function with respect to the weight vector is derived as

$$\begin{aligned} & \frac{\partial J_{2m}(\mathbf{c})}{\partial \mathbf{c}} \\ &= \frac{\text{sgn}(C_{m,m} \{z(t)\}) \times m \times \text{cum}\{z(t):m, z^*(t):m-1, \mathbf{x}^*(t)\}}{(\mathbf{c}^T \mathbf{R} \mathbf{c}^*)^m} \\ & \quad - \frac{m \times |C_{m,m} \{z(t)\}|}{(\mathbf{c}^T \mathbf{R} \mathbf{c}^*)^{m+1}} \times \mathbf{R}^T \mathbf{c} \end{aligned} \quad (59)$$

where $\text{sgn}(C_{m,m} \{z(t)\})$ expresses the sign of the cumulant $C_{m,m} \{z(t)\}$. Because $J_{2m}(\mathbf{c})$ is a nonlinear function with respect to \mathbf{c} , the stationary points [5], [8], [11] of (58) can be obtained from

$$\frac{\partial J_{2m}(\mathbf{c})}{\partial \mathbf{c}} = \mathbf{0}. \quad (60)$$

We have the following theorem to show the relationship between the stationary points and the algorithm.

Theorem 4: The super-exponential iterative algorithm specified by (35a)–(35c) converges to the stationary points of the objective function (58).

Theorem 4 is proved in Appendix D.

The associated cumulant terms in (35a)–(35c) and (50) with order (2,2), (3,3), and (4,4) are derived in Appendix E.

From Appendix D, we see that the stationary points of the objective function cannot be expressed by an analytical solution, and the algorithm described by (35a) and (35b) satisfies (D2), which specifies the stationary points. Because the initial $\mathbf{c}^{(0)}$ could be different, the algorithm may converge to the different stationary points. In essence, Theorem 4 shows that the inverse filter criteria can be an objective function of the blind beamforming. Based on Theorem 4, the complete super-exponential blind adaptive beamforming algorithm can be cast into the following three steps.

- Step 1) Let $k = 1$, given the accuracy $\rho (= 10^{-5})$, $\mathbf{c}^{(0)}$ and the order of the cumulant.
- Step 2) Let $p = m$ and $q = p - 1$, perform (35a)–(35c), and obtain $\mathbf{c}^{(k)}$.
- Step 3) If $|J_{2m}(\mathbf{c}^{(k)}) - J_{2m}(\mathbf{c}^{(k-1)})| < \rho$, then stop; Else let $k = k + 1$, go to Step 1.

IV. COMPUTER SIMULATIONS

In order to confirm the effectiveness of the proposed algorithm, computer simulations are performed. In the simulations, a six-element uniform linear array with half-wavelength spacing is employed. The array is not calibrated, i.e., the complex gains of the array elements are perturbed from their nominal value. Herein, we assume that the perturbation gains are time-invariant, whereas their exact values are unknown. The respective nominal values of the amplitude and phase of the complex gains are set to 1 (0 dB) and 0 rad. We choose the perturbation gain from the samples of uniformly distributed random numbers, where the standard deviation of the amplitude is assumed as 0.1, and that of the phase is assumed as 0.1 rad. As a result, the perturbation gain vector used in the simulations is randomly generated, and the set of the samples

$$\{1.0236 - 0.1233j, 0.8636 - 0.0075j, 1.1470 + 0.1185j, 1.0491 - 0.0626j, 0.8854 - 0.1041j, 1.1536 + 0.1937j\}$$

is selected.

We consider the scenario where three signal sources illuminate the antenna array from the far field. All the source signals are quadrature phase shift keying signals. The strongest source is assumed to be the desired one, which is located at 12.3° from the broadside of the antenna array (the broadside here is considered as the reference direction-of-arrival signals). The power of the second source is assumed to be 0.5 dB lower than that of the desired signal, and its incident angle is 16.5° from the broadside. The power of the third source is 2 dB lower than that of the desired signal, and its incident angle is -1.7° from the broadside.

As mentioned in Sections II and III, different initial weight vectors may lead to different convergent results. For comparison, the following two initial weight vectors are considered.

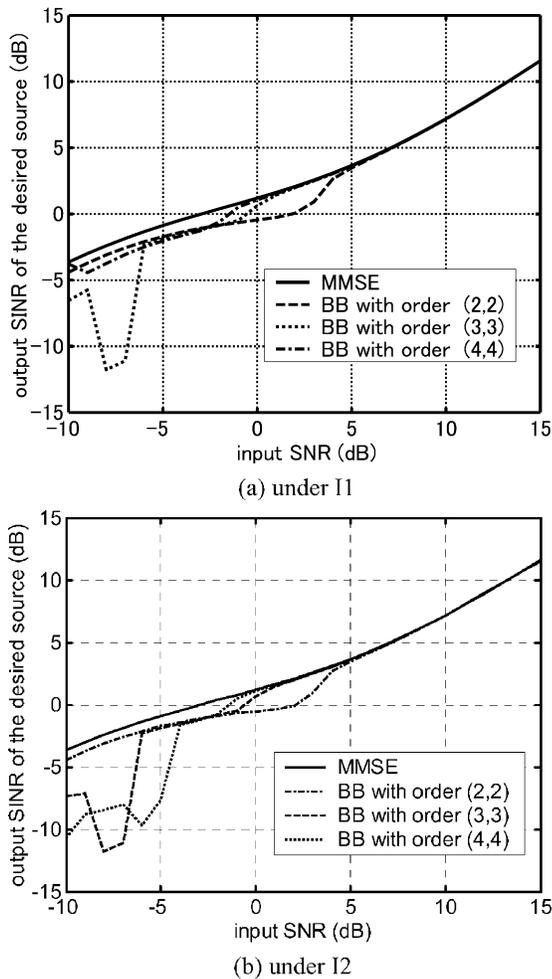


Fig. 1. Output SINR of the desired source.

I1)

$$\begin{aligned} \mathbf{c}^{(0)} &= \mathbf{a}(\theta)|_{\theta=12.3^\circ} \\ &= \begin{bmatrix} e^{-j(2\pi/\lambda)(5/2)d \sin(\theta)}, e^{-j(2\pi/\lambda)(3/2)d \sin(\theta)} \\ \dots, e^{j(2\pi/\lambda)(3/2)d \sin(\theta)}, e^{j(2\pi/\lambda)(5/2)d \sin(\theta)} \end{bmatrix}^T \Big|_{\theta=12.3^\circ} \end{aligned}$$

$$I2) \quad \mathbf{c}^{(0)} = [1, 0, \dots, 0]^T.$$

It is obvious that under I1, a beam toward the desired source is formed, whereas the array response (pattern) under I2 is omnidirectional. Although better beamformer performance is expected when I1 is used, nevertheless, I2 is often used in practice because the *a priori* information of the spatial signature of the desired source is usually not available.

In Fig. 1(a) and (b), we depict the output SINR versus the input SNR under the initial weight vectors I1 and I2, respectively. The number of the samples employed is 5000, and $\rho = 1 \times 10^{-5}$ is used. In these figures, “BB” stands for the proposed blind algorithm and “MMSE” stands for the nonblind MMSE optimal beamformer. It is evident from these figures that when the input SNR is not too low (input SNR > -5 dB), the higher the order of the cumulants used, the closer the performance of the blind algorithm approaches that of the MMSE beamformer. It is also seen that the results under I1 are slightly better than

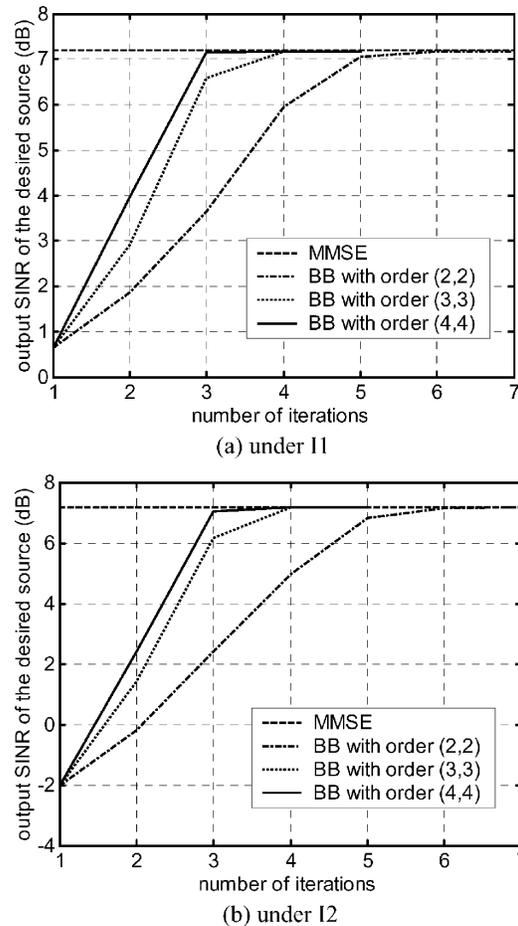


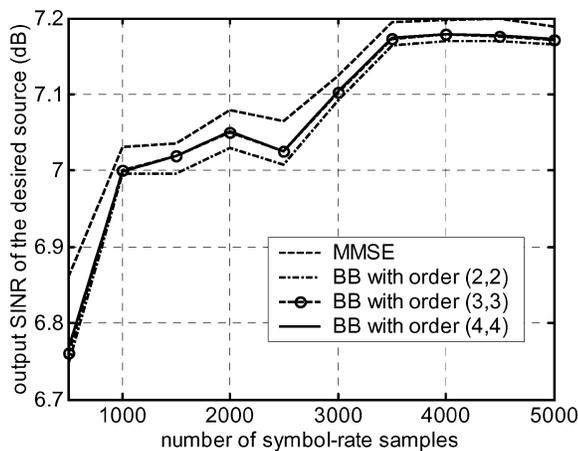
Fig. 2. Comparison of the convergence performance.

that under I2. When the input SNR is high, on the other hand, all the results of the blind algorithm with both initial weight vectors I1 and I2 and different cumulant orders of (2,2), (3,3), and (4,4) approach that of the MMSE beamformer with negligible difference.

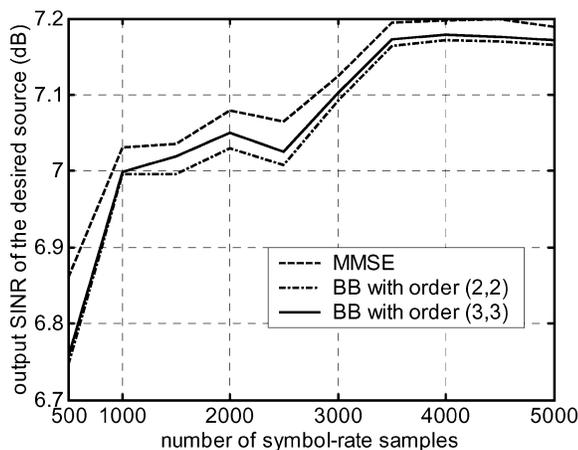
Fig. 2(a) and (b) show the convergence performance of the blind algorithm under I1 and I2, respectively, where the input SNR is 10 dB, and 5000 data samples are employed. From these figures, we see that the convergence rates under both the two initial weight vectors are basically the same.

To investigate the effects of the number of the array signal samples on the performance of the proposed algorithm, Fig. 3 shows the output SINR versus the number of the array data samples, where the input SNR is fixed to 10 dB. It is seen that the longer the samples, the more closely the blind algorithm under both I1 and I2 will converge to the optimum result. However, because the power of the second source signal is close to that of the desired source signal, it would be possible for the algorithm to mistakenly converge to the result corresponding to the second source when I2 is employed. Fig. 3(c) shows the phenomenon for the case of the sample size equal to 1000. This is natural because all blind algorithms are blind to the order of all the corresponding signals, provided that the *a priori* information is available to distinguish them.

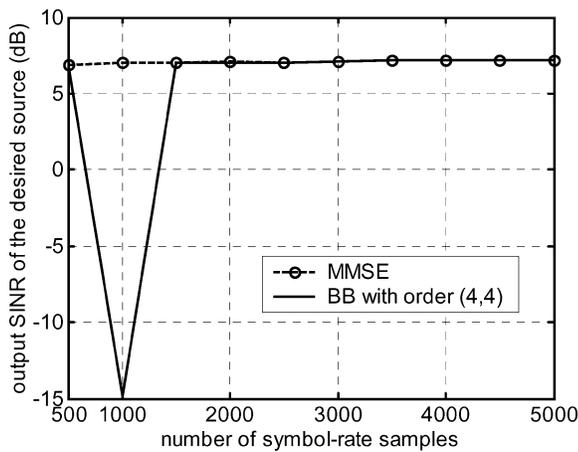
Finally, we show the beam patterns in Fig. 4 with both initial conditions, where 5000 samples are used, and the input SNR



(a) under I1



(b) under I2

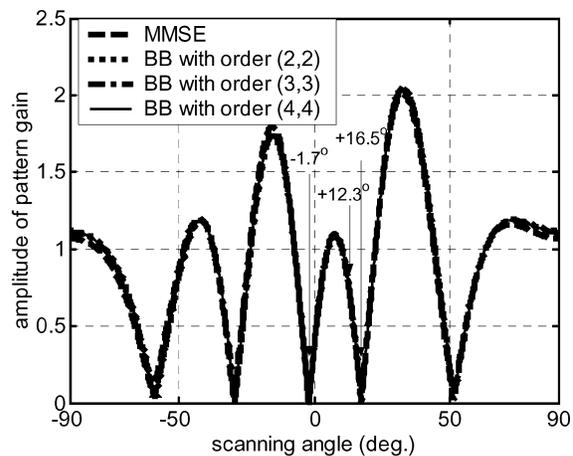


(c) under I2

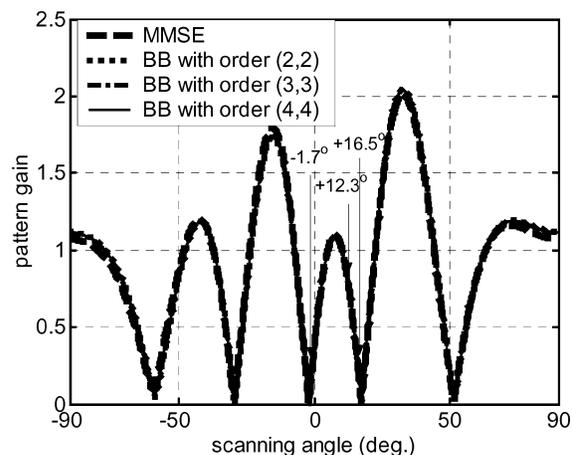
Fig. 3. Effects of number of samples on performance.

is 10 dB. It is seen that the convergent beam patterns corresponding to both initial conditions are very close to that of the optimal nonblind MMSE beamformer.

As a conclusion of the simulation results depicted in Figs. 1–4, we see that the proposed beamforming algorithm is robust and effective. The simulation results demonstrated that the cumulant orders, i.e., (2, 2), (3, 3), and (4, 4) are applicable, and using the order (2, 2) leads to lower computations. Although using higher order cumulants is helpful to improve the convergence rate, the algorithm will be more sensitive to the



(a) under I1



(b) under I2

Fig. 4. Convergent patterns.

accuracy of cumulant estimation when a higher order is used. Regarding the initial array patterns, the simulation results show that the algorithm under the omnidirectional pattern by I2 is robust in most situations, even when there are small differences in power between the desired source signal and the undesired source signals.

V. CONCLUSIONS

In this paper, we have presented the super-exponential blind adaptive beamforming algorithm, which is an extension from the super-exponential blind deconvolution theory and the inverse filter criteria. This extension theoretically considers the presence of noise such that the proposed beamforming algorithm is applicable, regardless of whether the number of independent sources exceeds the system degrees of freedom or not. We have also proved two special conditions under which the performance of the proposed blind algorithm approaches that of the optimal nonblind MMSE beamformer. Simulation results have shown that the proposed algorithm is effective and robust, even when the initial weight vector corresponding to the omnidirectional pattern is used. It should be noted that using the higher order cumulants increases the computational complexity and the sensitivity to the estimation accuracy of the cumulants and decreases the robustness to the model assumptions.

APPENDIX A
PROOF OF THEOREM 1

A. Proof

When $p + q \geq 2$, (20) can be expressed as

$$s_n^{(k)} = \frac{\left(\frac{|\beta_{n_0}|}{|\beta_n|}\right)^{1/(p+q-1)} \left(\frac{|a_n||\beta_n|^{1/(p+q-1)}}{|a_{n_0}||\beta_{n_0}|^{1/(p+q-1)}}\right)^{(p+q)^k} e^{j(p-q)^k \varphi_n + j \psi_n \sum_{l=0}^{k-1} (p-q)^l}}{\sqrt{1 + \sum_{\substack{l=1 \\ l \neq n_0}}^{P+N} \left(\frac{|\beta_{n_0}|}{|\beta_l|}\right)^{2/(p+q-1)} \left(\frac{|a_l||\beta_l|^{1/(p+q-1)}}{|a_{n_0}||\beta_{n_0}|^{1/(p+q-1)}}\right)^{2(p+q)^k}}}. \quad (\text{A1})$$

Under the condition (22), that is

$$\frac{|a_n||\beta_n|^{1/(p+q-1)}}{|a_{n_0}||\beta_{n_0}|^{1/(p+q-1)}} < 1, n = 1, \dots, P + N, n \neq n_0 \quad (\text{A2})$$

where $\beta_n = |\beta_n| e^{j\psi_n}$, $|\beta_n| < \infty$, $n = 1, \dots, P + N$, and the nominator of $|s_n^{(k)}|$ super-exponentially $((p+q)^k)$ converges to zero as $k \rightarrow \infty$, i.e., as in (A3) and (A4), shown at the bottom of the page. Therefore, as $k \rightarrow \infty$

$$|s_n^{(k)}| \rightarrow 0, n = 1, \dots, P + N, n \neq n_0 \quad (\text{A5})$$

$$|s_{n_0}^{(k)}| \rightarrow 1. \quad (\text{A6})$$

Thus, the theorem has been proved. \square

APPENDIX B
PROOF OF THEOREM 2

A. Proof

By substituting (27) and (25) into (23) and by (2), the gain vector at the iteration k can be simplified as

$$\begin{aligned} \mathbf{s}^{(k)} &= \begin{bmatrix} s_1^{(k)} \\ \vdots \\ s_{P+N}^{(k)} \end{bmatrix} = \frac{1}{\Delta^{(k-1)}} \mathbf{H}^T (\mathbf{R}^T)^{-1} \mathbf{H}^* \mathbf{g}^{(k-1)} \\ &= \frac{1}{\Delta^{(k-1)}} \begin{bmatrix} \bar{\mathbf{A}}^T(\Theta) (\mathbf{R}^T)^{-1} \bar{\mathbf{A}}^*(\Theta) & \sigma_{\varpi} \bar{\mathbf{A}}^T(\Theta) (\mathbf{R}^T)^{-1} \\ \sigma_{\varpi} (\mathbf{R}^T)^{-1} \bar{\mathbf{A}}^*(\Theta) & \sigma_{\varpi}^2 (\mathbf{R}^T)^{-1} \end{bmatrix} \\ &\quad \times \begin{bmatrix} \beta_1 \left(s_1^{(k-1)}\right)^p \left(s_1^{(k-1)*}\right)^q \\ \vdots \\ \beta_{P+N} \left(s_{P+N}^{(k-1)}\right)^p \left(s_{P+N}^{(k-1)*}\right)^q \end{bmatrix} \quad (\text{B1}) \end{aligned}$$

where

$$\begin{aligned} \Delta^{(k-1)} &= \sqrt{(\mathbf{d}^{(k-1)})^H (\mathbf{R}^T)^{-1} \mathbf{d}^{(k-1)}} \\ &= \left\| \mathbf{H}^T (\mathbf{H}^* \mathbf{H}^T)^{-1} \mathbf{H}^* \mathbf{g}^{(k-1)} \right\|. \quad (\text{B2}) \end{aligned}$$

Because the cumulants of Gaussian noises are equal to zeros when $p + q \geq 2$ and $P \leq N$, (B1) thus becomes

$$\begin{aligned} \mathbf{s}^{(k)} &= \frac{1}{\Delta^{(k-1)}} \begin{bmatrix} \bar{\mathbf{A}}^T(\Theta) (\mathbf{R}^T)^{-1} \bar{\mathbf{A}}^*(\Theta) \\ \sigma_{\varpi} (\mathbf{R}^T)^{-1} \bar{\mathbf{A}}^*(\Theta) \end{bmatrix} \\ &\quad \times \begin{bmatrix} \beta_1 \left(s_1^{(k-1)}\right)^p \left(s_1^{(k-1)*}\right)^q \\ \vdots \\ \beta_P \left(s_P^{(k-1)}\right)^p \left(s_P^{(k-1)*}\right)^q \end{bmatrix}. \quad (\text{B3}) \end{aligned}$$

With the use of Sherman–Morrison–Woodbury formula [25], under assumption A3, the inversion of \mathbf{R}^T can be written as

$$\begin{aligned} (\mathbf{R}^T)^{-1} &= (\bar{\mathbf{A}}^*(\Theta) \bar{\mathbf{A}}^T(\Theta) + \sigma_{\varpi}^2 \mathbf{I})^{-1} \\ &= \frac{1}{\sigma_{\varpi}^2} \mathbf{I} - \frac{1}{\sigma_{\varpi}^4} \bar{\mathbf{A}}^*(\Theta) \left(\mathbf{I} + \frac{\bar{\mathbf{A}}^T(\Theta) \bar{\mathbf{A}}^*(\Theta)}{\sigma_{\varpi}^2} \right)^{-1} \bar{\mathbf{A}}^T(\Theta) \\ &= \frac{1}{\sigma_{\varpi}^2} \mathbf{I} - \frac{1}{\sigma_{\varpi}^4} \bar{\mathbf{A}}^*(\Theta) \left(\frac{\bar{\mathbf{A}}^T(\Theta) \bar{\mathbf{A}}^*(\Theta)}{\sigma_{\varpi}^2} \right)^{-1} \\ &\quad \times \left(\mathbf{I} + \left(\frac{\bar{\mathbf{A}}^T(\Theta) \bar{\mathbf{A}}^*(\Theta)}{\sigma_{\varpi}^2} \right)^{-1} \right)^{-1} \bar{\mathbf{A}}^T(\Theta). \quad (\text{B4}) \end{aligned}$$

Therefore

$$\begin{aligned} \bar{\mathbf{A}}^T(\Theta) (\mathbf{R}^T)^{-1} \bar{\mathbf{A}}^*(\Theta) &= \frac{\bar{\mathbf{A}}^T(\Theta) \bar{\mathbf{A}}^*(\Theta)}{\sigma_{\varpi}^2} \\ &\quad - \frac{1}{\sigma_{\varpi}^4} \bar{\mathbf{A}}^T(\Theta) \bar{\mathbf{A}}^*(\Theta) \left(\frac{\bar{\mathbf{A}}^T(\Theta) \bar{\mathbf{A}}^*(\Theta)}{\sigma_{\varpi}^2} \right)^{-1} \\ &\quad \times \left(\mathbf{I} + \left(\frac{\bar{\mathbf{A}}^T(\Theta) \bar{\mathbf{A}}^*(\Theta)}{\sigma_{\varpi}^2} \right)^{-1} \right)^{-1} \bar{\mathbf{A}}^T(\Theta) \bar{\mathbf{A}}^*(\Theta) \\ &= \left(\mathbf{I} + \left(\frac{\bar{\mathbf{A}}^T(\Theta) \bar{\mathbf{A}}^*(\Theta)}{\sigma_{\varpi}^2} \right)^{-1} \right)^{-1}. \quad (\text{B5}) \end{aligned}$$

Similarly

$$\begin{aligned} \sigma_{\varpi} (\mathbf{R}^T)^{-1} \bar{\mathbf{A}}^*(\Theta) &= \frac{\bar{\mathbf{A}}^*(\Theta)}{\sigma_{\varpi}} \\ &\quad - \frac{1}{\sigma_{\varpi}^3} \bar{\mathbf{A}}^*(\Theta) \left(\mathbf{I} + \left(\frac{\bar{\mathbf{A}}^T(\Theta) \bar{\mathbf{A}}^*(\Theta)}{\sigma_{\varpi}^2} \right)^{-1} \right)^{-1} \\ &\quad \times \left(\frac{\bar{\mathbf{A}}^T(\Theta) \bar{\mathbf{A}}^*(\Theta)}{\sigma_{\varpi}^2} \right)^{-1} \bar{\mathbf{A}}^T(\Theta) \bar{\mathbf{A}}^*(\Theta) \\ &= \frac{\bar{\mathbf{A}}^*(\Theta)}{\sigma_{\varpi}} \left(\mathbf{I} + \frac{\bar{\mathbf{A}}^T(\Theta) \bar{\mathbf{A}}^*(\Theta)}{\sigma_{\varpi}^2} \right)^{-1}. \quad (\text{B6}) \end{aligned}$$

$$\left(\frac{|\beta_{n_0}|}{|\beta_n|}\right)^{1/(p+q-1)} \left(\frac{|a_n||\beta_n|^{1/(p+q-1)}}{|a_{n_0}||\beta_{n_0}|^{1/(p+q-1)}}\right)^{(p+q)^k} \rightarrow 0 \quad (\text{A3})$$

$$\sqrt{1 + \sum_{\substack{l=1 \\ l \neq n_0}}^{P+N} \left(\frac{|\beta_{n_0}|}{|\beta_l|}\right)^{2/(p+q-1)} \left(\frac{|a_l||\beta_l|^{1/(p+q-1)}}{|a_{n_0}||\beta_{n_0}|^{1/(p+q-1)}}\right)^{2(p+q)^k}} \xrightarrow{k \rightarrow \infty} 1. \quad (\text{A4})$$

When $\sigma_{b_i}^2/\sigma_{\varpi}^2 \rightarrow +\infty, i = 1, \dots, P$, (B5) and (B6) approach the following results:

$$\bar{\mathbf{A}}^T(\Theta) (\mathbf{R}^T)^{-1} \bar{\mathbf{A}}^*(\Theta) \rightarrow \mathbf{I} \quad (\text{B7})$$

$$\sigma_{\varpi} (\mathbf{R}^T)^{-1} \bar{\mathbf{A}}^*(\Theta) \rightarrow \mathbf{0} \quad (\text{B8})$$

and substituting (B7) and (B8) into (B3), we have

$$\mathbf{s}^{(k)} = \frac{1}{\Delta^{(k-1)}} \begin{bmatrix} \beta_1 \left(s_1^{(k-1)} \right)^p \left(s_1^{(k-1)*} \right)^q \\ \vdots \\ \beta_P \left(s_P^{(k-1)} \right)^p \left(s_P^{(k-1)*} \right)^q \\ \mathbf{0} \end{bmatrix} \quad (\text{B9})$$

where

$$\Delta^{(k-1)} = \sqrt{\sum_{n=1}^P \left| \beta_n \left(s_n^{(k-1)} \right)^p \left(s_n^{(k-1)*} \right)^q \right|^2}. \quad (\text{B10})$$

Because (B9) actually corresponds to a special case of (18) considered in Theorem 1, applying the conclusion of Theorem 1 proves Theorem 2. \square

APPENDIX C PROOF OF THEOREM 3

A. Proof

Similar to the proof of Theorem 2, when $P \leq N$ and the associated noises are Gaussian distributed, we rewrite the iterative gain vector of (B3) as (C1), shown at the bottom of the page. Define

$$\bar{\mathbf{A}}_i(\Theta) = [\sigma_{b_1} \mathbf{a}(\theta_1), \dots, \sigma_{b_{i-1}} \mathbf{a}(\theta_{i-1}), \sigma_{b_{i+1}} \mathbf{a}(\theta_{i+1}), \dots, \sigma_{b_P} \mathbf{a}(\theta_P)] \quad (\text{C2})$$

$$\mathbf{R}_i = \bar{\mathbf{A}}_i(\Theta) \bar{\mathbf{A}}_i^H(\Theta) + \sigma_{\varpi}^2 \mathbf{I}. \quad (\text{C3})$$

The condition that $\mathbf{a}^T(\theta_n) \mathbf{a}^*(\theta_{n_0}) = 0, n = 1, \dots, P, n \neq n_0$ implies

$$\bar{\mathbf{A}}_{n_0}^T(\Theta) \mathbf{a}^*(\theta_{n_0}) = \mathbf{0}. \quad (\text{C4})$$

Because of (6), \mathbf{R} can be expressed as

$$\mathbf{R} = \mathbf{R}_{n_0} + \sigma_{b_{n_0}} \mathbf{a}(\theta_{n_0}) (\sigma_{b_{n_0}} \mathbf{a}(\theta_{n_0}))^H. \quad (\text{C5})$$

According to the Sherman–Morrison–Woodbury formula [25], the inversion of \mathbf{R} can be denoted as

$$(\mathbf{R}^T)^{-1} = (\mathbf{R}_{n_0}^T)^{-1} - \frac{(\mathbf{R}_{n_0}^T)^{-1} \sigma_{b_{n_0}} \mathbf{a}^*(\theta_{n_0}) \sigma_{b_{n_0}} \mathbf{a}^T(\theta_{n_0}) (\mathbf{R}_{n_0}^T)^{-1}}{1 + \sigma_{b_{n_0}}^2 \mathbf{a}^T(\theta_{n_0}) (\mathbf{R}_{n_0}^T)^{-1} \mathbf{a}^*(\theta_{n_0})}. \quad (\text{C6})$$

Similarly, by employing the Sherman–Morrison–Woodbury formula to $(\mathbf{R}_{n_0}^T)^{-1}$ and with (C4), we obtain

$$(\mathbf{R}_{n_0}^T)^{-1} \mathbf{a}^*(\theta_{n_0}) = \frac{1}{\sigma_{\varpi}^2} \mathbf{a}^*(\theta_{n_0}). \quad (\text{C7})$$

Therefore

$$(\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_{n_0}) = \frac{\frac{\mathbf{a}^*(\theta_{n_0})}{\sigma_{\varpi}^2}}{1 + \frac{\|\mathbf{a}(\theta_{n_0})\|^2 \sigma_{b_{n_0}}^2}{\sigma_{\varpi}^2}}. \quad (\text{C8})$$

Substituting (C8) into (C1) and by (C4), we have (C9), shown at the bottom of the next page, where

$$\begin{aligned} & \Delta^{(k-1)} \\ &= \left(\left| \mu_0 \beta_{n_0} \left(s_{n_0}^{(k-1)} \right)^p \left(s_{n_0}^{(k-1)*} \right)^q \right|^2 \right. \\ & \quad \left. + \sum_{\substack{l=1 \\ l \neq n_0}}^P \left| \sum_{\substack{i=1 \\ i \neq n_0}}^P \sigma_{b_l} \sigma_{b_i} \mathbf{a}^T(\theta_l) (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_i) \beta_i \left(s_i^{(k-1)} \right)^p \left(s_i^{(k-1)*} \right)^q \right|^2 \right. \\ & \quad \left. + \left\| \sum_{i=1}^P \sigma_{\varpi} \sigma_{b_i} (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_i) \beta_i \left(s_i^{(k-1)} \right)^p \left(s_i^{(k-1)*} \right)^q \right\|^2 \right)^{1/2} \end{aligned} \quad (\text{C9})$$

$$\begin{aligned} & \mu_0 \\ &= \frac{\frac{\|\mathbf{a}(\theta_{n_0})\|^2 \sigma_{b_{n_0}}^2}{\sigma_{\varpi}^2}}{1 + \frac{\|\mathbf{a}(\theta_{n_0})\|^2 \sigma_{b_{n_0}}^2}{\sigma_{\varpi}^2}} < 1. \end{aligned} \quad (\text{C11})$$

According to (C9), the following inequalities stand:

$$\begin{aligned} & \left| s_n^{(k)} \right| \\ & \leq \frac{1}{\Delta^{(k-1)}} \sum_{\substack{i=1 \\ i \neq n_0}}^P \left| \sigma_{b_n} \sigma_{b_i} \mathbf{a}^T(\theta_n) (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_i) \right| \left| \beta_i \right| \left| s_i^{(k-1)} \right|^{p+q} \\ & \quad n = 1, \dots, P, n \neq n_0 \end{aligned} \quad (\text{C12a})$$

$$\begin{aligned} & \left| s_{n_0}^{(k)} \right| \\ &= \frac{1}{\Delta^{(k-1)}} \mu_0 \left| \beta_{n_0} \right| \left| s_{n_0}^{(k-1)} \right|^{p+q} \end{aligned} \quad (\text{C12b})$$

$$\begin{aligned} & \left| s_{P+l}^{(k)} \right| \\ &= \frac{1}{\Delta^{(k-1)}} \sum_{i=1}^P \left(\sigma_{\varpi} \sigma_{b_i} (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_i) \right)_l \beta_i \left(s_i^{(k-1)} \right)^p \left(s_i^{(k-1)*} \right)^q \\ & \quad l = 1, \dots, N. \end{aligned} \quad (\text{C12c})$$

It can be seen that $s_{P+l}^{(k)}, l = 1, \dots, N$ depends on $s_n^{(k)}, n = 1, \dots, P$, whereas we are only concerned with the iterative behavior of $s_n^{(k)}, n = 1, \dots, P$. Let $\left| s_{m_0}^{(0)} \right| = \max \left\{ \left| s_n^{(0)} \right|, n = 1, \dots, P, n \neq n_0 \right\}$, i.e., for $n = 1, \dots, P$

$$\left| s_n^{(0)} \right| \leq \left| s_{m_0}^{(0)} \right|. \quad (\text{C13})$$

$$\mathbf{s}^{(k)} = \frac{1}{\Delta^{(k-1)}} \begin{bmatrix} \sigma_{b_1}^2 \mathbf{a}^T(\theta_1) (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_1) & \cdots & \sigma_{b_1} \sigma_{b_P} \mathbf{a}^T(\theta_1) (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_P) \\ \vdots & \ddots & \vdots \\ \sigma_{b_P} \sigma_{b_1} \mathbf{a}^T(\theta_P) (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_1) & \cdots & \sigma_{b_P}^2 \mathbf{a}^T(\theta_P) (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_P) \\ \sigma_{b_1} \sigma_{\varpi} (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_1) & \cdots & \sigma_{b_P} \sigma_{\varpi} (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_P) \end{bmatrix} \times \begin{bmatrix} \beta_1 \left(s_1^{(k-1)} \right)^p \left(s_1^{(k-1)*} \right)^q \\ \vdots \\ \beta_P \left(s_P^{(k-1)} \right)^p \left(s_P^{(k-1)*} \right)^q \end{bmatrix}. \quad (\text{C1})$$

By (C13), (C12a) and (C12b), the following relations can be easily inferred:

$$\begin{aligned} & \left| s_n^{(1)} \right| \\ & \leq \frac{1}{\Delta^{(0)}} \left| s_{m_0}^{(0)} \right|^{p+q} \sum_{\substack{i=1 \\ i \neq n_0}}^P \left| \sigma_{b_n} \sigma_{b_i} \mathbf{a}^T(\theta_n) (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_i) \right| |\beta_i| \\ & \quad n = 1, \dots, P, n \neq n_0 \end{aligned} \quad (\text{C14a})$$

$$\begin{aligned} & \left| s_{n_0}^{(1)} \right| \\ & = \frac{1}{\Delta^{(0)}} \mu_0 |\beta_{n_0}| \left| s_{n_0}^{(0)} \right|^{p+q} \\ & \quad \Delta^{(0)} \end{aligned} \quad (\text{C14b})$$

$$\begin{aligned} & = \left(|\beta_{n_0}|^2 \left| s_{n_0}^{(0)} \right|^{2(p+q)} \right. \\ & \quad \left. + \sum_{\substack{l=1 \\ l \neq n_0}}^P \left| \sum_{\substack{i=1 \\ i \neq n_0}}^P \sigma_{b_l} \sigma_{b_i} \mathbf{a}^T(\theta_l) (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_i) \beta_i \left(s_i^{(0)} \right)^p \left(s_i^{(0)*} \right)^q \right|^2 \right. \\ & \quad \left. + \left\| \sum_{i=1}^P \sigma_{\omega} \sigma_{b_i} (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_i) \beta_i \left(s_i^{(0)} \right)^p \left(s_i^{(0)*} \right)^q \right\|^2 \right)^{1/2}. \end{aligned} \quad (\text{C15})$$

Define

$$\begin{aligned} \tilde{\beta}_n & = \sum_{\substack{i=1 \\ i \neq n_0}}^P \left| \sigma_{b_n} \sigma_{b_i} \mathbf{a}^T(\theta_n) (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_i) \right| |\beta_i| \\ & \quad n = 1, \dots, P, n \neq n_0 \end{aligned} \quad (\text{C16a})$$

$$\tilde{\beta}_{n_0} = \mu_0 |\beta_{n_0}| \quad (\text{C16b})$$

$$\tilde{\beta}_{r_0} = \max \left\{ \tilde{\beta}_n, n = 1, \dots, P, n \neq n_0 \right\}, 1 \leq r_0 \leq P \quad (\text{C16c})$$

so we have

$$\left| s_n^{(1)} \right| \leq \frac{1}{\Delta^{(0)}} \tilde{\beta}_n \left| s_{m_0}^{(0)} \right|^{p+q} \leq \frac{1}{\Delta^{(0)}} \tilde{\beta}_{r_0} \left| s_{m_0}^{(0)} \right|^{p+q}, n = 1, \dots, P. \quad (\text{C17})$$

$$\begin{aligned} \left| s_n^{(2)} \right| & \leq \frac{1}{\Delta^{(1)}} \sum_{\substack{i=1 \\ i \neq n_0}}^P \left| \sigma_{b_n} \sigma_{b_i} \mathbf{a}^T(\theta_n) (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_i) \right| |\beta_i| \left| s_i^{(1)} \right|^{p+q} \\ & \leq \frac{1}{\Delta^{(1)}} \tilde{\beta}_n \left(\tilde{\beta}_{r_0} \left| s_{m_0}^{(0)} \right|^{p+q} \right)^{p+q} \\ & \leq \frac{1}{\Delta^{(1)}} \tilde{\beta}_{r_0}^{p+q+1} \left| s_{m_0}^{(0)} \right|^{(p+q)^2}, n = 1, \dots, P, n \neq n_0 \end{aligned} \quad (\text{C18a})$$

$$\left| s_{n_0}^{(2)} \right| = \frac{1}{\Delta^{(1)}} \tilde{\beta}_{n_0}^{p+q+1} \left| s_{n_0}^{(0)} \right|^{(p+q)^2} \leq \frac{1}{\Delta^{(1)}} \tilde{\beta}_{r_0}^{p+q+1} \left| s_{m_0}^{(0)} \right|^{(p+q)^2}. \quad (\text{C18b})$$

Repeatedly, we have the iterative equations

$$\begin{aligned} & \left| s_{n_0}^{(k)} \right| \\ & = \frac{1}{\Delta^{(k-1)}} \tilde{\beta}_{n_0}^{((p+q)^k - 1)/(p+q-1)} \left| s_{n_0}^{(0)} \right|^{(p+q)^k} \\ & \leq \frac{1}{\Delta^{(k-1)}} \tilde{\beta}_{r_0}^{((p+q)^k - 1)/(p+q-1)} \left| s_{m_0}^{(0)} \right|^{(p+q)^k} \end{aligned} \quad (\text{C19a})$$

$$\begin{aligned} & \left| s_n^{(k)} \right| \\ & \leq \frac{1}{\Delta^{(k-1)}} \tilde{\beta}_{r_0}^{((p+q)^k - 1)/(p+q-1)} \left| s_{m_0}^{(0)} \right|^{(p+q)^k} \\ & \quad n = 1, \dots, P + N, n \neq n_0 \end{aligned} \quad (\text{C19b})$$

$$\begin{aligned} & \Delta^{(k-1)} \\ & = \left(\tilde{\beta}_{n_0}^{2 \times ((p+q)^k - 1)/(p+q-1)} \left| s_{n_0}^{(0)} \right|^{2(p+q)^k} \right. \\ & \quad \left. + \sum_{\substack{l=1 \\ l \neq n_0}}^P \left| \sum_{\substack{i=1 \\ i \neq n_0}}^P \sigma_{b_l} \sigma_{b_i} \mathbf{a}^T(\theta_l) (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_i) \beta_i \left(s_i^{(k-1)} \right)^p \left(s_i^{(k-1)*} \right)^q \right|^2 \right. \\ & \quad \left. + \left\| \sum_{i=1}^P \sigma_{\omega} \sigma_{b_i} (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_i) \beta_i \left(s_i^{(k-1)} \right)^p \left(s_i^{(k-1)*} \right)^q \right\|^2 \right)^{1/2}. \end{aligned} \quad (\text{C19c})$$

$$\mathbf{s}^{(k)} = \frac{1}{\Delta^{(k-1)}} \begin{bmatrix} \sum_{\substack{i=1 \\ i \neq n_0}}^P \sigma_{b_1} \sigma_{b_i} \mathbf{a}^T(\theta_1) (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_i) \beta_i \left(s_i^{(k-1)} \right)^p \left(s_i^{(k-1)*} \right)^q \\ \vdots \\ \sum_{\substack{i=1 \\ i \neq n_0}}^P \sigma_{b_{n_0-1}} \sigma_{b_i} \mathbf{a}^T(\theta_{n_0-1}) (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_i) \beta_i \left(s_i^{(k-1)} \right)^p \left(s_i^{(k-1)*} \right)^q \\ \mu_0 \beta_{n_0} \left(s_{n_0}^{(k-1)} \right)^p \left(s_{n_0}^{(k-1)*} \right)^q \\ \sum_{\substack{i=1 \\ i \neq n_0}}^P \sigma_{b_{n_0+1}} \sigma_{b_i} \mathbf{a}^T(\theta_{n_0+1}) (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_i) \beta_i \left(s_i^{(k-1)} \right)^p \left(s_i^{(k-1)*} \right)^q \\ \vdots \\ \sum_{\substack{i=1 \\ i \neq n_0}}^P \sigma_{b_P} \sigma_{b_i} \mathbf{a}^T(\theta_P) (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_i) \beta_i \left(s_i^{(k-1)} \right)^p \left(s_i^{(k-1)*} \right)^q \\ \sum_{i=1}^P \sigma_{\omega} \sigma_{b_i} (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_i) \beta_i \left(s_i^{(k-1)} \right)^p \left(s_i^{(k-1)*} \right)^q \end{bmatrix} \quad (\text{C9})$$

Therefore

$$\left| s_n^{(k)} \right| \leq \left(\frac{\tilde{\beta}_{n_0}}{\tilde{\beta}_{r_0}} \right)^{1/(p+q-1)} \left(\frac{\tilde{\beta}_{r_0}^{1/(p+q-1)} |s_{m_0}^{(0)}|}{\tilde{\beta}_{n_0}^{1/(p+q-1)} |s_{n_0}^{(0)}|} \right)^{(p+q)^k} \quad n = 1, \dots, P, n \neq n_0. \quad (\text{C20})$$

When $\tilde{\beta}_{n_0}^{1/(p+q-1)} |s_{n_0}^{(0)}| > \tilde{\beta}_{r_0}^{1/(p+q-1)} |s_{m_0}^{(0)}|$, as $k \rightarrow \infty$, the following term will super-exponentially converges to zero:

$$\left(\frac{\tilde{\beta}_{r_0}^{1/(p+q-1)} |s_{m_0}^{(0)}|}{\tilde{\beta}_{n_0}^{1/(p+q-1)} |s_{n_0}^{(0)}|} \right)^{(p+q)^k} \xrightarrow{k \rightarrow \infty} 0. \quad (\text{C21})$$

Because of (C8), (C21) implies that

$$\left| s_n^{(k)} \right| \xrightarrow{k \rightarrow \infty} 0, n = 1, \dots, P, n \neq n_0 \quad (\text{C22a})$$

$$\left| s_{n_0}^{(k)} \right| \xrightarrow{k \rightarrow \infty} \frac{1}{\sqrt{1 + \frac{1}{\frac{\|\mathbf{a}(\theta_{n_0})\|^2 \sigma_{b_{n_0}}^2}{\sigma_{\tilde{\mathbf{c}}}^2}}}}}. \quad (\text{C22b})$$

Therefore, from (C9), we have

$$\begin{aligned} \mathbf{s}^{(k)} \xrightarrow{k \rightarrow \infty} \sigma_{b_{n_0}} \mathbf{H}^T (\mathbf{R}^T)^{-1} \mathbf{a}^*(\theta_{n_0}) \\ = \frac{\mathbf{H}^T \mathbf{a}^*(\theta_{n_0}) \sigma_{b_{n_0}}}{\sigma_{\tilde{\mathbf{c}}}^2} \\ = \frac{\mathbf{H}^T \mathbf{a}^*(\theta_{n_0}) \sigma_{b_{n_0}}}{1 + \frac{\|\mathbf{a}(\theta_{n_0})\|^2 \sigma_{b_{n_0}}^2}{\sigma_{\tilde{\mathbf{c}}}^2}}. \end{aligned} \quad (\text{C23})$$

It can be easily seen from (C23) that as k increases, $\mathbf{s}^{(k)}$ approaches the gain vector of the optimal nonblind MMSE beamformer up to a complex rotating phase factor. Thus, Theorem 3 has been proved. \square

APPENDIX D

PROOF OF THEOREM 4

A. Proof

When the super-exponential iterative algorithm converges, (35a)–(35c) can be written as

$$\mathbf{c} = \frac{(\mathbf{R}^T)^{-1} \mathbf{d}}{\sqrt{\mathbf{d}^H (\mathbf{R}^T)^{-1} \mathbf{d}}} \quad (\text{D1a})$$

$$\mathbf{d} = \text{cum}(z(t) : m; z^*(t) : m-1; \mathbf{x}^*(t)) \quad (\text{D1b})$$

$$z(t) = \mathbf{c}^T \mathbf{x}(t). \quad (\text{D1c})$$

Substituting (D1a) into (59), we have

$$\begin{aligned} \frac{\partial J_{2m}(\mathbf{c})}{\partial \mathbf{c}} &= m \times \left(\text{sgn}(C_{m,m}\{z(t)\}) \right. \\ &\times \text{cum}\{z(t) : m, z^*(t) : m-1, \mathbf{x}^*(t)\} - \frac{|C_{m,m}\{z(t)\}|}{\sqrt{\mathbf{d}^H (\mathbf{R}^T)^{-1} \mathbf{d}}} \mathbf{d} \left. \right) \end{aligned} \quad (\text{D2})$$

where $\mathbf{c}^T \mathbf{R} \mathbf{c} = 1$ is used. Based on the linearity property of the cumulants [7] and with (D1b) and (D1c)

$$\begin{aligned} C_{m,m}\{z(t)\} &= \mathbf{c}^H \text{cum}\{z(t) : m, z^*(t) : m-1, \mathbf{x}^*(t)\} \\ &= \mathbf{c}^H \mathbf{d} = \frac{\mathbf{d}^H (\mathbf{R}^T)^{-1} \mathbf{d}}{\sqrt{\mathbf{d}^H (\mathbf{R}^T)^{-1} \mathbf{d}}} > 0 \end{aligned} \quad (\text{D3})$$

which means

$$\text{sgn}(C_{m,m}\{z(t)\}) = +1. \quad (\text{D4})$$

Therefore

$$\frac{\partial J_{2m}(\mathbf{c})}{\partial \mathbf{c}} = \mathbf{0} \quad (\text{D5})$$

and Theorem 4 has been proved. \square

APPENDIX E

CUMULANTS $\text{cum}(z(t) : m; z^*(t) : m)$ FOR $m = 2, 3, 4$

Because analytical signals (complex envelope) in communications, radar, etc., are often of symmetrical distributions with odd order cumulants equal to zero, the following are formulae for computing $2m$ th-order cumulants of $z(t)$ for $m = 2, 3$, and 4 as defined in [7]:

$$\begin{aligned} &\text{cum}(z(t) : 2; z^*(t) : 2) \\ &= E(|z(t)|^4) - 2(E(|z(t)|^2))^2 \end{aligned} \quad (\text{E1a})$$

$$\begin{aligned} &\text{cum}(z(t) : 2; z^*(t) : 2; \mathbf{x}^*(t)) \\ &= E(|z(t)|^2 z(t) \mathbf{x}^*(t)) - 2E(|z(t)|^2) E(z(t) \mathbf{x}^*(t)) \end{aligned} \quad (\text{E1b})$$

$$\begin{aligned} &\text{cum}(z(t) : 3; z^*(t) : 3) \\ &= E(|z(t)|^6) - 9E(|z(t)|^4) E(|z(t)|^2) + 12(E(|z(t)|^2))^3 \end{aligned} \quad (\text{E2a})$$

$$\begin{aligned} &\text{cum}(z(t) : 3; z^*(t) : 2; \mathbf{x}^*(t)) \\ &= E(|z(t)|^4 z(t) \mathbf{x}^*(t)) - 3E(|z(t)|^4) E(z(t) \mathbf{x}^*(t)) \\ &\quad - 6E(|z(t)|^2) E(|z(t)|^2 z(t) \mathbf{x}^*(t)) \\ &\quad + 12(E(|z(t)|^2))^2 E(z(t) \mathbf{x}^*(t)) \end{aligned} \quad (\text{E2b})$$

$$\begin{aligned} &\text{cum}(z(t) : 4; z^*(t) : 4) \\ &= E(|z(t)|^8) - 16E(|z(t)|^6) E(|z(t)|^2) - 18(E(|z(t)|^4))^2 \\ &\quad + 144E(|z(t)|^4) (E(|z(t)|^2))^2 - 144(E(|z(t)|^2))^4 \end{aligned} \quad (\text{E3a})$$

$$\begin{aligned} &\text{cum}(z(t) : 4; z^*(t) : 3; \mathbf{x}^*(t)) \\ &= E(|z(t)|^6 z(t) \mathbf{x}^*(t)) - 4E(|z(t)|^6) E(z(t) \mathbf{x}^*(t)) \\ &\quad - 12E(|z(t)|^2) E(|z(t)|^4 z(t) \mathbf{x}^*(t)) \\ &\quad - 18E(|z(t)|^4) E(|z(t)|^2 z(t) \mathbf{x}^*(t)) \\ &\quad + 72E(|z(t)|^4) E(|z(t)|^2) E(z(t) \mathbf{x}^*(t)) \\ &\quad + 72(E(|z(t)|^2))^2 E(|z(t)|^2 z(t) \mathbf{x}^*(t)) \\ &\quad - 144(E(|z(t)|^2))^3 E(z(t) \mathbf{x}^*(t)). \end{aligned} \quad (\text{E3b})$$

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