# Approximated stochastic realization and model reduction methods applied to array processing by means of state space models 

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## ABSTRACT

The aim of this paper is to present new methods for passive array processing. The basic idea consists in using a state-space modelling of the sensors output. This paper copes with basic problems as: unknown noise correlations, approximation by a Toplitz matrix of lower rank, detection of small sources. The presented methods represent considerable improvements with respect to the usual ones and furthermore are quite feasible, some statistical results illustrate these claims.

## 1. Introduction

A modelling of the sensors output of an equally spaced linear array by a linear system is shown to be a powerful tool, in particular to take into account fundamental hypotheses about sources and noise. A general frame for approximated stochastic realization methods is presented and connections with existing work and array processing are especially considered.

A state space modelling of the sensors output uses sources and noise properties in a summarizing way and is a powerful means to separate the noise and sources contributions even in the case of an unknown correlated noise. This last property represents a considerable improvement against usual high resolution methods. For that purpose two methods of approximated stochastic realization are considered, the links between these ones and array processing are detailled.

However a fundamental problem remains: find an estimation of the sources subspace that preserves the plane wave structure. We are now facing the basic problem which consists in extracting a low rank matrix (of sources) from an estimated (full rank) one preserving the special structure (Tœplitz or Hankel) induced by plane wave and spatial stationnarity hypotheses. The Adamjan-Arov-Krein theory is the theoretical frame, the initial problem of functionnal analysis is translated into another simpler and finite dimensional problem by the use of the properties of finite structures and balanced realizations. A further advantage consists in using a $\mathrm{L}_{\infty}$-norm rather than an $\mathrm{L}_{2}$-norm ,much more sensitive to local variations of the spectra.

Furthermore, simulation results present comparisons between state space approaches and. usual high resolution methods (as eigenvectors methods) and try to precise their advantages.

## 2. Additive white noise

This case is especially simple but allows us to precise our definitions. In the whole paper, the sensors output of a linear equispaced array constituted by n sensors are considered at a given frequency (omitted for the sequel). Under plane wave assumption, the output at the i -th sensor may be written as:

$$
y_{i}=\sum_{k=1}^{m} \alpha_{k} \cdot \exp \left(-j \cdot \varphi_{k}\right)+v_{i} \quad \text { (1) } \quad\left(j^{2}=-1\right)
$$

where $m$ is the number of sources, $\alpha_{k}$ and $\varphi_{k}$ represent respectively the random level and the deterministic phase of the $k$-th source, and $\left\{v_{i}\right\}$ is the additive noise.

A rank $m$ linear state space model of (1) is given in the case of an additive white noise by [1]
$\left\{\begin{array}{l}X_{i+1}=F . X_{i} \\ y_{i}=h^{*} \cdot \mathbf{X}_{i}+v_{i}\end{array}\right.$
F: (mxm)
$y_{i=} h^{*} \cdot X_{i}+v_{i} \quad$ (2) $h^{*}:(1 \times m)$
where $F, h$ are deterministic parameters, $X_{i}$ is the state vector at sensor $i$ and $\left\{v_{i}\right\}$ is an additive white noise. The matrix $F$ translates the phasing properties of the propagation while $h$ represents its power effects. The phases $\varphi_{k}$ are determined from the eigenvalues of the transition matrix $F$ which is similar to $\operatorname{diag}\left[\exp \left(-j . \varphi_{k}\right)\right]$.

Let $\Gamma_{p}$ be the covariances of the sensors output, i.e. : $\Gamma_{p}=E\left(y_{i+p} y_{i}{ }^{*}\right)$ then

$$
\begin{equation*}
\Gamma_{p}=h^{*} \cdot \mathbf{F}^{p} \cdot \mathbf{P} \cdot \mathbf{h}+\mathbf{b}^{2} \cdot \delta(p) \tag{3}
\end{equation*}
$$

with $P=E\left(X_{i} \cdot X_{i}{ }^{*}\right), b^{2}=E\left(\left|v_{i}\right|^{2}\right)$.
Furthermore, let $\Gamma_{+}$be the exact covariance matrix of the sensors output, then

$$
\begin{equation*}
\Gamma_{+}=0 . P \cdot O^{*}+b^{2} \cdot I_{n} \tag{4}
\end{equation*}
$$

where $\mathbf{O}$ is the ( $\mathrm{n} \times \mathrm{n}$ ) observability matrix


The following property can be casily deduced from the structure of the observability matrix:

$$
\begin{equation*}
0^{\uparrow} . F=0^{\downarrow} \tag{5}
\end{equation*}
$$

( $\uparrow$ : select the ( $n-1$ ) last rows, $\downarrow$ : select the ( $n-1$ ) first rows, $0^{\uparrow}$ and $0^{\downarrow}$ are ( $\left.(\mathrm{n}-1) \times \mathrm{m}\right)$ matrices), and therefore;

$$
\begin{equation*}
F=\left(0^{\uparrow}\right)^{\#} \cdot 0^{\downarrow} \tag{6}
\end{equation*}
$$

(\#: denoting pseudo-inverse matrix)
Elsewhere, the covariance matrix can also be written in terms of array processing as:

$$
\begin{equation*}
\Gamma_{+}=D \cdot \gamma \cdot D^{*}+b^{2} \cdot I_{n} \tag{7}
\end{equation*}
$$

with $\mathbf{D}$ : source steering matrix
then the plane wave hypothesis yields: $D^{\uparrow} \cdot F=D^{\downarrow}$ or $F=\left(D^{\uparrow}\right)^{\#} \cdot D^{\downarrow}$
leading to property 1 [2], and a geometric meaning of the observability matrix.

Prop.1: in the asymptotic case, the eigenvalues of the transition matrix $F$ are identical to these of the matrix $\left[\left(U_{1}{ }^{\uparrow}\right)^{\#} . U_{1}{ }^{\downarrow}\right]$ where $U_{1}$ is the ( $\mathrm{n} \times \mathrm{m}$ ) matrix
formed froni the eigenvectors corresponding to the $m$ largest eigenvalues of $\Gamma_{+}$.

In practical situation however, the matrix $\mathrm{U}_{1}$ is replaced by $\widehat{U}_{1}$ an estimate (from $\widehat{\Gamma}_{+}$), therefore the presented method "forgets" the special structure of the obscrvability matrix.
3. Approximated realization methods, additive correlated noise.
In the general case, it is possible to take into account the noise correlations by modelling the sensors output as the outputs of a minimum phase innovation model, i.e. :

$$
\left\{\begin{array}{l}
\mathbf{x}_{i+1}=F \cdot X_{i}+\mathbf{T} \cdot v_{i}  \tag{9}\\
y_{i}=h^{*} \cdot \mathbf{x}_{i}+v_{i}
\end{array}\right.
$$

Then the approximated realization methods rely upon consideration of two vectors $Y_{+}$and $Y_{\text {. (future }}$ and past at sensor $n / 2$ ) defined as follows:

$$
\left\{\begin{array}{l}
\mathbf{Y}_{-} \equiv\left(y_{n_{2}}, y_{n_{2}-1}, \ldots, y_{1}\right)  \tag{10}\\
\mathbf{Y}_{+} \equiv\left(y_{n_{2}+1}, y_{n_{2}+2}, \ldots, y_{n}\right)
\end{array}\right.
$$

The aim of the methods consists in summarizing the past in order to obtain the more efficient (for a given criterion) prediction of the future. The stochastic approach [3] determines a state vector $\mathbf{X}$ (of a given dimension) which sums up the more pertinent part of $Y_{-}$in order to predict $Y_{+}$, i.e.:

$$
\begin{equation*}
X=A . Y . \text { with } A(p \times n) \tag{11}
\end{equation*}
$$

Theoretically one has:

$$
\begin{equation*}
A=\left[\mathbf{T},\left(\mathbf{F}-\mathbf{T} \cdot \mathbf{h}^{*}\right) \cdot \mathbf{T}, \ldots,\left(\mathbf{F}-\mathbf{T} \cdot \mathbf{h}^{*}\right)^{\mathbf{n} \mathbf{2}^{-1}} \cdot \mathbf{T}\right] \tag{12}
\end{equation*}
$$

Furthermore, in the case of a markovian system, the orthogonal projection of the future onto the past is given by:
$\mathbf{Y}_{+} \mid \mathbf{Y}_{-}=\mathbf{0} . \mathbf{X} \quad(\mathbf{O}:$ observability matrix) (13)
The determination of A allows us to estimate 0 and therefore the system parameters. The crucial step of these methods consists in the estimation of the matrix A, for that purpose various methods can be considered.
3.1. The predictive efficiency criterion (AK method)

The matrix $A$ is the matrix which minimizes the variance of the predictor error, i.e. [4] :
$\operatorname{Min}\left\{\operatorname{tr}\left[\operatorname{cov}\left(\mathbf{Y}_{+}-\mathbf{Y}_{+} \mid \mathbf{X}\right)\right]\right\}$
$\mathbf{X}=\mathbf{A} . \mathbf{Y}_{\text {- }}$
yielding to:

$$
\operatorname{Min}\left\{\operatorname{tr}\left[\Gamma_{+}-H \cdot A^{*} \cdot\left(A \cdot \Gamma_{-} \cdot A^{*}\right) \cdot A \cdot H^{*}\right]\right\}
$$

A
$\left(\Gamma_{+}=\mathrm{E}\left(\mathbf{Y}_{+} \cdot \mathbf{Y}_{+}{ }^{*}\right), \Gamma_{-} \equiv \mathrm{E}\left(\mathbf{Y}_{-} . \mathbf{Y}_{-}^{*}\right), \mathbf{H} \equiv \mathrm{E}\left(\mathbf{Y}_{+} . \mathbf{Y}_{-}^{*}\right)\right)$
forgetting special structure of $A$, the solution of (14) is obtained by elementary algebra, i.e.:

Prop. 2: the solution of the problem (14) has the general form: $\quad \mathrm{A}=\mathrm{B} \cdot \Sigma_{1}{ }^{-1} \cdot \mathrm{U}_{1}^{*} \cdot \mathrm{H} \cdot \Gamma_{-}^{-1} \quad$ (15) where $B$ is any ( $p \times p$ ) invertible matrix and ( $\Sigma_{1}{ }^{2}, U_{1}$ ) are the $p$ principal components of the matrix H. . ${ }_{-}^{-1} \cdot H^{*}$.

### 3.2. Information criterion method (DP method)

This approach is based upon the fact that the whole information contained in $\mathbf{Y}_{+}$which is explained by $\mathbf{Y}_{.}$ can be expressed in terms of few parameters $\left\{\sigma_{k}\right\}_{k=1, p}$, called the canonical correlation coefficients:

$$
\mathcal{L}\left(\mathbf{Y}_{+} \mathbf{Y}_{-}\right)=-\sum_{\mathbf{k}=1}^{\mathbf{R}} \log \left(1-\boldsymbol{\sigma}_{\mathbf{k}}^{\mathbf{2}}\right)
$$

Denoting $\mathcal{L}\left(Y_{+}, X\right)$ the mutual information between $\mathbf{Y}_{+}$and $\mathbf{X}$ defined as follows:
$\boldsymbol{L}\left(\mathbf{Y}_{+}, \mathbf{X}\right)=\boldsymbol{H}\left(\mathbf{Y}_{+}\right)+\boldsymbol{H}(\mathbf{X})-\boldsymbol{H}\left\lceil\frac{\mathbf{Y}_{+}}{\mathbf{X}}\right\rceil$
where $\mathcal{H}(X)$ is the entropy of the gaussian vector $X$, then the estimation of $X=A . Y$. amouts to minimize the following criterion [5]:

$$
\begin{equation*}
\operatorname{Min}\left\{\operatorname{det}\left[\Gamma_{+}-H \cdot A^{*} \cdot\left(A \cdot \Gamma_{-} \cdot A^{*}\right), A \cdot H^{*}\right]\right\} \tag{17}
\end{equation*}
$$

A
Note that both criteria differ only by the functional (tr or det), the minimization of this criterion is achieved by means of elementary algebra leading to:

Prop. 3: the solution of the problem (17) is given by:

$$
\begin{equation*}
A=B \cdot V_{1}^{*} \cdot \Gamma_{-}^{-1 / 2} \tag{18}
\end{equation*}
$$

where $V_{1}$ is the matrix formed from the p right principal singular vectors of the matrix $\left\{\Gamma_{+}{ }^{-1 / 2} \cdot \mathrm{H} \cdot\left(\Gamma_{-}^{1 / 2}\right)^{*}\right\}$.

Therefore the solution of (17) is nothing but the solution obtained by maximization of canonical correlations as advocated by Desai and Pal [5].

The matrix A being estimated by means of Prop. 2 or $3, F$ is staightforwardly deduced by two ways since:

- $0 . A=H . \Gamma_{-}^{-1}$, and $0^{\uparrow} . F=0^{\downarrow}$.
$-X_{i+1} \mid X_{i}=F . \mathbf{X}_{i}$.
The links with classical methods of array processing are much more subtle than in $\$ 2$. In the white noise case, one has [6]: H. $\Gamma_{-}^{-1} \cdot \mathbf{H}^{*}=\mathbf{D} . \Delta . \mathbf{D}^{*}$. therefore the matrix $H . \Gamma_{-}^{-1} \cdot H^{*}$ plays the role of $\Gamma_{+}$in Prop. 1; but this is not very relevant of the general case. For that purpose, it is more appropriate to consider the predictive efficiency criterion in terms of the intersection of two linear subspaces, i.e. [6] :

Prop. 4: the rows of the matrix $A$ approximate a basis of Range $\left(\mathbf{H}^{*} . H^{\prime}\right) \cap \operatorname{Range}\left(\Gamma_{-}\right)$.

Prop. 4 means that the rows of $A$ minimize the principal angles between Range( $\Gamma$.) and Range ( $\mathrm{H}^{*} . \mathrm{H}$ ).

Under the assumption of a shortly correlated additive noise, the matrix $H$ is perturbed only by a triangular matrix, the vectors of Range( $\left.H^{*} . H\right) \cap$ Range( $\Gamma$.) are mainly related to source parameters.

To sum up, the presented methods are efficient to estimate the source bearings in presence of an additive noise with unknown correlations, however they do not use fundamentally the plane wave hypothesis for the estimation of $A$. We shall now cope with this problem.

## 4. Optimal realization methods.

4.1. A general framework

The approximation of a linear system by a lower order one represents an important part of control systems litterature. However this field seems to be not very relevant for array processing, but the following general framework permits us to use the results of control systems theory.

The array outputs are modelled by a maximal order parametric model (e.g.: an AR model of order $n$ ), note that there is no infurmation loss since it is possible to construct the corresponding covariances ( $\Gamma$ ) from them. It is simply another representation of the covariances. A state space model can be easily obtained from the model and it is then possible to use the optimal model reduction methods [7].

We want an exact approximation, preserving the structure and avoiding the classical least squares approximations (leading generally to principal component analysis). This is basically a problem of approximation. In fact, the theorem of Adamjan, Arov and Krein (AAK) is the foundation of such an
approximation, although it is mainly theoretical and in particular involves infinite matrices. Fortunately it is possible to translate this problem into another finite dimensional problem (by use of finite structure) and to present a practical algorithm.

### 4.2. Problem formulation

Given a scalar transfert function (e.g. :
$\left.f(z)=h^{*} \cdot(z, I-F)^{-1} \cdot \mathbf{T}+1\right)$ we seek an approximation (for a certain norm) of this function by a lower order system. The use of $L_{2}$-norm is classical but is not as powerful as $\mathrm{L}_{\infty}$-norm for our problem where sources may have very different powers. The inclusion of special structure leads to the following problem:

Pb 1: given $f \in \mathrm{~L}_{\infty}$ (defined on $\mathrm{C}(0,1)$ ) and a positive integer $k$ find

$$
\begin{equation*}
\inf \left\{\|f-\varphi\|_{\infty}: \varphi \in \mathrm{H}_{\infty, \mathbf{k}}\right\} \tag{19}
\end{equation*}
$$

and a function $\varphi_{k}$ (of $H_{\infty, k}$ ) for which the infinimum is attained; $H_{\infty, k}$ being the set of meromorphic functions in $\mathbf{L}_{\infty}$ which can be written as:

$$
\begin{equation*}
\varphi(z)=\frac{g(z)}{\left(z-\alpha_{1}\right) \ldots\left(z-\alpha_{n}\right)} \tag{20}
\end{equation*}
$$

where : $\left(\alpha_{1}, \ldots, \alpha_{k}\right) \in D(0,1)$ and $g \in H_{\infty} ; H_{\infty}$ being the subspace of $L_{\infty}$ for which $c(n)=0$ for $n<0$ (c(n) being the $n$-th Fourier coefficient of $f$ w.r. to the complete basis ( $\mathrm{z}^{\mathrm{n}}$ ), $\mathrm{z}=\exp (\mathrm{j} . \boldsymbol{\theta})$ ).

This problem has been solved by AAK, its solution consists in using the Hankel matrix of a function $f$ w.r. to the basis $z^{-1}, z^{-2}, \ldots$ in $L_{2} \Theta H_{2}$ defined as:
then a first result is available:
Res. I: $\left\|H_{f}\right\|_{\infty}=\operatorname{dist}_{L_{\infty}}\left(\mathbf{f}, \mathrm{H}_{\infty}\right)$
the function of $H_{\infty}$ which minimizes that distance being given by Nehari's theorem.

Furthermore, the Kronecker's theorem allows us to precise the problem:

Kronecker's theorem : let $f \in L_{\infty}$, then $H_{f}$ has finite rank $\leq k$ if and only if $f \in H_{\infty, k}$.

The theorem of AAK gives an explicit solution to Pb . 1 .
Theorem (AAK): let $f \in L_{\infty}$ and $k \geq 0$. Then $\operatorname{dist}_{L_{\infty}}\left(\mathrm{f}, \mathrm{H}_{\infty, k}\right)=\sigma_{\mathrm{k}+1}\left(\mathrm{H}_{\mathrm{f}}\right)$, where the $\sigma_{\mathrm{i}}$ are the singular values of $\mathbf{H}_{f}$ ordered in decreasing value; furthermore this distance is attained at a unique function $\varphi_{k} \in H_{\infty, k}$ and if $v_{k+1}$ is the singular vector of $H_{f}$ corresponding to $\sigma_{\mathrm{k}+1}\left(\mathrm{H}_{\mathrm{f}}\right)$ then:

$$
\begin{equation*}
\varphi_{\mathbf{k}}(\mathrm{z})=\mathrm{f}(\mathrm{z})-\left[\mathrm{H}_{\mathbf{f}} \cdot \mathrm{v}_{\mathbf{k}+1}\right](\mathrm{z}) / \mathrm{v}_{\mathbf{k}+1}(\mathrm{z}) \tag{22.}
\end{equation*}
$$

where $\mathbf{H}_{\mathbf{f}} \cdot \mathbf{v}_{\mathbf{k}+1}=\boldsymbol{\Pi}_{\mathbf{-}}\left(\mathbf{v}_{\mathbf{k}+1} \cdot \mathrm{f}\right)$
Defining the projection operators $\Pi_{+}$and $\Pi_{\text {. }}$ b
$\Pi_{-}\left(\sum_{i=-\infty}^{+\infty} \alpha_{i} \cdot z^{-i}\right)=\sum_{i=1}^{+\infty} \alpha_{i} \cdot z^{-i} \quad, \Pi_{+}=I-\Pi_{-}$, oneobtains

$$
\varphi_{k}(z)=f(z)-\sigma_{k+1} \cdot \frac{\mu_{k+1}(z)}{v_{k+1}(z)}
$$

with:

$$
\mu_{k+1}(z)=\sum_{j=1}^{+\infty} \mathbf{u}_{j}^{k+1} \cdot z^{j}
$$

$$
v_{k+1}(z)=\sum_{j=1}^{+\infty} v_{j}^{k+1} \cdot z^{1-j}
$$

these results are interesting but not directly useful because they involve mainly the infinite Hankel matrix $H_{f}$ and its singular vectors. The implication of finite structure allows us to solve explicitely the initial Pb . 1 .
4.3. Implication of finite structure

Firstly, recall the classical equality: $\mathrm{H}=\mathrm{O} . \mathrm{C}$, therefore the square of the $i$-th singular value $\sigma_{i}$ of the (infinite) matrix $H$ is
$\sigma_{i}\left(H^{*}, H\right)=\sigma_{i}\left(C^{*}, 0^{*}, \mathbf{O} \cdot C\right)=\sigma_{i}\left(C \cdot C^{*} \cdot 0^{*}, 0\right)$

$$
=\sigma_{i}(P, Q)
$$

where $P$ and $Q$ are respectively the controllability and observability gramians. The matrices $P$ and $Q$ depend strongly on state-space coordinates but not the eigenvalues of the product $\mathbf{P}, \mathbf{Q}$, furthermore it has the great advantage to be finite-dimensional.

Let the strictly proper part of the transfert function be: $\quad f(z)=h^{*} .(z \cdot I-F)^{-1} \cdot \mathbf{T}=\frac{n(z)}{d(z)}$ where $d(z)=\operatorname{det}(z . I-F)$. By use of balanced system transformation, it is possible to obtain simpler expressions for thez functions $\mu_{i}(z)$ and $v_{i}(z)$.

More precisely, let $H=H(f)$ and $B$ a balancing
transform:
$\left(F, h^{*}, T\right) \xrightarrow{B}\left(F_{b}, h_{b}{ }^{*}, T_{b}\right)$
(balanced triple). Hence,
if $\mathrm{H}=\mathrm{U} \cdot \Sigma, \mathrm{V}$, then $\mathrm{O}_{\mathrm{b}}=\mathrm{U} \cdot \Sigma^{-1 / 2}$ and $\mathrm{C}_{\mathrm{b}}=\Sigma^{1 / 2} \cdot \mathrm{~V}$, therefore :

$$
\begin{equation*}
\mu_{i}(z)=\sigma_{i}^{-1 / 2} \cdot h_{b}^{*} \cdot\left(I_{n} \cdot z \cdot F_{b}\right)^{-1} \cdot e_{i} \tag{26}
\end{equation*}
$$

where $e_{i}$ denotes the $i$-th column of the ( $n \times n$ ) identity matrix.

Similarly, one obtains:
$v_{i}(z)=\sigma_{i}^{-1 / 2} \cdot e_{i}^{*} \cdot\left(z \cdot I_{n}-F_{b}\right)^{-1} \cdot T_{b}$
This last equality can be usefully transformed by means of the following property:

Prop. 5: for a balanced SISO system, the following equalities are satisfied: $\quad F_{b}{ }^{t}=\mathbf{Q} \cdot F_{b} \cdot Q^{*} ; T=Q^{*} \cdot h$, $Q$ being a unitary diagonal matrix.

From (26) we see that $\mu_{i}(z)$ has a rational form, i.e.:

$$
\mu_{i}(z)=\frac{m(z)}{d(z)}
$$

then (27) and Prop. 5 yield the following fundamental property:

$$
v_{i}(z)=q_{i} \cdot z^{-1} \cdot \mu_{i}\left(z^{-1}\right)=q_{i} \cdot \frac{\tilde{m}(z)}{d(z)} \quad\left(\text { with }\left|q_{i}\right|^{2}=1\right)
$$

where $q_{i}$ is the $i$-th element of the diagonal $Q$ and $\tilde{m}(z)$, ( $d(z)$ ) the reversed polynomials of $m(z)$ (resp. $d(z)$ ). i.e. $\quad \tilde{m}(z)=z^{n-1} \cdot \bar{m}\left(z^{-1}\right)$

Using AAK theorem the best $\mathbf{L}_{\infty}$ approximation of order $\mathbf{k}$ is given by:

$$
\begin{align*}
\varphi_{k}(z) & =\frac{n(z)}{d(z)} \cdot q_{k+1} \cdot \sigma_{k+1} \cdot \frac{m(z) \cdot \tilde{d}(z)}{\tilde{m}(z) \cdot d(z)} \\
& =\frac{q(z)}{d(z) \cdot \tilde{m}(z)} \tag{28}
\end{align*}
$$

Furthermore it can be easily proved that the rational fraction $q(z) / d(z)$ is actually a polynomial $p(z)$ of degree $<\mathbf{n}$, as a consequence the following equality is valid:
with :

$$
\begin{gathered}
\varphi_{k}(z)=\frac{p(z)}{\tilde{m}(z)} \\
p(z)=\frac{n(z) \cdot \tilde{m}(z)-\lambda \cdot \tilde{d}(z) \cdot m(z)}{d(z)} \\
\left(\lambda=q_{k+1} \cdot \sigma_{k+1}\right)
\end{gathered}
$$

This last form of $\varphi_{k}(z)$ leads to determine the coefficients of $m(z)$ by solving the following
polynomial equation:

$$
\begin{equation*}
p(z) \cdot d(z)=n(z) \cdot \tilde{m}(z)-\lambda \cdot d(z) \cdot m(z) \tag{29}
\end{equation*}
$$ equating the coefficients of powers of $z$ [7], one obtains the two following matricial equations (30) and (31);

$J . S_{2} \cdot J \cdot p=R_{2} \cdot \dot{m}-\lambda \cdot J \cdot \dot{S}_{1} \cdot \bar{J}, m \quad\left(z^{i}: i=0, \ldots, n-1\right)$

$$
S_{1} \cdot p=R_{1} \cdot \bar{m}-\lambda \cdot \dot{S}_{2} \cdot m \quad\left(z^{i}: i=n, \ldots, 2 n-1\right)
$$

where :
$n(z)=c_{n}+c_{n-1} \cdot z+\ldots+c_{1} \cdot z^{n-1}$
$d(z)=a_{n}+a_{n-1} \cdot z+\ldots+a_{1} \cdot z^{n-1}+z^{n}$
$m(z)=m_{n-1}+m_{n-2} \cdot z+\ldots+m_{0} \cdot z^{n-1}$
$p(z)=p_{n-1}+p_{n-2} \cdot z+\ldots+p_{0} \cdot z^{n-1}$
and:

$\mathbf{R}_{\mathbf{1}} \equiv\left[\begin{array}{llll}0 & c_{1} & \ldots & c_{\mathbf{n - 1}} \\ \cdots & \ddots & \ddots & \vdots \\ \cdots & \ddots & \cdots \\ \cdots & & c_{1} \\ 0 & \cdots & \cdots & 0\end{array}\right] \mathbf{R}_{\mathbf{2}}=\left[\begin{array}{llll}c_{\mathbf{n}} & & \cdots & 0 \\ \vdots & \ddots & & \\ \vdots & & & \\ c_{2} & & \ddots & \vdots \\ c_{1} & c_{2} & \cdots & c_{n}\end{array}\right]$
and $J$ the antidiagonal matrix
from (30) and (31) one obtains (33):
$\left(\mathrm{J} . \mathrm{S}_{2} \cdot \mathrm{~J} \cdot \mathrm{~S}_{1}{ }^{-1} \cdot \mathrm{R}_{1}-R_{2}\right) \cdot \ddot{m}=\lambda .\left(\mathrm{J} . \mathrm{S}_{2} \cdot \mathrm{~J} . \mathrm{S}_{1}{ }^{-1} \cdot \bar{S}_{2}-\mathrm{J} \cdot \bar{S}_{1} \cdot \mathrm{~J}\right) \cdot \mathrm{m}$ $m^{t}=\left(m_{0}, m_{1}, \ldots, m_{n-1}\right)$
which appears as a generalized eigenvalueeigenvector equation. In order to solve (33), numerous subroutines are available. Once the coefficients of $m(z)$ have been calculated, those of $p(z)$ are straightforwardly deduced.

$$
\begin{aligned}
\mathbf{J} \cdot \mathbf{S}_{2} \cdot \mathbf{J} \cdot \mathbf{p} & =\mathbf{R}_{2} \cdot \bar{m}-\lambda \cdot J \cdot \bar{S}_{1} \cdot J \cdot m \quad\left(z^{i}: i=0, \ldots, n-1\right) \\
S_{1} \cdot p & =R_{1} \cdot \bar{m}-\lambda \cdot \bar{S}_{2} \cdot m \quad\left(z^{i}: i=n, \ldots, 2 n-1\right)
\end{aligned}
$$

4.4. Practical implementation:

Practical utilization of the previous method needs a prior state-space model which will be reduced. The more "random" model (or maximum entropy) 'seems to be the most convenient, it summarizes all the available informations (i.c. : the estimated covariance of the outputs) and gives the more "random" extension of them. For a linear equispaced array of sensors it is similar with an AR modelling. The practical algorithm takes the following form after DFT of sensors:

1) estimation of the covariances of the sensors output $(\underline{n}+1$ sensors); i.e. : $\hat{r}(0), \hat{r}(1), \ldots, \hat{r}(n)$.
2) estimation of the corresponding $A R$ model coefficients: $\quad\left(1, \hat{a}_{1}, \ldots, \hat{a}_{n} ; b^{2}\right)$.
3) solve the generalized eigenvalue-eigenvector problem (33) with:

$$
\begin{aligned}
& a_{1}=\hat{a}_{1}, \ldots, a_{n-1}=\hat{a}_{n-1}, a_{n}=\hat{a}_{n} \\
& c_{1}=0, \ldots, c_{n-1}=0, c_{n}=1
\end{aligned}
$$

select the $(k+1)$-th generalized eigenvector (to obtain a rank $k$ approximation) and calculate $m_{k+1}(z)$.
4) compute the roots of $m_{k+1}(z)$ and their arguments (note that by definition of $H_{\infty, k}$ there are $k$ roots in $\mathbf{D}(0,1)$ ).

## 5. Simulation results.

Due to the lack of space, we present briefly simulation results, details are available in [6].
5.1. Approximated realization methods (ARM):

The case of an additive white noise is lirslly considered. Three uncorrelated sources impinge a linear array constituted by 16 equispaced sensors (at $\lambda / 2$ ). The covariance matrix of the outputs is extimated by averaged periodogram (BT $=100$ ). The source bearings are $56,60,65$ degrees with respective levels $-1,-11,-1 \mathrm{~dB}$. The assumed sources number is 4 , the results are summarized below:

| Detection of a veak source | MUSIC | ARM |
| :---: | :--- | :--- |
| 100 trials | $15 / 100$ | $40 / 100$ |

5.3. Optimal realization method (ORM):

Simulation 1: the additive noise is spatially white, 3 sources are simulated with respective powers and bearings ( $30 \mathrm{deg}, 0 \mathrm{~dB}$ ), ( $36 \mathrm{deg},-10 \mathrm{~dB}$ ), $(45 \mathrm{deg}, 7 \mathrm{~dB}$ ); the number of averaged snapshots (BT) is equal to 300 . Afterwhat the poles of Music method and ORM are computed, the results are presented below (for ten trials)

| Detection | MUSIC | ORM |
| :---: | :---: | :---: |
| of the reak source | $0 / 10$ | $8 / 10$ |

Note that for MUSIC method the pole associated to the weak source ( -10 dB ) can not be seen generally or at best is very small ( $\cong 0.5$ ), conversely the pole associated to that source is very close to the strongest ones for ORM (typically 0.98 !).
Simulation 2: in order to judge ability of ORM to separate 2 close sources, two sources of same powers (0) dB ) have been simulated at bearings 10 and 13 deg (BT $=$ 300 ), with spatially white noise; results are summarized below:

| Separation | MUSIC | ORM |
| :---: | :---: | :---: |
| of the two sources | $4 / 10$ | $9 / 10$ |

Simulation 3: we consider now the simulation of $\S 5.2$, the additive noise is strongly correlated. For ten trials, the sources are correctly seen at their exact bearings for ORM.

## 6. Conclusion.

The simulation results have demonstrate the effectiveness of a state-space modelling of the sensors output. Robustness versus noise correlations is greatly enhanced by use of ARM (w.r. "usual " high resolution method), it is their main advantage. Detection of weak source can be slightly enhanced also but estimation of sources subspace remains based upon eigenvector decomposition

Conversely, ORM takes advantage of the estimation of sources subspace by use of their special structure and present considerable improvements for angular separation of close sources, detection of weak sources. Furthermore, the method is very robust w.r. noise correlations, its computation cost is similar to MUSIC method.

Obviously the presented tools are powerful for a linear array; however they can not be directly (if it is possible) extended to a general shape.

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