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The rank of the covariance matrix of an evanescent field

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ABSTRACT

Evanescent random fields arise as a component of the 2D Wold decomposition of homogeneous random fields. Besides their theoretical importance, evanescent random fields have a number of practical applications, such as in modeling the observed signal in the space–time adaptive processing (STAP) of airborne radar data. In this paper we derive an expression for the rank of the low-rank covariance matrix of a finite dimension sample from an evanescent random field. It is shown that the rank of this covariance matrix is completely determined by the evanescent field spectral support parameters, alone. Thus, the problem of estimating the rank lends itself to a solution that avoids the need to estimate the rank from the sample covariance matrix. We show that this result can be immediately applied to considerably simplify the estimation of the rank of the interference covariance matrix in the STAP problem.

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1. Introduction

1.1. The evanescent random field

The problem of linear prediction of stationary processes is a classic problem in time-series analysis. One of the most fundamental results in this field is the Wold decomposition [1], that states that a regular one-dimensional wide-sense stationary processes indexed by \mathbb{Z} may be decomposed into two stationary and orthogonal components: the purely-indeterministic process (that produces the innovations) and the deterministic process. This decomposition can be equivalently reformulated using spectral notations: the spectral measure of the purely-indeterministic process is absolutely continuous with respect to the Lebesgue measure, and the spectral measure of the deterministic process is singular. In other words, the spectral measures of the orthogonal components of Wold decomposition yield the Lebesgue decomposition of the spectral measure of the process.

Homogeneous random fields, (also called doubly stationary series), are the two-dimensional (indexed by \mathbb{Z}^2) generalization of the one-dimensional wide-sense stationary process. Unfortunately, unlike the one-dimensional case, in multiple dimensions there is no natural order definition and terms such as “past” and “future” are meaningless unless defined with respect to a specific order. Linear prediction of homogeneous random fields was first rigorously formulated by Helson and Lowdenslager in [2]. The problem of defining “past” and “future” on the 2D lattice (i.e., \mathbb{Z}^2) was resolved in [2] in terms of “half-plane” total-ordering. The trivial example of a half-plane total order on \mathbb{Z}^2 is a usual lexicographic order: $(k, l) \preceq (n, m)$ iff $k < n$ or $(k = n$ and $l < m)$. Lexicographic order can be considered as a linear order induced by Non-Symmetric (delimited by a broken straight line) Half-Plane (NSHP), (see Fig. 1).

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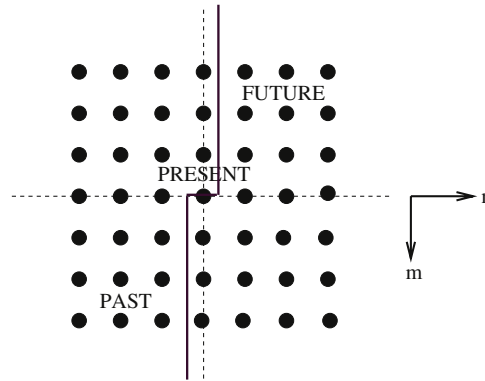


Fig. 1. Non-symmetric half-plane.

Further analysis of the prediction problem led to a generalization of the Wold decomposition [3]. When we consider random processes indexed by a group we obtain a Wold decomposition with respect to any given total order on the group. When the group is not \mathbb{Z} (like \mathbb{R} or \mathbb{Z}^2) the deterministic process can have as a direct summand a deterministic process of a special type, the *evanescent process*. In order to provide some intuition on the characteristics of the evanescent process we next state some basic definitions and present an example of an evanescent field defined with respect to a vertical total order, which is simply a lexicographic order on \mathbb{Z}^2 :

A homogeneous random field $\{y(n, m)\}$ is called *regular* with respect to the *lexicographic order* if for every (n, m) , $E[y(n, m) - \hat{y}(n, m)]^2 = \sigma^2 > 0$ where $\hat{y}(n, m)$ is the projection of $y(n, m)$ on the *c.l.m.* $[\{y(k, l) : k < n, l \in \mathbb{Z}\} \cup \{y(n, l) : l < m\}]$, where *c.l.m.* denote a closed linear manifold. Thus, a regular homogeneous random field has a non-zero innovation at every lattice point. A homogeneous random field $\{z(n, m)\}$ is called *deterministic* with respect to the lexicographic order if it can be perfectly linearly predicted from its past in mean-square sense, i.e., for every (n, m) we have $z(n, m) \in \text{c.l.m.} [\{z(k, l) : k < n, l \in \mathbb{Z}\} \cup \{z(n, l) : l < m\}]$.

Although a deterministic field can be perfectly predicted from its past with respect to lexicographic order, it may still possess a non-zero innovation when prediction is based on samples in previous columns only. We then say that the field $\{z(n, m)\}$ has vertical *column-to-column innovations* if $I(n, m) := z(n, m) - \hat{z}(n, m)$ (the *innovation*) is not 0, where $\hat{z}(n, m)$ is the orthogonal projection of $z(n, m)$ on the closed subspace generated by $\{z(k, l) : k < n, l \in \mathbb{Z}\}$. In other words, if a deterministic field has non-zero column-to-column innovations it cannot be perfectly linearly predicted from previous columns.

When $z(n, m)$ is the deterministic component of the decomposition of a regular random field with respect to a NSHP total-ordering, the vertical evanescent component $z_e(n, m)$ is the orthogonal projection of $z(n, m)$ on the closed subspace generated by the (orthogonal) column-to-column innovations $\{I(k, l) : k \leq n, l \in \mathbb{Z}\}$. Thus, an evanescent field spans a Hilbert space identical to the one spanned by column-to-column innovations. In other words, the evanescent field is a component of the deterministic field which represents column-to-column innovations. Horizontal column-to-column (row-to-row) innovations and evanescent components are similarly defined.

Evanescent processes were first introduced in [3] (on \mathbb{R}). In Korezlioglu and Loubaton [4], “horizontal” and “vertical” total orders and the corresponding horizontally and vertically evanescent components of a homogeneous random field on \mathbb{Z}^2 are defined. In Kallianpur [5], as well as in Chiang [6], similar techniques are employed to obtain four-fold orthogonal decompositions of regular (non-deterministic) homogeneous random fields. In Francos et al. [7] this decomposition of random fields on \mathbb{Z}^2 was further extended. This is done by considering *all* the Rational Non-Symmetrical Half-Plane (RNSHP) linear orders, each inducing a different partitioning of the 2D lattice into two sets by a broken straight line of rational slope. Intuitively, the usual lexicographic order is not the only possible order definition of the 2D lattice. Each RNSHP linear order is induced by a “rotation” of the usual lexicographic order, such that the resulting non-symmetrical half-plane is delimited by a broken straight line with rational slope, and which leads to a different linear order definition. Consequently, terms such as “past” and “future” are redefined with respect to a specific RNSHP linear order (see, for example, Fig. 2).

More specifically, each Rational Non-Symmetrical Half-Plane is defined in terms of two co-prime integers (a, b) , such that the past $P_{a,b}$ is defined by

$$P_{a,b} = \{(n, m) \in \mathbb{Z}^2 : na + mb < 0, \text{ or } na + mb = 0 \text{ and } m \leq 0\}. \tag{1}$$

Then $P = P_{a,b}$ satisfies

$$(i) P \cap (-P) = \{0\}, \quad (ii) P \cup (-P) = \mathbb{Z}^2, \quad (iii) P + P \subset P \text{ (usual addition).}$$

By (i)–(iii), P induces on \mathbb{Z}^2 a linear order, which is defined by $(k, l) \preceq (n, m)$ if and only if $(k - n, l - m) \in P$.

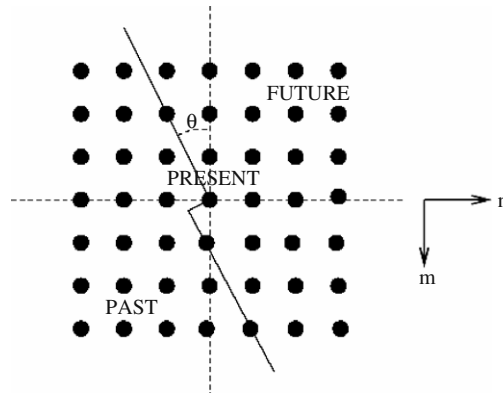


Fig. 2. RNSHP support.

Clearly, there are countably many such linear orders. Each such order induces a different definition of the term “column”, and correspondingly different definitions of column-to-column innovations and evanescent field.

The Wold decomposition of a regular random field into purely-indeterministic and deterministic components is invariant to the choice of a RNSHP order. The decomposition in [7] further asserts that we can represent the deterministic component of the field as a mutually orthogonal sum of a “half-plane deterministic” field and a countable number of evanescent fields. The half-plane deterministic field has no innovations, nor column-to-column innovations, with respect to any RNSHP linear order. On the other hand, each of the evanescent fields can be revealed only by using the corresponding RNSHP linear order, i.e., with respect to specific definitions of “columns” and column-to-column innovations. This decomposition yields a corresponding spectral decomposition, i.e., we can decompose the spectral measure of the deterministic part into a countable sum of mutually singular spectral measures, such that the spectral measure of each evanescent component is concentrated on a line with a rational slope.

Based on these results, a parametric model of the homogeneous random field was derived in [7]. The purely-indeterministic component of the field is modeled by a white innovations driven 2D moving average process with respect to some RNSHP linear order. This component contributes the absolutely continuous part of the spectral measure of the regular field. One of the components of the half-plane deterministic component that is often found in practical applications is the 2D harmonic random field which is the sum of a countable number of exponential components, each having a constant spatial frequency and random amplitude. This component contributes the 2D delta functions in the spectral domain. The number of evanescent components of the regular field is countable. The model of the evanescent field with respect to specific order is presented below:

Let (a, b) be a pair of co-prime integers ($a \geq 0$) which defines a specific RNSHP linear order according to (1). Then, the model of the evanescent field which corresponds to this order is

$$e_{(a,b)}(n, m) = \sum_{i=1}^{I_{(a,b)}} s_i^{(a,b)}(na + mb) \exp\left(j\omega_i^{(a,b)}(nc + md)\right), \tag{2}$$

where c and d are co-prime integers satisfying $ad - bc = 1$. For the case where $(a, b) = (0, 1)$ we have $(c, d) = (1, 0)$, and for $(a, b) = (1, 0)$ we have $(c, d) = (0, 1)$. We further note that in this notation $na + mb$ is the “column” index and $nc + md$ defines a “row”. The modulating process $\{s_i^{(a,b)}(na + mb)\}$ is a 1D purely-indeterministic, complex-valued processes, and $\omega_i^{(a,b)}$ is a modulation frequency. Thus, $e_{(a,b)}$ has no innovations, with respect to “rows”, and has non-zero column-to-column innovation (expressed by the modulating process $s_i^{(a,b)}$) with respect to its “columns”. $I_{(a,b)}$ denotes the number of different evanescent components that correspond to the same RNSHP defined by (a, b) . The different components are such that their 1D modulating processes $\{s_i^{(a,b)}\}$ and $\{s_j^{(a,b)}\}$, are *mutually orthogonal* and their modulation frequencies are different $\omega_i^{(a,b)} \neq \omega_j^{(a,b)}$ for all $1 \leq i \neq j \leq I_{(a,b)}$.

The “spectral density function” of each evanescent field has the form of a sum of 1D delta functions which are supported on lines of rational slope in the 2D spectral domain. The amplitude of each of these delta functions is determined by the spectral density of the 1D modulating process. Since the spectral density of the modulating process can rapidly decay to zero, so will the “spectral density” of the evanescent field, and hence the name “evanescent”.

1.2. Practical applications

Besides its fundamental theoretical importance, the Wold decomposition of a regular random field has various applications in image processing and wave propagation problems. For example, the parametric model that results from

these orthogonal decompositions, naturally arises as the physical model in problems of texture modeling, estimation and synthesis [8].

Another application is space–time adaptive processing of airborne radar data [9]. Space–time adaptive processing (STAP) is an increasingly popular radar signal processing technique for detecting slow-moving targets. The space dimension arises from the use of array of multiple antenna elements and the time dimension arises from the use of coherent train of radar pulses. The power of STAP comes from the joint processing along the space and time dimensions. Comprehensive analysis of STAP appears in [10,11].

In [9] it is shown that the same parametric model that results from the 2D Wold orthogonal decomposition naturally arises as the physical model in the problem of space–time processing of airborne radar data. This correspondence is exploited to derive computationally efficient detection algorithms. More specifically, the target signal is modeled as a random amplitude complex exponential where the exponential is defined by a space–time steering vector that has the target’s angle and Doppler. Thus, in the space–time domain the target contribution is the half–plan deterministic component of the observed field. The sum of the white noise field due to the internally generated receiver amplifier noise, and the sky noise contribution, is the purely–indeterministic component of the space–time field decomposition.

The presence of a jammer (a foe interference source, transmitting high power noise aimed at “blinding” the radar system) results in a barrage of noise localized in angle and uniformly distributed over all Doppler frequencies (since the transmitted noise is white). Hence, in the space–time domain each jammer is modeled as an evanescent component with $(a, b) = (0, 1)$ such that its 1D modulating process $s_i^{(0,1)}(m)$ is the random process of the jammer amplitudes. The jammer samples from different pulses are uncorrelated. In the angle–Doppler domain each jammer contributes a 1D delta function, parallel to the Doppler axis and located at a specific angle $\omega_i^{(0,1)}$ using the notation of (2).

Since the ground clutter is received from all angles and is spread in Doppler frequency due to platform motion it results in an additional evanescent component of the observed 2D space–time field. The aircraft platform motion produces a very special structure of the clutter due to the dependence of the Doppler frequency on angle. The clutter’s echo from a single ground patch has a Doppler frequency that linearly depends on its aspect with respect to the platform. As the platform moves, identical clutter observations are repeated by different antenna elements on different pulses, which defines a specific linear locus in the angle–Doppler domain, commonly referred as the “clutter ridge”. Thus, the clutter ridge, which also represents clutter from all angles, is supported on a diagonal line (that generally wraps around) in the angle–Doppler domain. The clutter received by the stationary platform can be modeled as an evanescent component with $(a, b) = (1, 0)$ which results in 1D delta function parallel to angle domain and with zero Doppler frequency. Due to the physical properties of the problem the different components of the field are assumed to be mutually orthogonal. In the specific application of airborne radar, the evanescent components (the clutter, and jamming signals) are considered unknown interferences.

Although the data collected by STAP radars for different ranges can be viewed as a sequence of finite sample realizations from a homogeneous field, it is technically more convenient to represent each of the observations in a vector form and to statistically analyze them as multivariate vectors. Thus, if one uses a STAP system with N antenna elements and M pulses, the observed $N \times M$ STAP signal is treated as $NM \times 1$ multivariate random variable. These vectors are commonly called “snapshots”.

The STAP processor goal is to solve a detection problem, *i.e.*, to establish whether a hypothetical target is present or not. It adaptively weights the available data in order to achieve high gain at the target’s angle and Doppler and maximal mitigation along both the jamming and clutter lines. The adaptive weight vector is computed from the inverse of the interference–plus–noise covariance matrix [10,11]. It is shown in [12] that the dominant eigenvectors of the space–time covariance matrix contain all the information required to mitigate the interference. Thus, the weight vector is constrained to be in the subspace orthogonal to the dominant eigenvectors. Because the interference–plus–noise covariance matrix is unknown a priori, it is typically estimated using sample covariances obtained from averaging over a few range gates. This is the known as the *fully adaptive* STAP approach. The major drawback of this approach is its high computational complexity. The final detection of a target is performed by applying either Constant-False-Alarm-Rate (CFAR) detector, or Adaptive-Matched-Filter (AMF) detector, or Generalized-Likelihood-Ratio (GLR) detector. Usually the detector is embedded into the weight computations.

Fortunately, both the clutter and the jammers have low-rank covariance matrices. The clutter covariance matrix has a low rank due to the movement of the platform and presence of the clutter at all angles, as discussed above. The jammer covariance matrix has low rank since the jamming signal is spatially correlated between all antennas at each pulse. The low-rank structure of the interference covariance matrix may be exploited to achieve significant reduction in the adaptive problem dimensionality with little or no sacrifice in performance relative to the fully adaptive case. These methods are referred to as *partially adaptive* STAP.

Partially adaptive STAP methods require knowledge of the rank of the interference covariance matrix. However, it is a priori unknown, and unfortunately cannot be easily estimated from the sample covariance matrix due to the existence of a noise component which has a full rank covariance matrix. Hence, the problem of estimating the rank of the interference covariance matrix is critical in the implementation of many STAP algorithms.

In this work we consider the problem of determining the rank of the covariance matrix of a vectorized finite dimension sample from an evanescent random field. By using the evanescent field parametric model to model the interferences in the STAP problem, we considerably simplify the solution to the problem of estimating the rank of the low-rank interference covariance matrix. In fact, it turns out that in this parametric framework the well-known Brennan rule [11] for the rank of

the clutter covariance matrix, as well as the rank computation of the jammer, become special cases of the general result proved here. Hence, the provided derivation opens the way for new, computationally attractive, methods in parametric and non-parametric estimation of 2D random fields, with immediate applications in partially adaptive space–time adaptive processing of airborne radar data.

The rest of the paper is organized as follows: In Section 2 we formulate the problem we aim to solve. A formula for the rank of the covariance matrix of a complex-valued evanescent random field is derived in Section 3. In Section 4 we extend the obtained result to the case of a real-valued evanescent random field. Finally, in Section 5 we provide our conclusions.

The following notation is used throughout. Boldface upper case letters denote matrices, boldface lower case letters denote column vectors, and standard lower case letters denote scalars. The superscripts $(\cdot)^T$ and $(\cdot)^H$ denote the transpose and Hermitian transpose operators, respectively. By \mathbf{I} we denote the identity matrix and by $\mathbf{0}$ a matrix of zeros. The symbol \odot denotes an element by element product of the vectors (Hadamard product). Given a scalar function $f(\cdot)$ and a column vector \mathbf{v} , we denote by $f(\mathbf{v})$ a column vector consisting of the values of function $f(\cdot)$ evaluated for all the elements of \mathbf{v} . Finally $\text{diag}(\mathbf{v})$ denotes a square diagonal matrix with the elements of \mathbf{v} on its main diagonal.

2. Finite sample of an evanescent random field: Definitions and problem formulation

Let O denote the set of all possible pairs of different co-prime integers (a, b) , $a \geq 0$, where each pair defines a RNSHP order on the 2D lattice. Although in the case of an infinite 2D lattice the number of different RNSHP definitions is infinitely countable, in the finite sample case only a finite number of different linear orders can be defined. Moreover, in practical applications the number of different evanescent components for each order definition is finite as well. Therefore, we assume throughout this paper that $|O|$ and $I_{(a,b)}$ are finite integers.

Let $Q = \sum_{(a,b) \in O} I_{(a,b)}$ is a total number of evanescent components. Each of the Q evanescent components is uniquely defined by the triple $(a, b, \omega_i^{(a,b)})$ where $(a, b) \in O$ and $\omega_i^{(a,b)}$ is the frequency parameter. Let us denote the set of all possible triples by $O_Q = \{(a_1, b_1, \omega_1), (a_2, b_2, \omega_2), \dots, (a_Q, b_Q, \omega_Q)\}$. All triples are unique: they either have different support parameters $(a_i, b_i) \neq (a_j, b_j)$, or in case $(a_i, b_i) = (a_j, b_j)$ they have different frequencies such that $\omega_i \neq \omega_j$.

Finally, adapting (2) to the finite sample case we have that

$$e(n, m) = \sum_{q=1}^Q e_q(n, m), \tag{3}$$

where

$$e_q(n, m) = s_q(na_q + mb_q) \exp\left(j\omega_q(nc_q + md_q)\right) \tag{4}$$

such that $(a_q, b_q, \omega_q) \in O_Q$ and $\{s_q\}, c_q, d_q$ are defined as above.

We note that since the spectral measure of $\{e_q(n, m)\}$ is concentrated on a line (that may wrap around) whose slope is determined by a_q and b_q , we interchangeably refer to a_q and b_q as either the *spectral support parameters* of $\{e_q\}$ or as the *RNSHP slope parameters*.

Let $\{e(n, m) : (n, m) \in D\}$ where $D = \{(n, m) \in \mathbb{Z}^2 : 0 \leq n \leq N - 1, 0 \leq m \leq M - 1\}$ be the observed finite sample of the random field (3). Let \mathbf{e} denote an $NM \times 1$ vector-form representation of this finite sample:

$$\mathbf{e} = [e(0, 0), \dots, e(0, M - 1), e(1, 0), \dots, e(1, M - 1), \dots, e(N - 1, 0), \dots, e(N - 1, M - 1)]^T. \tag{5}$$

This is a multivariate representation of a finite sample of an evanescent random field.

Let Γ denote the $NM \times NM$ covariance matrix of the evanescent vector \mathbf{e} ,

$$\Gamma = E\left[\mathbf{e}(\mathbf{e}^H)\right]. \tag{6}$$

Due to the special structure of the evanescent field, many of the elements of \mathbf{e} are linearly dependent, and therefore Γ is low-rank. This property is easily demonstrated by considering a single evanescent component that corresponds to the vertical order $(a, b) = (0, 1)$ (single jammer source using the STAP nomenclature), with some arbitrary modulation frequency ω and modulating process $s(m)$. In that case

$$e_{(0,1)}(n, m) = s(m) \exp\left(j\omega n\right). \tag{7}$$

The vector-form representation of the finite sample of this evanescent field is

$$\mathbf{e}_{(0,1)} = [s(0), s(1), \dots, s(M - 1), s(0) \exp(j\omega), \dots, s(M - 1) \exp(j\omega), \dots, s(0) \exp(j\omega(N - 1)) \dots, s(M - 1) \exp(j\omega(N - 1))]^T. \tag{8}$$

Since the modulating frequency ω is a deterministic constant, it is obvious that $\mathbf{e}_{(0,1)}$ is comprised of only M independent random variables. Therefore, the rank of $\mathbf{\Gamma}_{(0,1)} = E[\mathbf{e}_{(0,1)}(\mathbf{e}_{(0,1)})^H]$ is also M .

The aim of this paper is to derive an expression for the rank of the low-rank covariance matrix $\mathbf{\Gamma}$ of the evanescent vector \mathbf{e} , in the general case (3).

3. The rank of the covariance matrix of an evanescent field

In this section we derive an expression for the rank of the covariance matrix $\mathbf{\Gamma}$. In order to do so, we have to find and quantify the linear dependencies between the samples of \mathbf{e} . Unfortunately, for arbitrary spectral support parameters and multiple evanescent components, this task involves tedious calculations. The results of the entire analysis in this section can be summarized by the following theorem:

Theorem 1. Let \mathbf{e} be a vector-form representation of a finite sample from a sum of evanescent random fields, given by (3)–(5). Then, the rank of its covariance matrix, $\mathbf{\Gamma}$, is given by

$$\text{rank}(\mathbf{\Gamma}) = \min \left(NM, \left[N \sum_{q=1}^Q |a_q| + M \sum_{q=1}^Q |b_q| - \sum_{q=1}^Q |a_q| \sum_{q=1}^Q |b_q| \right] \right). \tag{9}$$

Even though the evaluations in the next subsections are technical in nature, the obtained result is surprisingly interesting. Hence, before addressing the proof itself let us make some comments. From Theorem 1, it is clear that the rank of the covariance matrix of a finite sample from an evanescent random field is completely determined by the spectral support parameters (a_q, b_q) of the different evanescent components, while it is independent of the other parameters of the evanescent fields, such as the parameters of the modulating processes, $\{s_q\}$, or the modulation frequencies, ω_q .

Moreover, one can easily observe that the well-known Brennan rule for the rank of the low-rank clutter covariance matrix in the STAP framework, [11] as well as the rank of the covariance matrix of the jamming signals are special cases of this theorem. The Brennan rule states that the rank of the clutter covariance $\mathbf{\Gamma}_{clut}$ is given by:

$$\text{rank}(\mathbf{\Gamma}_{clut}) = \lfloor N + M\beta - \beta \rfloor \tag{10}$$

where β is the slope of the clutter ridge orientation in the angle–Doppler domain, and $\lfloor \cdot \rfloor$ denotes rounding to the nearest integer. It is easy to see that this formula is a special case of the above theorem when only a single evanescent field is observed, and its spectral support parameters are $(a, b) = (1, \beta)$. The rank of the jamming covariance matrix is, [11]:

$$\text{rank}(\mathbf{\Gamma}_{jam}) = MJ \tag{11}$$

where J is a number of sources. Since the spectral support of a single jammer in the angle–Doppler domain is a line parallel to the Doppler axis, and since all jammers are mutually orthogonal, they can be modeled as J vertical evanescent components with spectral support parameters $(a, b) = (0, 1)$, such that the rank of the covariance matrix of each individual jammer is M as in the above example.

3.1. Rank derivations

In this subsection we prove Theorem 1. The derivation provides an insight into the structure of the covariance matrix, and explains the nature of its low-rank. Moreover, we explicitly show how columns of the covariance matrix, that can be represented as linear combinations of other columns, are formed, which yields its low-rank.

Rewriting (3) in a vector form we have $\mathbf{e} = \sum_{q=1}^Q \mathbf{e}_q$, where

$$\mathbf{e}_q = [e_q(0, 0), \dots, e_q(0, M - 1), e_q(1, 0), \dots, e_q(1, M - 1), \dots, e_q(N - 1, 0), \dots, e_q(N - 1, M - 1)]^T. \tag{12}$$

Let

$$\begin{aligned} \boldsymbol{\xi}_q = [& s_q(0), s_q(b_q), \dots, s_q((M - 1)b_q), s_q(a_q), s_q(a_q + b_q), \dots, s_q(a_q + (M - 1)b_q), \\ & \dots, s_q((N - 1)a_q), s_q((N - 1)a_q + b_q), \dots, s_q((N - 1)a_q + (M - 1)b_q)]^T \end{aligned} \tag{13}$$

be the vector whose elements are the observed samples from the 1D modulating process $\{s_q\}$. Define

$$\begin{aligned} \mathbf{v}_q = [& 0, d_q, \dots, (M - 1)d_q, c_q, c_q + d_q, \dots, c_q + (M - 1)d_q, \dots, \\ & (N - 1)c_q, (N - 1)c_q + d_q, \dots, (N - 1)c_q + (M - 1)d_q]^T. \end{aligned} \tag{14}$$

Let

$$\mathbf{D}_q = \text{diag} \left(\exp(-j\omega_q \mathbf{v}_q) \right). \tag{15}$$

be an $NM \times NM$ diagonal matrix. Thus, using (4), we have that

$$\mathbf{e}_q = \mathbf{D}_q^H \boldsymbol{\xi}_q. \tag{16}$$

Let \mathbf{s}_q be a $(N - 1)|a_q| + (M - 1)|b_q| + 1$ dimensional column vector of consecutive samples of the 1D modulating process $\{s_q\}$. For the case in which $a_q > 0$ and $b_q < 0$, \mathbf{s}_q is defined as

$$\mathbf{s}_q = [s_q((M - 1)b_q), \dots, s_q((N - 1)a_q)]^T, \tag{17}$$

while for the case in which $a_q \geq 0$ and $b_q \geq 0$, \mathbf{s}_q is defined as

$$\mathbf{s}_q = [s_q(0), \dots, s_q((N - 1)a_q + (M - 1)b_q)]^T. \tag{18}$$

Thus for any (a_q, b_q) we have that

$$\boldsymbol{\xi}_q = \mathbf{A}_q^T \mathbf{s}_q \tag{19}$$

and

$$\mathbf{e}_q = \mathbf{D}_q^H \mathbf{A}_q^T \mathbf{s}_q, \tag{20}$$

where \mathbf{A}_q is a real-valued $[(N - 1)|a_q| + (M - 1)|b_q| + 1] \times NM$ rectangular matrix where each of its columns has a single element whose value is “1”, while all the others are zero. Thus, each column of \mathbf{A}_q “chooses” the single element from the vector \mathbf{s}_q that contributes to the corresponding element of the vector $\boldsymbol{\xi}_q$. Due to boundary effects, resulting from the finiteness of the observation, not *all* of the elements of the vector \mathbf{s}_q contribute to the vector $\boldsymbol{\xi}_q$, unless $|a_q| \leq 1$ or $|b_q| \leq 1$. Hence some rows of the matrix \mathbf{A}_q may contain only zeros. On the other hand, whenever $na_q + mb_q = ka_q + \ell b_q$ for some integers n, m, k, ℓ such that $0 \leq n, k \leq N - 1$ and $0 \leq m, \ell \leq M - 1$, the same sample from the modulating process $\{s_q\}$ is duplicated in the elements of $\boldsymbol{\xi}_q$. Therefore, the number of *distinct* columns in \mathbf{A}_q is equal to the number of elements of \mathbf{s}_q that appear in $\boldsymbol{\xi}_q$, i.e., the number of *distinct* samples from the random process $\{s_q\}$ that are found in an observed finite sample of an evanescent field of dimensions $N \times M$. The matrix \mathbf{A}_q depends only on (a_q, b_q) and is independent of the modulation frequency ω_q or the modulating process $\{s_q\}$.

The rank of covariance matrix $\boldsymbol{\Gamma}$ is strongly related to the number of *distinct* samples from the random processes $\{s_q\}$ for all $1 \leq q \leq Q$ which can be found in the evanescent vector \mathbf{e} . Therefore, the rank of $\boldsymbol{\Gamma}$ is tightly related to the ranks of the matrices \mathbf{A}_q , $1 \leq q \leq Q$.

Let \mathbf{R}_q be the covariance matrix of the vector \mathbf{s}_q i.e.,

$$\mathbf{R}_q = E \left[\mathbf{s}_q (\mathbf{s}_q)^H \right]. \tag{21}$$

The matrix \mathbf{R}_q is full rank positive definite since the process $\{s_q\}$ is purely-indeterministic. Since the evanescent components $\{e_q\}$ are mutually orthogonal we conclude that $\boldsymbol{\Gamma}$, the covariance matrix of \mathbf{e} , has the form

$$\boldsymbol{\Gamma} = E \left[\mathbf{e} (\mathbf{e})^H \right] = \sum_{q=1}^Q \boldsymbol{\Gamma}_q, \tag{22}$$

where $\boldsymbol{\Gamma}_q$ is the covariance matrix of \mathbf{e}_q . Using (20) and (21) we find that

$$\boldsymbol{\Gamma}_q = E \left[\mathbf{e}_q (\mathbf{e}_q)^H \right] = \mathbf{D}_q^H \mathbf{A}_q^T \mathbf{R}_q \mathbf{A}_q \mathbf{D}_q. \tag{23}$$

Finally,

$$\boldsymbol{\Gamma} = \sum_{q=1}^Q \mathbf{D}_q^H \mathbf{A}_q^T \mathbf{R}_q \mathbf{A}_q \mathbf{D}_q. \tag{24}$$

One can rewrite the above expression in a block-matrix form

$$\boldsymbol{\Gamma} = \mathbf{C}_Q^H \mathbf{R} \mathbf{C}_Q \tag{25}$$

where

$$\mathbf{C}_Q = [\mathbf{D}_1^H \mathbf{A}_1^T \cdots \mathbf{D}_Q^H \mathbf{A}_Q^T]^H, \tag{26}$$

and

$$\mathbf{R} = \text{diag}([\mathbf{R}_1 \cdots \mathbf{R}_Q]) \tag{27}$$

is a block-diagonal matrix with the matrices \mathbf{R}_q , $1 \leq q \leq Q$ on its diagonal, and zeros elsewhere.

Since the covariance matrices \mathbf{R}_q are full rank positive definite, the block-diagonal matrix \mathbf{R} is full rank positive definite as well. Hence by observation 7.1.6 [13] we have

$$\text{rank}(\boldsymbol{\Gamma}) = \text{rank}(\mathbf{C}_Q). \tag{28}$$

The matrix \mathbf{C}_Q has exactly NM columns, such that each one of its columns corresponds to an entry in the evanescent vector \mathbf{e} , or similarly, each one of its columns corresponds to a point on the original $N \times M$ lattice $D = \{(n, m) \in \mathbb{Z}^2 : 0 \leq n \leq N - 1, 0 \leq m \leq M - 1\}$. More specifically, the $n(M - 1) + m$ column of \mathbf{C}_Q corresponds to the $n(M - 1) + m$ element of \mathbf{e} which represents the evanescent field sample at the (n, m) lattice point. (Note that we enumerate the columns starting from zero). In the following we will adopt the abbreviation $[\mathbf{n}, \mathbf{m}]$ for indexing the $n(M - 1) + m$ column of a matrix.

To gain more understanding on the structure of \mathbf{C}_Q let us examine the different matrices \mathbf{C}_Q is comprised of. We begin with \mathbf{A}_q for some $1 \leq q \leq Q$: By construction (see (20) and the following explanation) all columns of \mathbf{A}_q are unit vectors, where the single “1” entry in each column chooses the single element from the vector \mathbf{s}_q that contributes to $e(n, m)$ – the evanescent field sample at (n, m) . The single non-zero entry in the $[n, m]$ column of \mathbf{A}_q is located in the k th row where $na_q + mb_q = k$ (we allow negative indexed rows in the case where $b_q < 0$). For example if $a_q > 0$ and $b_q > 0$, the matrix \mathbf{A}_q is given by

$$\begin{matrix}
 & [0, 0] & \cdots & [n, m] & \cdots & [N - 1, M - 1] \\
 0 & \left[\begin{array}{cccccc}
 1 & \cdots & \cdots & 0 & \cdots & \cdots & \cdots & 0 \\
 \vdots & & & \vdots & & & & \vdots \\
 k & \vdots & & 1 & & & & \vdots \\
 \vdots & & & \vdots & & & & \vdots \\
 (N - 1)a_q + (M - 1)b_q & 0 & \cdots & 0 & \cdots & \cdots & \cdots & 1
 \end{array} \right] & & & & & & &
 \end{matrix} \tag{29}$$

Let (n^*, m^*) be a solution to the linear Diophantine equation $na_q + mb_q = k$. Then, the equation is also satisfied by $n = n^* + tb_q$ and $m = m^* - ta_q$, where t is an arbitrary integer. Since (a_q, b_q) are co-prime integers these are the only possible solutions. It means that as soon as $(n + tb_q, m - ta_q) \in D$, the corresponding $[n + tb_q, m - ta_q]$ column of \mathbf{A}_q will be equal to its $[n, m]$ column. To find the rank of \mathbf{A}_q we have to evaluate the number of linearly independent columns, i.e., the number of distinct elements of $\mathbf{s}_q(na_q + mb_q)(n, m) \in D$ which contribute to \mathbf{e} .

Since \mathbf{D}_q is a diagonal matrix, the structure of $\mathbf{A}_q\mathbf{D}_q$ is similar to the structure of \mathbf{A}_q with the only difference being that instead “1” in each column, we have the appropriate exponential coefficient. Therefore, each column of the matrix \mathbf{C}_Q has exactly Q non-zero elements.

Next, let us concatenate the matrices \mathbf{A}_p and \mathbf{A}_q , where $1 \leq p \neq q \leq Q$ and examine the structure of resulting matrix

$$\tilde{\mathbf{C}}_{pq} = [\mathbf{A}_p^T \mathbf{A}_q^T]^T \tag{30}$$

As before, let us consider the structure of some arbitrary $[n, m]$ column of this matrix. It has two non-zero entries: On the k_p row of \mathbf{A}_p and on the k_q row of \mathbf{A}_q , where (n, m) satisfies

$$na_p + mb_p = k_p, \tag{31}$$

and

$$na_q + mb_q = k_q. \tag{32}$$

Next, we note that the pair $(n + tb_p, m - ta_p)$ satisfies the linear Diophantine (31) for any integer t . Therefore, for t_p such that $(n + t_p b_p, m - t_p a_p) \in D$, the $[n + t_p b_p, m - t_p a_p]$ column of \mathbf{A}_p has a “1” entry, at the same row as the $[n, m]$ column. However, $(n + t_p b_p, m - t_p a_p)$ also satisfies the linear Diophantine equation

$$(n + t_p b_p)a_q + (m - t_p a_p)b_q = \ell_q. \tag{33}$$

Hence, the $[n + t_p b_p, m - t_p a_p]$ column of \mathbf{A}_q has a “1” entry on its ℓ_q row.

Similarly, since $(n + tb_q, m - ta_q)$ satisfies the linear Diophantine Eq. (32) for any integer t , for t_q such that $(n + t_q b_q, m - t_q a_q) \in D$ we have that the $[n + t_q b_q, m - t_q a_q]$ column of \mathbf{A}_q has “1” at the same row as the $[n, m]$ column. Since,

$$(n + t_q b_q)a_p + (m - t_q a_q)b_p = \ell_p, \tag{34}$$

the $[n + t_q b_q, m - t_q a_q]$ column of \mathbf{A}_p has “1” on its ℓ_p row. Moreover, one can observe that for a pair of integers (t_p, t_q) such that $(n + t_p b_p + t_q b_q, m - t_p a_p - t_q a_q) \in D$, the pair $(n + t_p b_p + t_q b_q, m - t_p a_p - t_q a_q)$ simultaneously satisfies (33) and (34):

$$\begin{aligned}
 (n + t_p b_p + t_q b_q)a_p + (m - t_p a_p - t_q a_q)b_p &= \ell_p \\
 (n + t_p b_p + t_q b_q)a_q + (m - t_p a_p - t_q a_q)b_q &= \ell_q.
 \end{aligned} \tag{35}$$

Therefore the $[n + t_p b_p + t_q b_q, m - t_p a_p - t_q a_q]$ column of \mathbf{A}_p has “1” on its ℓ_p row, and the same column of \mathbf{A}_q has “1” on its ℓ_q row.

Finally, we can represent the $[n, m]$ column of $\tilde{\mathbf{C}}_{pq}$ by a linear combination of its other columns:

$$[n, m] = [n + t_p b_p, m - t_p a_p] + [n + t_q b_q, m - t_q a_q] - [n + t_p b_p + t_q b_q, m - t_p a_p - t_q a_q], \tag{36}$$

or in a more detailed form by

$$k_p \begin{bmatrix} 1 \\ 1 \end{bmatrix} = k_p \begin{bmatrix} 1 \\ 1 \end{bmatrix} + \ell_p \begin{bmatrix} 1 \\ 1 \end{bmatrix} - \ell_p \begin{bmatrix} 1 \\ 1 \end{bmatrix}. \tag{37}$$

Let $T_{pq}^{(n,m)}$ be the set of all the integer pairs (t_p, t_q) such that $(n + t_p b_p, m - t_p a_p)$, $(n + t_q b_q, m - t_q a_q)$, $(n + t_p b_p + t_q b_q, m - t_p a_p - t_q a_q) \in D$. Clearly, the set $T_{pq}^{(n,m)}$ is non-empty since $(0, 0) \in T_{pq}^{(n,m)}$, and it corresponds to a trivial representation of the $[n, m]$ column by itself. If $|T_{pq}^{(n,m)}| > 1$ then the $[n, m]$ column has non-trivial linear representation by other columns.

Recall however, that the matrix \mathbf{C}_Q is comprised of blocks where each block is of the form $\mathbf{A}_q \mathbf{D}_q$. Consider next the concatenation of two such blocks $\mathbf{A}_p \mathbf{D}_p$ and $\mathbf{A}_q \mathbf{D}_q$,

$$\mathbf{C}_{pq} = [\mathbf{D}_p^H \mathbf{A}_p^T \quad \mathbf{D}_q^H \mathbf{A}_q^T]^H. \tag{38}$$

Keeping in mind that by definition $a_p d_p - b_p c_p = 1$ and $a_q d_q - b_q c_q = 1$, it is easy to check that the replacement of the “1” in the columns of $\tilde{\mathbf{C}}_{pq}$ by exponentials as in (38) will only affect the coefficients of the linear combination. Indeed, the linear combination of columns in this case has the form

$$[n, m] = [n + t_p b_p, m - t_p a_p] \exp(j\omega_p t_p) + [n + t_q b_q, m - t_q a_q] \exp(j\omega_q t_q) - [n + t_p b_p + t_q b_q, m - t_p a_p - t_q a_q] \exp(j[\omega_p t_p + \omega_q t_q]). \tag{39}$$

One may also notice that if $(a_p, b_p) = (a_q, b_q)$ we have that $k_p = k_q = \ell_p = \ell_q$. However, since in this case $\omega_p \neq \omega_q$ the linear combination in (39) is still valid and non-trivial.

It is clear that the linear dependencies of columns of \mathbf{C}_Q are governed by the same simple laws: Let $T_Q^{(n,m)}$ be a set of all Q -tuples of integers (t_1, \dots, t_Q) defined as follows: For any $1 \leq q \leq Q$, let (i_1, \dots, i_q) be a set of q indices, such that $1 \leq i_k \leq Q$ for all $1 \leq k \leq q$, and

$$(n + t_1 b_{i_1} + \dots + t_q b_{i_q}, m - t_1 a_{i_1} - \dots - t_q a_{i_q}) \in D.$$

Clearly, the set $T_Q^{(n,m)}$ is non-empty since $(0, \dots, 0) \in T_Q^{(n,m)}$. Let $(n, m) \in D$ be an arbitrary lattice point and let $[n, m]$ be its corresponding column in \mathbf{C}_Q . Then, the $[n, m]$ column can be represented by the linear combination

$$[n, m] = \sum_{q=1}^Q (-1)^{q-1} \underbrace{\sum_{i_1=1}^{Q-q+1} \dots \sum_{i_q=i_{q-1}+1}^Q}_{q \text{ sums}} [n + t_1 b_{i_1} + \dots + t_q b_{i_q}, m - t_1 a_{i_1} - \dots - t_q a_{i_q}] \times \exp(j[\omega_{i_1} t_1 + \dots + \omega_{i_q} t_q]), \tag{40}$$

where $(t_1, \dots, t_Q) \in T_Q^{(n,m)}$. The details of this derivation are presented in [Appendix](#).

Following the foregoing analysis of the linear dependencies between the columns of \mathbf{C}_Q , we next count its linearly independent columns in order to derive the rank of \mathbf{C}_Q . Let us first count the number of independent columns of $\mathbf{A}_q \mathbf{D}_q$. As mentioned earlier, this number is equal the number of distinct samples from $s_q(na_q + mb_q)$, $(n, m) \in D$ that contribute to \mathbf{e} . In other words, this is the number of different indices k , such that $na_q + mb_q = k$ where $(n, m) \in D$, and it can be easily calculated based on the dimensions of D (see [Fig. 3](#) for an illustrative example). Indeed, a new sample from the random process $\{s_q\}$ may be introduced only on the first a_q rows (since $a_q \geq 0$) and the last (first) $|b_q|$ columns (last if $b_q \geq 0$ and first if $b_q < 0$) of the observed finite dimensional field, while the rest of the field is filled by replicas of these samples. We thus count Na_q distinct samples in the first a_q rows and $M|b_q|$ distinct samples in the first (last) $|b_q|$ columns. However, on the intersection of these rows and columns $|a_q b_q|$ samples are counted twice. Finally, the total number of *distinct* samples from the random process $\{s_q\}$ that are found in an observed field of dimensions $N \times M$ (which is equal to the rank of \mathbf{A}_q and the rank of $\mathbf{A}_q \mathbf{D}_q$) is given by

$$r_q = Na_q + M|b_q| - |a_q b_q|. \tag{41}$$

Similarly, it can be shown that the number of linearly independent columns of $\mathbf{A}_p \mathbf{D}_p$ is $Na_p + M|b_p| - |a_p b_p|$. Let us next count the number of linearly independent columns of \mathbf{C}_{pq} .

Since r_p columns of $\mathbf{A}_p \mathbf{D}_p$ are linearly independent, the same columns of the concatenated matrix \mathbf{C}_{pq} are linearly independent as well. The remaining $NM - r_p = (N - |b_p|)(M - |a_p|)$ columns may be considered to correspond to an $(N - |b_p|) \times (M - |a_p|)$ rectangular sub-lattice $D_1 = \{(n, m) \in \mathbb{Z}^2 : 0 \leq n \leq N - 1 - |b_p|, |a_p| \leq m \leq M - 1\}$, which is a subset of the original rectangular lattice (or similarly, one can define $D_1 = \{(n, m) \in \mathbb{Z}^2 : |b_p| \leq n \leq N - 1, |a_p| \leq m \leq M - 1\}$

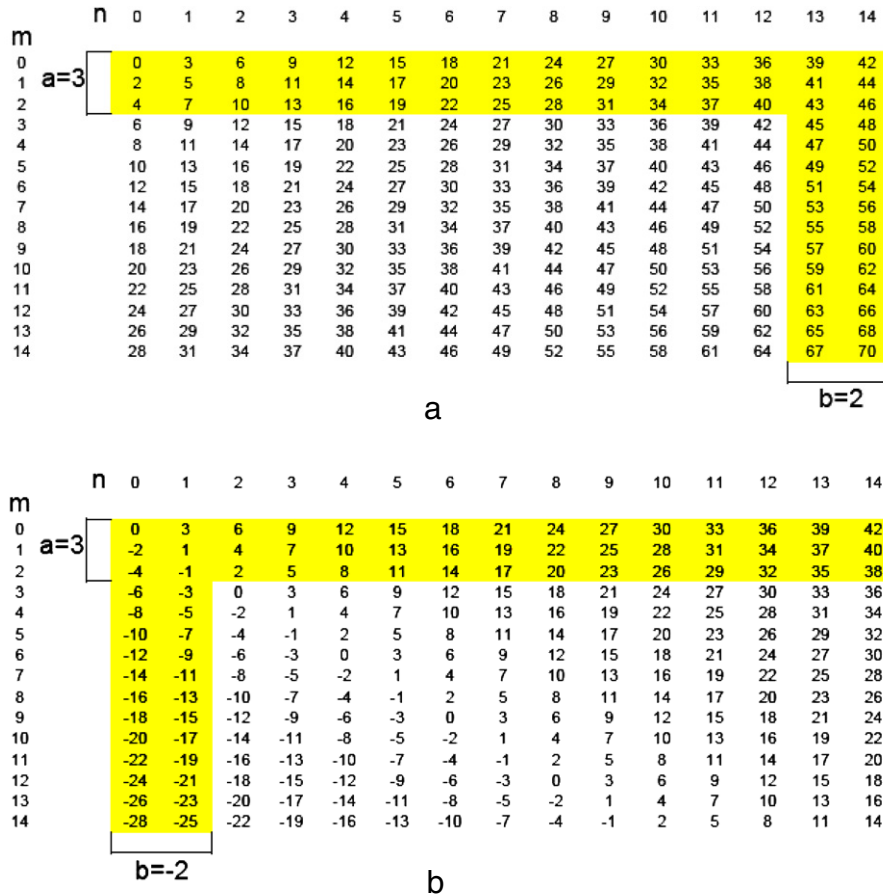


Fig. 3. $N = M = 15$, $a = 3$, $b = \pm 2$: (a) The indices $k = na + mb$ of the observation on $\{s^{(3,2)}(k)\}$; (b) The indices $k = na + mb$ of the observation on $\{s^{(3,-2)}(k)\}$. The sets of distinct indices of $\{s^{(3,2)}(k)\}$ and $\{s^{(3,-2)}(k)\}$ are marked in yellow (gray). Every other sample in the field is identical to some sample in the yellow area.

which does not change the reasoning of our arguments and only depends on a sign of b_p). See Fig. 4 as an example. Repeating the same arguments as those made above, one can show that the number of *distinct* samples from the random process $\{s_q\}$ that are found in a sub-lattice D_1 is

$$\tilde{r}_q = (N - |b_p|)|a_q| + (M - |a_p|)|b_q| - |a_q b_q|. \tag{42}$$

This is the number of linearly independent columns which can be found in C_{pq} in addition to the first r_p columns.

Let D_2 be the set of $NM - r_p - \tilde{r}_q = (N - |b_p| - |b_q|)(M - |a_p| - |a_q|)$ lattice points that remain after the removal from D of the $r_p + \tilde{r}_q$ points corresponding to the linearly independent columns of C_{pq} (for simplicity and without limiting of the generality of the results, we will discuss the case where $b_p > 0$ and $b_q > 0$, as illustrated in Fig. 4(a) (uncolored area)). Thus, $D_2 = \{(n, m) \in \mathbb{Z}^2 : 0 \leq n \leq N - 1 - |b_p| - |b_q|, |a_p| + |a_q| \leq m \leq M - 1\}$. It thus remains to be shown that all columns representing points in D_2 can be represented by a linear combination of columns that correspond to points in $D \setminus D_2$.

Since the “width” of $D \setminus D_2$ is $|b_p| + |b_q|$ along the n -axis and $|a_p| + |a_q|$ along the m -axis (colored areas in Fig. 4(a)), for every $(n, m) \in D_2$ we have $(n + b_p, m - a_p), (n + b_q, m - a_q), (n + b_p + b_q, m - a_p - a_q) \in D$. Thus, $(t_p, t_q) = (1, 1) \in T_{pq}^{(n,m)}$, and as we have shown above, we can represent $[n, m]$ by the linear combination

$$\begin{aligned} [n, m] &= [n + b_p, m - a_p] \exp(j\omega_p) + [n + b_q, m - a_q] \exp(j\omega_q) \\ &\quad - [n + b_p + b_q, m - a_p - a_q] \exp(j[\omega_p + \omega_q]). \end{aligned} \tag{43}$$

Continuing this construction recursively, it is obvious that for each point $(n, m) \in D_2$ we can find a pair $(t_p, t_q) \in T_{pq}^{(n,m)}$, and $(t_p, t_q) \neq (0, 0)$ such that $(n + t_p b_p, m - t_p a_p), (n + t_q b_q, m - t_q a_q), (n + t_p b_p + t_q b_q, m - t_p a_p - t_q a_q) \in D \setminus D_2$. On the other hand, for every point $(n, m) \in D \setminus D_2$ one can show that $T_{pq}^{(n,m)} = \{(0, 0)\}$, i.e., only the trivial linear combination exists. In other words, all the random variables indexed on $D \setminus D_2$ correspond to linearly independent columns. Therefore,

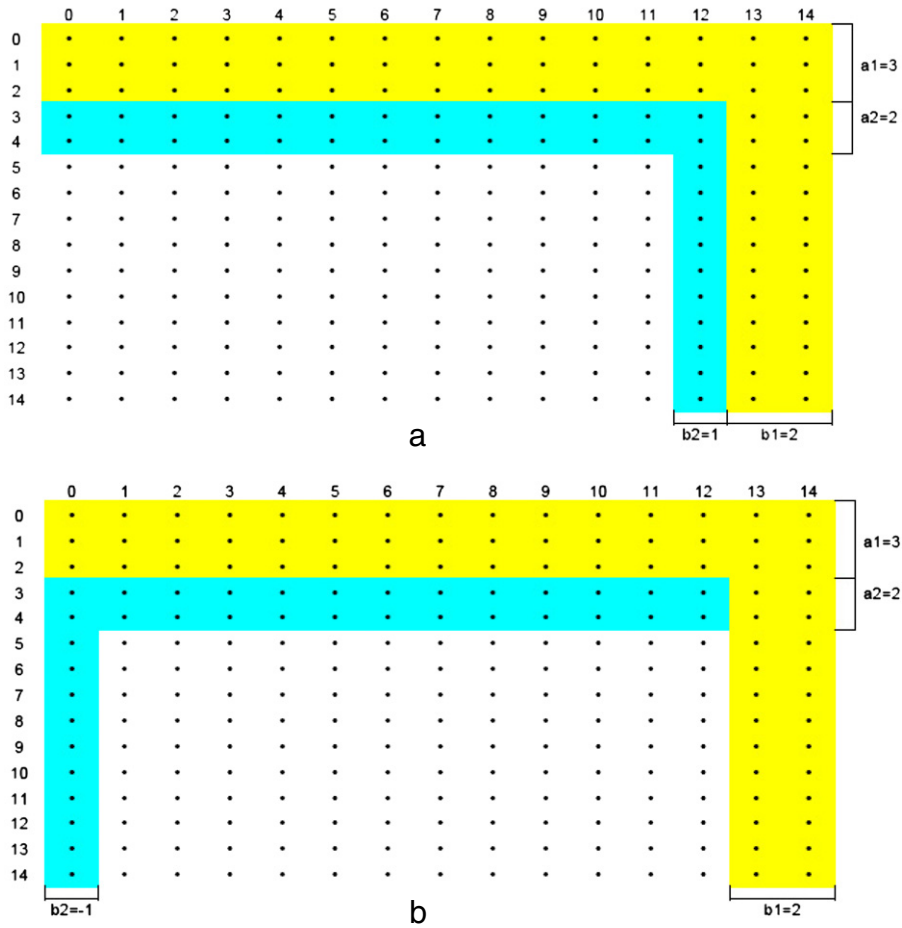


Fig. 4. $N = M = 15, a_1 = 3, b_1 = 2, a_2 = 2, b_2 = \pm 1$. (a): The set of distinct samples of $s^{(3,2)}(n, m)$ (in yellow) and $s^{(2,1)}(n, m)$ (in blue). (b): The set of distinct samples of $s^{(3,2)}(n, m)$ (in yellow or lightgray) and $s^{(2,-1)}(n, m)$ (in blue or darkgray). In both cases, every other sample in the field is a linear combination of samples in the colored areas.

the number of linearly independent columns in \mathbf{C}_{pq} is

$$\begin{aligned} \text{rank}(\mathbf{C}_{pq}) &= |D \setminus D_2| = r_p + \tilde{r}_q \\ &= N(|a_p| + |a_q|) + M(|b_p| + |b_q|) - (|a_p| + |a_q|)(|b_p| + |b_q|). \end{aligned} \tag{44}$$

The construction described above can be easily extended to the general case where we concatenate *all* the matrices which \mathbf{C}_Q is comprised of. See Fig. 5 for an example of a three component case. If one chooses a subset of the original lattice,

$$D_Q = \left\{ (n, m) \in \mathbb{Z}^2 : 0 \leq n \leq N - 1 - \sum_{q=1}^Q |b_q|, \sum_{q=1}^Q |a_q| \leq m \leq M - 1 \right\},$$

which remains after the removal of

$$N \sum_{q=1}^Q |a_q| + M \sum_{q=1}^Q |b_q| - \sum_{q=1}^Q |a_q| \sum_{q=1}^Q |b_q| \tag{45}$$

lattice points (similarly to D_2 which remains after the removal of the $r_p + \tilde{r}_q$ lattice points corresponding to the independent columns of \mathbf{C}_{pq}), one can repeat the same considerations as above and show that columns of \mathbf{C}_Q that correspond to the lattice points in $D \setminus D_Q$ are the *only* linearly independent columns of \mathbf{C}_Q . Thus,

$$\text{rank}(\mathbf{C}_Q) = |D \setminus D_Q| = N \sum_{q=1}^Q |a_q| + M \sum_{q=1}^Q |b_q| - \sum_{q=1}^Q |a_q| \sum_{q=1}^Q |b_q|. \tag{46}$$

Finally, Since the rank of $\mathbf{\Gamma}$ cannot exceed NM (the dimension of the covariance matrix), NM is an upper bound on the rank of $\mathbf{\Gamma}$. Combining this and (46) the rank of $\mathbf{\Gamma}$ is given by (9), which completes the proof.

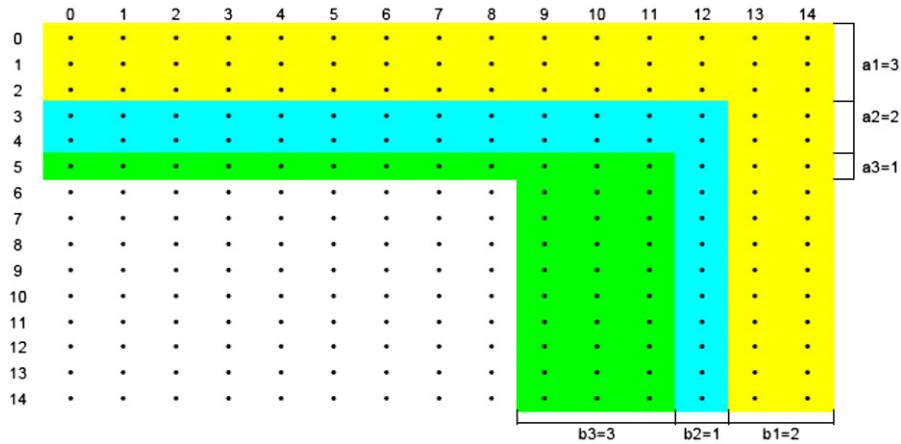


Fig. 5. Sets of distinct samples in the case of three evanescent components. $N = M = 15$, $a_1 = 3$, $b_1 = 2$, $a_2 = 2$, $b_2 = 1$, $a_3 = 1$, $b_3 = 3$.

4. The case of a real-valued evanescent field

In the case where a real-valued evanescent field is considered, we have

$$e_q(n, m) = s_q(na_q + mb_q) \cos(\omega_q(nc_q + md_q)) + t_q(na_q + mb_q) \sin(\omega_q(nc_q + md_q)), \tag{47}$$

where the 1D purely-indeterministic processes $\{s_q\}$, $\{s_p\}$, $\{t_q\}$, $\{t_p\}$ are mutually orthogonal for all $1 \leq p \neq q \leq Q$, and for all q the processes $\{s_q\}$ and $\{t_q\}$ have an identical autocorrelation function. Let \mathbf{t}_q be defined similarly to \mathbf{s}_q in (18). Using similar notations as in (20) we have

$$\mathbf{e}_q = \mathcal{R}(\mathbf{D}_q^H) \mathbf{A}_q^T \mathbf{s}_q + \mathcal{I}(\mathbf{D}_q^H) \mathbf{A}_q^T \mathbf{t}_q, \tag{48}$$

where \mathcal{R} and \mathcal{I} denote real and imaginary parts respectively. Finally, since the processes $\{s_q\}$ and $\{t_q\}$ are mutually orthogonal and have an identical autocorrelation function we find that

$$\mathbf{\Gamma}_q = E \left[\mathbf{e}_q (\mathbf{e}_q)^T \right] = \mathcal{R}(\mathbf{D}_q^H) \mathbf{A}_q^T \mathbf{R}_q \mathbf{A}_q \mathcal{R}(\mathbf{D}_q) + \mathcal{I}(\mathbf{D}_q^H) \mathbf{A}_q^T \mathbf{R}_q \mathbf{A}_q \mathcal{I}(\mathbf{D}_q) \tag{49}$$

where

$$\mathbf{R}_q = E \left[\mathbf{s}_q (\mathbf{s}_q)^T \right] = E \left[\mathbf{t}_q (\mathbf{t}_q)^T \right] \tag{50}$$

is positive definite since $\{s_q\}$ and $\{t_q\}$ are purely-indeterministic.

Similarly to the case of a complex-valued evanescent field, the covariance matrix $\mathbf{\Gamma}$ is given by

$$\mathbf{\Gamma} = \sum_{q=1}^Q \mathbf{\Gamma}_q. \tag{51}$$

The derivation of the rank of the covariance matrix (51) follows exactly the same lines as in the previous section, and the next corollary is immediate:

Corollary 1. Let \mathbf{e} be a vector-form representation of a finite sample from a sum of real-valued evanescent random fields, given by (3), (5) and (47). Then, the rank of its covariance matrix, $\mathbf{\Gamma}$, is given by

$$\text{rank}(\mathbf{\Gamma}) = \min \left(NM, \left[N \sum_{q=1}^Q 2|a_q| + M \sum_{q=1}^Q 2|b_q| - \sum_{q=1}^Q 2|a_q| \sum_{q=1}^Q 2|b_q| \right] \right). \tag{52}$$

5. Conclusion

We have considered the problem of evaluating the rank of the covariance matrix of a finite sample from an evanescent random field. We have analytically derived the rank formula and have shown that the rank of the covariance matrix of this finite sample from the evanescent random field is completely determined by the evanescent field spectral support parameters and is independent of all other parameters of the field. Thus, for example, the problem of evaluating the rank of

for canceling $k \binom{Q}{k}$ undesired elements created in the previous step. Due to this subtraction/addition new $\binom{Q}{k} (Q - k) = (k + 1) \binom{Q}{k+1}$ undesired elements are created. Clearly, when we subtract/add $\binom{Q}{Q-1}$ columns exactly $Q = Q \binom{Q}{Q}$ undesired elements are created. These may be canceled by subtraction of a single vector. By subtraction/addition of the last vector, $[n + t_1 b_1 + \dots + t_Q b_Q, m - t_1 a_1 - \dots - t_Q a_Q]$ the process terminates, since we are canceling the last Q undesired elements and remain with Q elements – exactly those of the $[n, m]$ column, i.e.,

$$\begin{aligned}
 [n, m] &= \sum_{i=1}^Q [n + t_i b_i, m - t_i a_i] \exp(j\omega_i t_i) - \sum_{i=1}^{Q-1} \sum_{j=i+1}^Q [n + t_i b_i + t_j b_j, m - t_i a_i - t_j a_j] \exp(j[\omega_i t_i + \omega_j t_j]) + \dots \\
 &\quad + (-1)^{Q-1} [n + t_1 b_1 + \dots + t_Q b_Q, m - t_1 a_1 - \dots - t_Q a_Q] \exp(j[\omega_1 t_1 + \dots + \omega_Q t_Q]) \\
 &= \sum_{q=1}^Q (-1)^{q-1} \underbrace{\sum_{i_1=1}^{Q-q+1} \dots \sum_{i_q=i_{q-1}+1}^Q}_{q \text{ sums}} [n + t_{i_1} b_{i_1} + \dots + t_{i_q} b_{i_q}, m - t_{i_1} a_{i_1} - \dots - t_{i_q} a_{i_q}] \\
 &\quad \times \exp(j[\omega_{i_1} t_{i_1} + \dots + \omega_{i_q} t_{i_q}]). \tag{58}
 \end{aligned}$$

Clearly, this linear combination will be meaningful only if $(t_1, \dots, t_Q) \in T_Q^{(n,m)}$, i.e., $(n + t_{i_1} b_{i_1} + \dots + t_{i_q} b_{i_q}, m - t_{i_1} a_{i_1} - \dots - t_{i_q} a_{i_q}) \in D$, for any $1 \leq q \leq Q$, and where (i_1, \dots, i_q) is such that $1 \leq i_k \leq Q$ for all $1 \leq k \leq q$.

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