

Blind and Semiblind Detections of OFDM Signals in Fading Channels

Ming-Xian Chang, *Member, IEEE*, and Yu T. Su, *Member, IEEE*

Abstract—This paper considers the problem of blind joint channel estimation and data detection for orthogonal frequency-division multiplexing (OFDM) systems in a fading environment. Employing a regression model for a time-varying channel, we convert the problem into one that finds the data sequence \mathbf{x} whose associated least-squares (LS) channel estimate $\mathbf{z}(\mathbf{x})$ is closest to the space of some regression curves (surfaces). We apply the branch-and-bound principle to solve the nonlinear integer programming problem associated with finding the curve that fits a subchannel in the LS sense. A recursive formula for fast metric update is obtained by exploiting the intrinsic characteristic of our objective function. The impacts of reordering the data sequence and selective detection are addressed. By employing a preferred order along with a selective detection method, we greatly reduce the detector complexity while giving up little performance loss. Both the complete and the reduced-complexity algorithms can be used for blind and semiblind detections of OFDM signals in a subchannel-by-subchannel manner. To further reduce the complexity and exploit the frequency-domain channel correlation, we suggest a two-stage approach that detects a few selected positions in some subchannels first, and then, treating the detected symbols as pilots, determines the remaining symbols within a properly chosen time-frequency block by a two-dimensional model-based pilot-assisted algorithm. The proposed methods do not require the information of the channel statistics like signal-to-noise ratio or channel correlation function. Performance of differential modulations like differential quaternary phase-shift keying and STAR 16-ary quadrature amplitude modulation are provided. Both blind and semiblind schemes yield satisfactory performance.

Index Terms—Estimation, fading channels, signal detection.

I. INTRODUCTION

BROADBAND wireless transmission often calls for complicated equalization circuits to compensate for severe frequency-selective fading. By transforming a wideband signal into an array of orthogonal narrowband signals, the orthogonal frequency-division multiplexing (OFDM) scheme converts a wideband frequency-selective fading channel into a bank of

subchannels that suffer from frequency-nonselective fading only.

A transmitted OFDM symbol period usually consists of a guard interval (cyclic prefix) and a regular symbol interval. We shall assume that the intersymbol interference (ISI) due to channel multipath is completely eliminated, after removing the received baseband samples in the guard interval and taking the discrete Fourier transform (DFT) on the remaining samples in the regular symbol interval. The received sample of the m th subchannel at the n th regular symbol interval is then given by [1]–[5]

$$Y_{mn} = H_{mn}X_{mn} + N_{mn} \quad (1)$$

where Y_{mn} and X_{mn} are, respectively, the received and transmitted symbols, H_{mn} accounts for the channel effect, or the channel (frequency) response (CR), and $\{N_{mn}\}$ are independent and identically distributed (i.i.d.) complex zero-mean Gaussian random variables. The above model also assumes that the receiver's frequency and timing recovery subsystem is such that, within the time span of interest, the effect of the residual error can be accounted for by the CR H_{mn} . To detect the transmitted data sequence $\{X_{mn}\}$, it is essential that the receiver obtain reliable estimates of the CRs.

Many pilot-based methods for estimating H_{mn} have been proposed [1]–[5]. Proposals of blind and semiblind OFDM channel estimation schemes are also plentiful [6]–[11]. [6] and [7] exploited the cyclostationarity, while [8] used the finite-alphabet property of the transmitted symbols and virtual carriers to identify the channel impulse response. The finite-alphabet property is also profitably used by Zhou and Giannakis [9] to develop efficient channel estimators. These statistical blind algorithms often require a period of hundreds or even thousands of symbol blocks to obtain the corresponding estimator, and they assume that the CR is static during this period. In many situations, such an assumption is invalid, and H_{mn} may, in fact, vary from block to block. The solution suggested by Luise *et al.* [10] is applicable for time-varying channels that vary from block to block, but it works for constant module signals only, and requires a few pilot symbols at both ends of a block. Their proposal is based on the so-called super-trellis and the per-survivor processing (PSP) algorithm [12] that requires relatively high computing complexity. The model-based approach of [5] and [11] can also handle time-varying channels and, like [10], does not need channel statistics such as noise level and channel correlation function. This paper presents a new modest-complexity blind joint OFDM channel estimation and data detection algorithm, adopting the model-based approach

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of [5] and [11]. We remark that [10] also uses a similar channel model in the frequency domain to predict CR for use in the PSP Viterbi detection algorithm.

Using a regression model to describe the CR process, the following section shows how the task of joint channel estimation and data detection becomes that of selecting the regression coefficients and the associated data sequence to minimize a cost function. As the optimal regression coefficient vector is a function of the data block of concern, we only have to deal with the remaining blind data detection issue that can be formulated as a nonlinear integer programming problem. In Section III, we develop an efficient algorithm to solve this problem, applying the branch-and-bound principle to the associated detection tree and deriving a simple recursive relation between node values of successive tree nodes. Section IV discusses the ambiguity issue and the influence of the tree-searching order. We show that the algorithm's complexity can be much reduced by selecting a properly ordered subset of data positions for tree searching. We then extend our investigation to the semiblind case, where only sparse pilot symbols are available. Section V gives some simulation examples of the performance of various OFDM systems and detectors. Finally, Section VI summarizes our main results.

II. BLIND DATA DETECTION

A. Problem Formulation

The two-dimensional (2-D) process $\{H_{mn}\}$ in (1) accounts for the channel (fading) effect that can be modeled as a 2-D regression surface [4], [5]

$$H_{mn} = \hat{H}_{mn} + g(m, n) = F(m, n, \mathbf{c}) + g(m, n) \quad (2)$$

where $\hat{H}_{mn} = F(m, n, \mathbf{c})$ is our model for H_{mn} within the time-frequency block, $0 \leq n < L$, $0 \leq m < M_0$, M_0 being less than the total subchannel (subcarrier) number, and $g(m, n)$ represents the modeling (approximation) error. A candidate 2-D surface is

$$F(m, n, \mathbf{c}) = c_0 \cdot m^2 + c_1 \cdot mn + c_2 \cdot n^2 + c_3 \cdot m + c_4 \cdot n + c_5 \quad (3)$$

where $\mathbf{c} = (c_0, c_1, \dots)^T$ is the coefficient vector. When the data symbols $\{X_{mn}\}$ are not known, we should select $\{X_{mn}\}$ that solves

$$\min_{\mathbf{X}} \min_{\mathbf{c}} \sum_{(m,n) \in \mathcal{I}} \left| \frac{Y_{mn}}{X_{mn}} - F(m, n, \mathbf{c}) \right|^2 \quad (4)$$

where $\mathcal{I} = \{(m, n) | 0 \leq m < M_0, 0 \leq n < L\}$, $\mathbf{X} = \{X_{mn} | (m, n) \in \mathcal{I}\}$. For a given data set \mathbf{X} , the corresponding optimal \mathbf{c} is

$$\mathbf{c} = \left(\sum_{(m,n) \in \mathcal{I}} \mathbf{q}_{mn} \mathbf{q}_{mn}^T \right)^{-1} \left(\sum_{(m,n) \in \mathcal{I}} \mathbf{q}_{mn} \frac{Y_{mn}}{X_{mn}} \right) \quad (5)$$

where $\mathbf{q}_{mn} \triangleq (m^2, mn, n^2, m, n, 1)^T$ if (3) is adopted as the channel model. Define

$$\mathbf{P} = \left(\sum_{(m,n) \in \mathcal{I}} \mathbf{q}_{mn} \mathbf{q}_{mn}^T \right)^{-1}. \quad (6)$$

Substituting (5) into (4), we have the optimal data (vector) estimate

$$\hat{\mathbf{X}} = \arg \left\{ \min_{\mathbf{X}} \sum_{(m,n) \in \mathcal{I}} \left| \frac{Y_{mn}}{X_{mn}} - \sum_{(k,l) \in \mathcal{I}} \mathbf{q}_{mn}^T \mathbf{P} \mathbf{q}_{kl} \frac{Y_{kl}}{X_{kl}} \right|^2 \right\} \quad (7)$$

which can be obtained by performing an exhaustive search on all possible \mathbf{X} . However, such a primitive approach is practical only if the cardinality of \mathbf{X} is very small.

Although we have reduced the twofold uncertainty $\{\mathbf{c}, \mathbf{X}\}$ to the singlefold uncertainty $\{\mathbf{X}\}$, the task of estimating the solution data matrix \mathbf{X} , (7), is still enormous. We shall only consider the case of one-dimensional (1-D) data sequence detection in the sequel. The procedure of solving the 2-D problem, (4), can be conveniently divided into two stages. In the first stage, we try to solve the 1-D version of (4) for some selected subchannels and some selected data positions. In the second stage, we employ a pilot-assisted, 2-D regression-model-based method [4], using the symbols detected in the first stage as pilots to determine the remaining symbols; see Section V-C.

For a given subchannel (or time slot), setting the corresponding parameter m (or n) to a constant, we obtain a projection of the surface (3), a polynomial in n (or m). The sampled CR process of a subchannel in an interval of L OFDM symbols can thus be modeled as

$$\hat{H}_n = F(n, \mathbf{c}) = c_0 \cdot n^2 + c_1 \cdot n + c_2, \quad 0 \leq n < L \quad (8)$$

where we have neglected the channel index m for conciseness. Equations (4) and (8) indicate that blind equalization of a subchannel needs to solve

$$\min_{\mathbf{x}} \min_{\mathbf{c}} \sum_{n=0}^{L-1} \left| \frac{Y_n}{X_n} - F(n, \mathbf{c}) \right|^2 = \min_{\mathbf{x}} \min_{\mathbf{c}} \sum_{n=0}^{L-1} \left| \frac{Y_n}{X_n} - \mathbf{q}_n^T \mathbf{c} \right|^2 \quad (9)$$

where $\mathbf{x} \triangleq (X_0, X_1, \dots, X_{L-1})^T$, $\mathbf{q}_n \triangleq (n^2, n, 1)^T$, and $\mathbf{c} = (c_0, c_1, c_2)^T$. Defining $\mathbf{z} = \mathbf{z}(\mathbf{x}) \triangleq (Y_0/X_0, Y_1/X_1, \dots, Y_{L-1}/X_{L-1})^T$, and $\mathbf{Q} \triangleq (\mathbf{q}_0^T, \mathbf{q}_1^T, \dots, \mathbf{q}_{L-1}^T)^T$, we rewrite (9) as

$$\min_{\mathbf{x}} \min_{\mathbf{c}} \|\mathbf{z} - \mathbf{Q} \mathbf{c}\|^2 \quad (10)$$

where $\|\cdot\|$ denotes the vector norm. For a fixed \mathbf{x} (and hence, \mathbf{z}), we have

$$\mathbf{c}_{\text{opt}} = \arg \left(\min_{\mathbf{c}} \|\mathbf{z} - \mathbf{Q} \mathbf{c}\|^2 \right) = (\mathbf{Q}^T \mathbf{Q})^{-1} \mathbf{Q}^T \mathbf{z}. \quad (11)$$

Substituting (11) into (10) and defining the constant symmetrical real matrix

$$\mathbf{V} \triangleq \mathbf{Q} (\mathbf{Q}^T \mathbf{Q})^{-1} \mathbf{Q}^T \quad (12)$$

we obtain

$$\begin{aligned} \hat{\mathbf{x}} &= \arg \left(\min_{\mathbf{z}(\mathbf{x})} \|\mathbf{z}(\mathbf{x}) - \mathbf{V} \mathbf{z}(\mathbf{x})\|^2 \right) \\ &= \arg \left[\min_{\mathbf{x}} (\mathbf{z}^H \mathbf{z} - \mathbf{z}^H \mathbf{V} \mathbf{z}) \right]. \end{aligned} \quad (13)$$

Note that if \mathbf{x} is known, then $\mathbf{z}(\mathbf{x})$ is the corresponding least-squares (LS) estimate of the CRs. The above equation therefore

implies that the desired solution is the data vector \mathbf{x} that minimizes the squared Euclidean distance between the LS channel estimate associated with it, $\mathbf{z}(\mathbf{x})$, and the corresponding orthogonal projection $\mathbf{V}\mathbf{z}$ onto the space of curves defined by complex polynomials of degree equal or less than d [e.g., $d = 2$ in (8)]. It is obvious that the same conclusion can be extended to the case of the space of higher degree (> 2) polynomial curves or other proper linear spaces if we use other regression models, parametric or not. Similar interpretation can be given to (7) with the space now consisting of quadratic surfaces instead. The elements X_k 's of \mathbf{x} belong to a finite symbol set \mathcal{S} that depends on the modulation scheme. A straightforward approach to solve (13) is trying all candidate data sequences \mathbf{x} but, as in the 2-D case, such an exhaustive search quickly becomes impractical as L becomes large. Before addressing the issue of solving (13), let us exploit some properties of the orthogonal projection operator \mathbf{V} .

B. Eigenspace Decomposition

Note that each column of \mathbf{Q} is an eigenvector of \mathbf{V} with eigenvalue one, since

$$\mathbf{V}\mathbf{Q} = \mathbf{Q}(\mathbf{Q}^T\mathbf{Q})^{-1}\mathbf{Q}^T\mathbf{Q} = \mathbf{Q}. \quad (14)$$

If the regression polynomial is of degree d , then \mathbf{Q} is an $L \times (d+1)$ real matrix, where the block length L is often much larger than the polynomial degree d . As (12) implies that the rank of \mathbf{V} is, at most, $(d+1)$, the span of the $(d+1)$ independent columns of \mathbf{Q} forms an eigenspace of \mathbf{V} in the L -dimensional real space \mathcal{R}^L , the corresponding eigenvalue being one. This eigenspace, \mathcal{P}_L^d , consists of "curves" defined by the polynomials whose degrees are less than $d+1$, and whose domain is the discrete set $\{0, 1, \dots, L\}$, i.e.,

$$\mathcal{P}_L^d \triangleq \left\{ (y_0, y_1, \dots, y_{L-1}) \mid y_n = \sum_{i=0}^d c_i n^i; c_i \in \mathcal{R}, \right. \\ \left. 0 \leq n < L, 0 \leq i \leq d \right\}.$$

On the other hand, the null space of the linear operator defined by the matrix \mathbf{V} is the eigenspace associated with the eigenvalue zero. It can easily be seen that \mathbf{V} is an orthogonal projection that maps a vector in \mathcal{R}^L into a vector (curve) in \mathcal{P}_L^d , a $(d+1)$ -dimensional space of polynomial curves. Hence, the real and imaginary part of $\mathbf{V}\mathbf{z}$ are, respectively, the projection of the real and imaginary part of \mathbf{z} .

By performing the Gram-Schmidt process on the columns of \mathbf{Q} , we obtain a set of orthonormal eigenvectors $\{\mathbf{e}_0, \mathbf{e}_1, \dots, \mathbf{e}_d\}$ of \mathbf{V} . This set can be extended to form an orthonormal basis of \mathcal{R}^L , $\{\mathbf{e}_0, \mathbf{e}_1, \dots, \mathbf{e}_{L-1}\}$, where $\{\mathbf{e}_0, \dots, \mathbf{e}_d\}$ is a basis of the range of \mathbf{V} while $\{\mathbf{e}_{d+1}, \dots, \mathbf{e}_{L-1}\}$ is a basis for the null space. Let

$$\mathbf{z}(\mathbf{x}) = \sum_{k=0}^d c_k \mathbf{e}_k + \sum_{k=d+1}^{L-1} c_k \mathbf{e}_k = \sum_{k=0}^d c_k \mathbf{e}_k + \mathbf{w}(\mathbf{x}) \quad (15)$$

where \mathbf{w} represents the component of \mathbf{z} in the null space. Since \mathbf{V} has the representation, $\mathbf{V} = \sum_{k=0}^d \mathbf{e}_k \mathbf{e}_k^T$ we can rewrite (13) as

$$\hat{\mathbf{x}} = \arg \left[\min_{\mathbf{x}} (\mathbf{z}^H(\mathbf{x})\mathbf{z}(\mathbf{x}) - \mathbf{z}^H(\mathbf{x})\mathbf{V}\mathbf{z}(\mathbf{x})) \right] \\ = \arg \left[\min_{\mathbf{x}} (\|\mathbf{w}(\mathbf{x})\|^2) \right] \quad (16)$$

i.e., choosing \mathbf{x} that minimizes the norm of the orthogonal projection of its LS channel estimate, \mathbf{z} , in the null space of \mathbf{V} , i.e., the orthogonal complement of \mathcal{P}_L^d .

The above discussion suggests a simple algorithm to find an approximate solution of (16). Obviously, without restricting the domain of \mathbf{x} , there will be an infinite number of solutions, i.e., those \mathbf{x} whose corresponding LS channel estimate $\mathbf{z}(\mathbf{x})$ lies in the range space of \mathbf{V} . We, therefore, seek for the closest approximation of \mathbf{z} , i.e., $\mathbf{z} \approx \sum_{k=0}^d c_k \mathbf{e}_k$. The solution of the approximation

$$\tilde{\mathbf{z}} = c_0 \begin{bmatrix} | \\ \mathbf{e}_0 \\ | \end{bmatrix} + c_1 \begin{bmatrix} | \\ \mathbf{e}_1 \\ | \end{bmatrix} + \dots + c_d \begin{bmatrix} | \\ \mathbf{e}_d \\ | \end{bmatrix} \approx \begin{bmatrix} \frac{Y_0}{X_0} \\ \vdots \\ \frac{Y_{L-1}}{X_{L-1}} \end{bmatrix} \quad (17)$$

can be readily found by the following steps.

- 1) Pick proper $(d+1)$ positions $i_0, \dots, i_d \in \{0, 1, \dots, L-1\}$. The proper positions should be uniform in $0, \dots, L$; see Section IV-B.
- 2) Find the coefficients (c_0, \dots, c_d) in (17) for the candidate subsequence $(X_{i_0}, \dots, X_{i_d})$ and compute the associated $\tilde{\mathbf{z}}$.
- 3) Determine the remaining X_k 's by performing hard decision on Y_i/\tilde{z}_i , where \tilde{z}_i is the i th component of $\tilde{\mathbf{z}}$.
- 4) Compute the magnitude of the resulting orthogonal projection \mathbf{w} in (15).
- 5) Repeat the above procedure 1)-4) for all candidate subsequences $(X_{i_0}, \dots, X_{i_d})$'s and choose the one whose associated "complete" sequence, \mathbf{x} , has the minimum norm orthogonal complement.

Numerical experiments show that if we pick $(d+1)$ proper positions and if the signal-to-noise ratio (SNR) is high, the above algorithm usually leads to the correct estimate. In general, a judicious choice of the initial positions i_0, \dots, i_d is needed to ensure a good approximation. When $d = 2$ or 3 is sufficient to model the subchannel variations within a properly selected time interval, which is true for many cases of interest, the number of candidate subsequences is far smaller than that required by an exhaustive search on the class of "complete" sequences. The utility of this algorithm, however, lies in the fact that it gives us a good initial estimate and a reference metric value to be used in the main algorithm, presented in the following section.

III. BLIND DETECTION ALGORITHM

As the elements of a candidate data sequence \mathbf{x} in our blind detection problem belong to a finite set of modulation symbols, the optimization of (13) is analogous to an integer programming problem. In this section, we translate a general integer programming problem into a tree-searching problem and present an efficient tree-searching algorithm. The algorithm is based on the

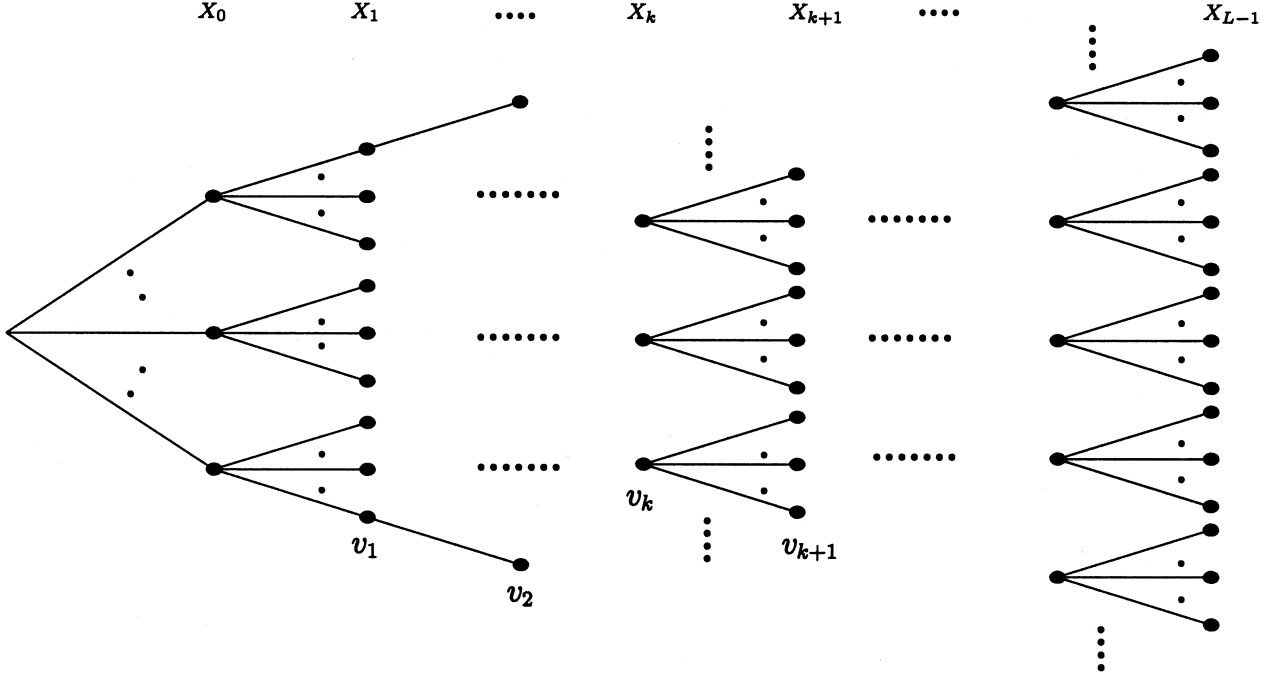


Fig. 1. Tree that represents all possible candidate sequences and the node values, $\{v_k, k = 0, 1, \dots\}$, associated with nodes in the tree.

basic principle of branch-and-bound [16], and uses an unconstrained solution obtained by using some inherent characteristic of the associated objective function. We also derive a recursive formula for the unconstrained minimum value, greatly reducing the complexity of our algorithm.

A. General Problem and Algorithm

Consider the general problem. Find the minimum value of a multivariate real-valued objective function $f(X_0, X_1, \dots, X_{L-1})$, where $X_k \in \mathcal{S} = \{s_1, s_2, \dots, s_{|\mathcal{S}|}\}$, $|\mathcal{S}| < \infty$ being the cardinality of \mathcal{S} . We call a sequence $\mathbf{x} = (X_0, X_1, \dots, X_{L-1})^T$, $X_k \in \mathcal{S}$ a candidate sequence. A straightforward approach to obtain the solution is evaluating $f(\mathbf{x})$ for all candidate sequences $\mathbf{x} = (X_0, \dots, X_{L-1})^T$. Fig. 1 plots a tree related to this search problem. Every node in the tree has $|\mathcal{S}|$ branches stemming from it, each connecting to a different descendent node. The $|\mathcal{S}|$ descendent nodes of the root (i.e., the leftmost node of the tree) represent all possible values of X_0 taken from \mathcal{S} , and are called nodes of depth 0. In general, a node of depth $k-1$ has $|\mathcal{S}|$ descendent nodes of depth k , each denotes a possible value of X_k , and the depth of the tree associated with the objective function f is $L-1$. A complete path in the tree is a path that starts at the root and ends at the tree tip. Each complete path contains nodes of depth 0 to $L-1$ and represents a candidate sequence $\mathbf{x} = (X_0, X_1, \dots, X_{L-1})^T$, where X_k corresponds to the depth- k node of this path.

Let $\mathbf{x}_0^k \triangleq (X_0, \dots, X_k)$ be a partial path of length $k+1$, and define the node value v_k of a node of depth k associated with this partial path as

$$v_k \triangleq \min_{X_{k+1}, \dots, X_{L-1}} f(X_0, \dots, X_k, X_{k+1}, \dots, X_{L-1})$$

$$\text{with } \begin{cases} X_0, X_1, \dots, X_k & \text{known and fixed} \\ X_{k+1}, \dots, X_{L-1} & \in \mathcal{C} \end{cases} \quad (18)$$

where \mathcal{C} denotes the set of complex numbers. Note that X_{k+1}, \dots, X_{L-1} are “free” variables whose values are allowed to be any complex numbers without constraining them to be in \mathcal{S} . For every descendent node of the last node of \mathbf{x}_0^k , we compute the corresponding node value according to

$$v_{k+1} \triangleq \min_{X_{k+2}, \dots, X_{L-1}} f(X_0, \dots, X_{k+1}, X_{k+2}, \dots, X_{L-1})$$

$$\text{with } \begin{cases} X_0, X_1, \dots, X_{k+1} & \text{known and fixed} \\ X_{k+2}, \dots, X_{L-1} & \in \mathcal{C}. \end{cases} \quad (19)$$

Comparing (18) and (19), we see that $v_k \leq v_{k+1}$ as (18) has less constraint (or a larger degree of freedom). Hence, the node values of a path satisfy

$$v_0 \leq v_1 \leq v_2 \leq \dots \leq v_{L-1}. \quad (20)$$

Equation (18) implies that the node value v_{L-1} of a node at one of the tips of a tree is equal to $f(\mathbf{x})$ where $\mathbf{x} = (X_0, X_1, \dots, X_{L-1})^T$ is the corresponding candidate sequence. The problem of finding the minimum value of $f(\mathbf{x})$ is equivalent to finding the minimum node value amongst all tip nodes of the corresponding tree. As our tree structure implies that there is only one path from the root to a tip node, our problem becomes that of searching for the path whose last node has the minimum node value.

Before the search, we can assign a “reference solution” \mathbf{x}_o , which can be obtained either by the approximation solution of Section II-B or by other proper methods. In the process of searching for the “correct” path of a tree, if a node whose node value v is larger than the reference metric $f(\mathbf{x}_o)$, then (20) implies that for those candidate sequences \mathbf{x} 's whose corresponding paths pass through this node, $f(\mathbf{x}_o) < f(\mathbf{x})$. Thus, there is no need to search the subtree stemming from this node. On the other hand, when we visit a node of depth $L-1$ at the

tip of the tree and find $f(\mathbf{x}_1) < f(\mathbf{x}_o)$, where \mathbf{x}_1 is the corresponding candidate sequence, then we replace $f(\mathbf{x}_o)$ by $f(\mathbf{x}_1)$ and use it as the new reference metric in subsequent visits. At the end of the search process, the surviving \mathbf{x}_o will be the desired solution. The above procedure avoids searching all paths and makes the complexity of finding the solution computationally feasible.

B. Branch-and-Bound Algorithm

For our blind detection problem, the associated multivariate objective function is given by (13)

$$f(\mathbf{x}) = \|\mathbf{z} - \mathbf{V}\mathbf{z}\|^2 \quad (21)$$

where \mathbf{z} is a function of \mathbf{x} . The initial ‘‘reference solution’’ \mathbf{x}_0 , as just mentioned, can be obtained by the approximation method of Section II-B.

Now consider the calculation of the node value v_k of a node of depth k . When such a node is visited, the partial path \mathbf{x}_0^k has been determined. The definition of v_k , (18), requires that we find $(X_{k+1}, \dots, X_{L-1}) \in \mathcal{C}^{L-k-1}$ that minimizes $\|\mathbf{z} - \mathbf{V}\mathbf{z}\|^2$. Let $\mathbf{z}_a = (Z_0, Z_1, \dots, Z_k)^\top$, $\mathbf{z}_b = (Z_{k+1}, \dots, Z_{L-1})^\top$, and

$$\mathbf{V} = \begin{bmatrix} \mathbf{V}_a & \mathbf{V}_c \\ \mathbf{V}_c^\top & \mathbf{V}_b \end{bmatrix} \quad (22)$$

where \mathbf{V}_a is a $(k+1) \times (k+1)$ matrix, \mathbf{V}_b an $(L-k-1) \times (L-k-1)$ matrix, and \mathbf{V}_c a $(k+1) \times (L-k-1)$ matrix. Then

$$\begin{aligned} \|\mathbf{z} - \mathbf{V}\mathbf{z}\|^2 &= \mathbf{z}^\text{H} (\mathbf{I} - \mathbf{V}) \mathbf{z} \\ &= \left\| \begin{bmatrix} \mathbf{z}_a \\ \mathbf{z}_b \end{bmatrix} - \begin{bmatrix} \mathbf{V}_a & \mathbf{V}_c \\ \mathbf{V}_c^\top & \mathbf{V}_b \end{bmatrix} \begin{bmatrix} \mathbf{z}_a \\ \mathbf{z}_b \end{bmatrix} \right\|^2 \\ &= \mathbf{z}_a^\text{H} (\mathbf{I}_a - \mathbf{V}_a) \mathbf{z}_a - \mathbf{z}_a^\text{H} \mathbf{V}_c \mathbf{z}_b - \mathbf{z}_b^\text{H} \mathbf{V}_c^\top \mathbf{z}_a \\ &\quad + \mathbf{z}_b^\text{H} (\mathbf{I}_b - \mathbf{V}_b) \mathbf{z}_b \end{aligned} \quad (23)$$

where we have used the property $\mathbf{V}^2 = \mathbf{V}$ [see (12)]. While \mathbf{z}_a is determined by \mathbf{x}_0^k , \mathbf{z}_b is allowed to be any complex vector of length $(L-1-k)$, as there is no constraint on the associated data vector. Taking derivative with respect to \mathbf{z}_b , we see that the objective function (23) is minimized when

$$\mathbf{z}_b = \mathbf{A}\mathbf{z}_a \quad (24)$$

and the corresponding minimum value is

$$v_k = \mathbf{z}_a^\text{H} \mathbf{B}\mathbf{z}_a \quad (25)$$

where the constant matrices \mathbf{A} and \mathbf{B} are given by

$$\mathbf{A} = (\mathbf{I}_b - \mathbf{V}_b)^{-1} \mathbf{V}_c^\top \quad (26)$$

$$\mathbf{B} = (\mathbf{I}_a - \mathbf{V}_a) - \mathbf{V}_c (\mathbf{I}_b - \mathbf{V}_b)^{-1} \mathbf{V}_c^\top \quad (27)$$

and \mathbf{I}_a and \mathbf{I}_b are identity matrices of the same dimensions as those of \mathbf{V}_a and \mathbf{V}_b , respectively. Note that when (24) is satisfied, we have

$$\begin{aligned} &\begin{bmatrix} \mathbf{I}_a - \mathbf{V}_a & -\mathbf{V}_c \\ -\mathbf{V}_c^\top & \mathbf{I}_b - \mathbf{V}_b \end{bmatrix} \begin{bmatrix} \mathbf{z}_a \\ \mathbf{z}_b \end{bmatrix} \\ &= \begin{bmatrix} (\mathbf{I}_a - \mathbf{V}_a) \mathbf{z}_a - \mathbf{V}_c \mathbf{z}_b \\ -\mathbf{V}_c^\top \mathbf{z}_a + (\mathbf{I}_b - \mathbf{V}_b) \mathbf{z}_b \end{bmatrix} \\ &= \begin{bmatrix} (\mathbf{I}_a - \mathbf{V}_a) \mathbf{z}_a - \mathbf{V}_c \mathbf{A}\mathbf{z}_a \\ -\mathbf{V}_c^\top \mathbf{z}_a + (\mathbf{I}_b - \mathbf{V}_b) \mathbf{A}\mathbf{z}_a \end{bmatrix} \\ &= \begin{bmatrix} (\mathbf{I}_a - \mathbf{V}_a) \mathbf{z}_a - \mathbf{V}_c (\mathbf{I}_b - \mathbf{V}_b)^{-1} \mathbf{V}_c^\top \mathbf{z}_a \\ \mathbf{0} \end{bmatrix} \end{aligned} \quad (28)$$

and then

$$\begin{aligned} v_k &= \begin{bmatrix} \mathbf{z}_a^\text{H} & \mathbf{z}_b^\text{H} \end{bmatrix} \begin{bmatrix} \mathbf{I}_a - \mathbf{V}_a & -\mathbf{V}_c \\ -\mathbf{V}_c^\top & \mathbf{I}_b - \mathbf{V}_b \end{bmatrix} \begin{bmatrix} \mathbf{z}_a \\ \mathbf{z}_b \end{bmatrix} \\ &= \mathbf{z}_a^\text{H} [(\mathbf{I}_a - \mathbf{V}_a) - \mathbf{V}_c (\mathbf{I}_b - \mathbf{V}_b)^{-1} \mathbf{V}_c^\top] \mathbf{z}_a \end{aligned} \quad (29)$$

which is (25).

C. A Recursive Formula

To reduce the calculation of node values, we develop a recursive relation between v_k and v_{k-1} , the node value of a node of depth $k-1$ and that of one of its immediate descendent nodes of depth k . For convenience and preciseness, we add a superscript k to each vector and matrix appearing in (22)–(27) to signify that these quantities are related to nodes of depth k . From (25) and (27), we obtain

$$\begin{aligned} v_{k-1} &= \mathbf{z}_a^{(k-1)\text{H}} \left[(\mathbf{I}_a^{(k-1)} - \mathbf{V}_a^{(k-1)}) \right. \\ &\quad \left. - \mathbf{V}_c^{(k-1)} (\mathbf{I}_b^{(k-1)} - \mathbf{V}_b^{(k-1)})^{-1} \mathbf{V}_c^{(k-1)\top} \right] \mathbf{z}_a^{(k-1)} \end{aligned} \quad (30)$$

$$\begin{aligned} v_k &= \mathbf{z}_a^{(k)\text{H}} \left[(\mathbf{I}_a^{(k)} - \mathbf{V}_a^{(k)}) \right. \\ &\quad \left. - \mathbf{V}_c^{(k)} (\mathbf{I}_b^{(k)} - \mathbf{V}_b^{(k)})^{-1} \mathbf{V}_c^{(k)\top} \right] \mathbf{z}_a^{(k)} \end{aligned} \quad (31)$$

where $\mathbf{z}_a^{(k-1)} = (Z_0, \dots, Z_{k-1})^\top$, and $\mathbf{z}_a^{(k)} = (Z_0, \dots, Z_{k-1}, Z_k)^\top$. Equation (24) reveals that the corresponding unconstrained sequence $\mathbf{z}_b^{(k-1)}$ for $\mathbf{z}_a^{(k-1)}$ is

$$\mathbf{z}_b^{(k-1)} = \mathbf{A}^{(k-1)} \mathbf{z}_a^{(k-1)} \triangleq (z_k, \dots, z_{L-1})^\top. \quad (32)$$

It can be shown (see Appendix) that

$$v_k = v_{k-1} + \mu_k |\Delta z_k|^2 \quad (33)$$

where $\Delta z_k = (Z_k - z_k)$, μ_k is the $(k+1, k+1)$ th element of the matrix $\mathbf{B}^{(k)}$, and z_k is derived from (32) via

$$z_k = \mathbf{p}^{(k-1)\top} \mathbf{z}_a^{(k-1)} \quad (34)$$

$\mathbf{p}^{(k-1)\top}$ being the first row of $\mathbf{A}^{(k-1)}$. Compared with (25), (33) clearly reduces much of the complexity of computing v_k .

D. Summary of the Algorithm

We summarize the proposed algorithm in the following. The constant parameters $\mathbf{p}^{(k)\top}$'s and μ_k 's are independent of the received vector. They can be predetermined and prestored before the blind detector is activated. We refer to this portion of the algorithm as the *precalculation* stage.

Precalculation:

- 1) Let the block length be L and the regression polynomial degree be d . Compute and store $\mathbf{q}_n = (n^d, \dots, n, 1)^\top$, $n = 0, 1, \dots, L-1$, \mathbf{Q} , and \mathbf{V} .
- 2) For $k = 1, 2, \dots, L-1$, compute and store $\mathbf{V}_a^{(k)}$, $\mathbf{V}_b^{(k)}$, and $\mathbf{V}_c^{(k)}$ according to (22), and $\mathbf{A}^{(k)}$ and $\mathbf{B}^{(k)}$ based on (26) and (27).
- 3) Determine $\mathbf{p}^{(k)\top}$ and μ_k , $k = 1, 2, \dots, L-1$. By definition, $\mathbf{p}^{(k)\top}$ is the first row of $\mathbf{A}^{(k)}$, and μ_k is the $(k+1, k+1)$ th entry of $\mathbf{B}^{(k)}$.

Tree search:

- 1) Determine an initial estimate \mathbf{x}_o by the approximation method of Section II-B, and compute $f(\mathbf{x}_o) = \|\mathbf{z}_o - \mathbf{V}\mathbf{z}_o\|^2 \triangleq v_o$, where $\mathbf{z}_o = \mathbf{z}(\mathbf{x}_o)$.
- 2) Starting at the root node, calculate the node values of all its immediate descendent nodes. Visit the immediate descendent node with the smallest node value.
- 3) When visiting a node, calculate the node values of all its immediate descendent nodes. Go visiting the immediate descendent node with the smallest node value, if this node value is smaller than v_o . Otherwise, discard all these immediate descendent nodes and the corresponding subtrees (i.e., all their descendent nodes).
- 4) When all descendent nodes of a node are either visited or discarded, go back to the immediate ancestor node. Among the unvisited immediate descendent nodes of this ancestor node, visit the one with the smallest node value, if this node value is smaller than v_o .
- 5) On visiting a node of depth $L - 1$ at the tip, if $v_{L-1} < v_o$, then replace v_o by v_{L-1} and \mathbf{x}_o by the corresponding sequence. Go back to the immediate ancestor node and follow the procedure of step 4.
- 6) The search process ends with \mathbf{x}_o as the desired solution when there are no more unvisited nodes.

Several remarks on the above algorithm are in order.

- R1. The stack algorithm used in decoding convolutional codes can also be used as an alternative search method, replacing the path metric by the node value.
- R2. Our simulation shows that, on the average, the number of nodes visited by the stack algorithm is about the same as that visited by the proposed algorithm. However, our algorithm does not have to perform stack sorting, and furthermore, for the stack algorithm, the number of stored nodes is unpredictable and may become very large at low SNRs, though the inequality (20) can be used to reduce the stack size. On the other hand, the number of stored nodes of the above algorithm is upper bounded by $L|\mathcal{S}|$.
- R3. The Fano metric has been proved to be optimal in additive white Gaussian noise (AWGN) channels in the sense it gives the best one-step prediction [14]. The stack algorithm, when using the Fano metric, usually renders very fast searching speed. The node value, v_k , plays a similar role to the Fano metric in AWGN, providing a (model-constrained) optimal one-step predictor in a channel with memory.
- R4. A good initial estimate (upper bound) \mathbf{x}_o is crucial for reducing the search time.
- R5. For a convolutional-coded OFDM system, blind decoding can be carried out in a similar manner with reduced complexity, by merging subtrees originating from the same state. But the corresponding code tree cannot be permuted, as suggested in Section IV-B, to accelerate the search speed.

IV. RELATED DESIGN ISSUES

We now address some related design issues and modifications, some of which lead to much less complex algorithms for both blind and semiblind (low pilot density) detections.

A. Ambiguity Problem

When the estimated data sequence for the m th subchannel $\mathbf{x}_m = (X_{m0}, X_{m1}, \dots, X_{m(L-1)})^T$ is obtained, we still have to resolve the phase ambiguity associated with any blind equalizer for signals with symmetric constellation. For example, if \mathbf{x}_o is a solution of (13), and the square 16-quadrature amplitude modulation (QAM) is used, then $-\mathbf{x}_o$ and $\pm i\mathbf{x}_o$ are also solutions of (13). We can avoid the ambiguity issue by employing differential modulation, e.g., differential M -ary phase-shift keying (MPSK) or differential 16 STAR-QAM [13]. Coding (e.g., rotational-invariant trellis codes) is another alternative to eliminate the ambiguity. The simplest solution is obtained by using a few pilot symbols. As pilots are used only for resolving the phase ambiguity, they can be much sparser than those used in pilot-aided equalizers. Let the inter-pilot distance be $D_p \cdot L$ symbols, where L is the block length. If $D_p > 1$, blocks without pilots can derive the phase information from the preceding pilot-inserted blocks.

Adjacent subchannels can also be used jointly to determine the phase rotation. Letting the estimated data sequence of the m th subchannel be \mathbf{x}_m , we should choose $\{\theta_m\}$ such that $\{e^{i\theta_m} \mathbf{x}_m\}_{m=0}^{M_0-1}$ best matches (4). More specifically, we select $\boldsymbol{\theta} = (\theta_0, \theta_1, \dots, \theta_{(M_0-1)})^T$ such that

$$\min_{\boldsymbol{\theta}} \sum_{m=0}^{M_0-1} \sum_{n=0}^{L-1} \left| \frac{Y_{mn}}{e^{i\theta_m} X_{mn}} - F(m, n, \mathbf{c}) \right|^2 \quad (35)$$

where θ_m is one of the candidate ambiguity angles.

B. Influence of the Search Order and Its Implications

The sequence shown on top of Fig. 1 indicates a tree search in natural order, i.e., from X_0 to X_{L-1} . Let $(\alpha_0, \alpha_1, \dots, \alpha_{L-1})$ be a permutation of $(0, 1, \dots, L-1)$. The corresponding rearranged sequence $\mathbf{x}' = (X_{\alpha_0}, X_{\alpha_1}, \dots, X_{\alpha_{L-1}})^T$ is now associated with a tree whose depth- k nodes are related to the possible values of X_{α_k} .

When performing the search process in this new tree, there is essentially no change of algorithmic complexity and logical procedure. We only have to replace \mathbf{V} by \mathbf{V}' , where $\mathbf{V}' = [V'_{ij}]$ is related to $\mathbf{V} = [V_{ij}]$ by

$$V'_{jk} = V_{\alpha_j \alpha_k}, \quad 0 \leq j, k, \alpha_j, \alpha_k \leq L-1 \quad (36)$$

and modify the corresponding parameters $\mathbf{p}^{(k)}$'s and μ_k 's accordingly. As our goal is to find a data sequence $\mathbf{x} = (X_0, X_1, \dots, X_{L-1})^T$ that minimizes (13), the order of finding the component values of the sequence is immaterial. Moreover, the final solution obtained by our algorithm is independent of the order of search. However, the search order does have some critical impacts.

Recall that the orthogonal projection \mathbf{V} in the objective function (21) is used to find the regression curve that fits the channel estimate \mathbf{z} in the LS sense (which is $\mathbf{V}\mathbf{z}$). On the one hand, we are looking for a smooth curve to fit \mathbf{z} ; on the other hand, the components of \mathbf{z} tend to dictate and bound the ‘‘shape’’ of the candidate LS regression curve. If the final curve is smooth enough, the projection of the subsequence that consists of, say, two end points and the middle point alone,

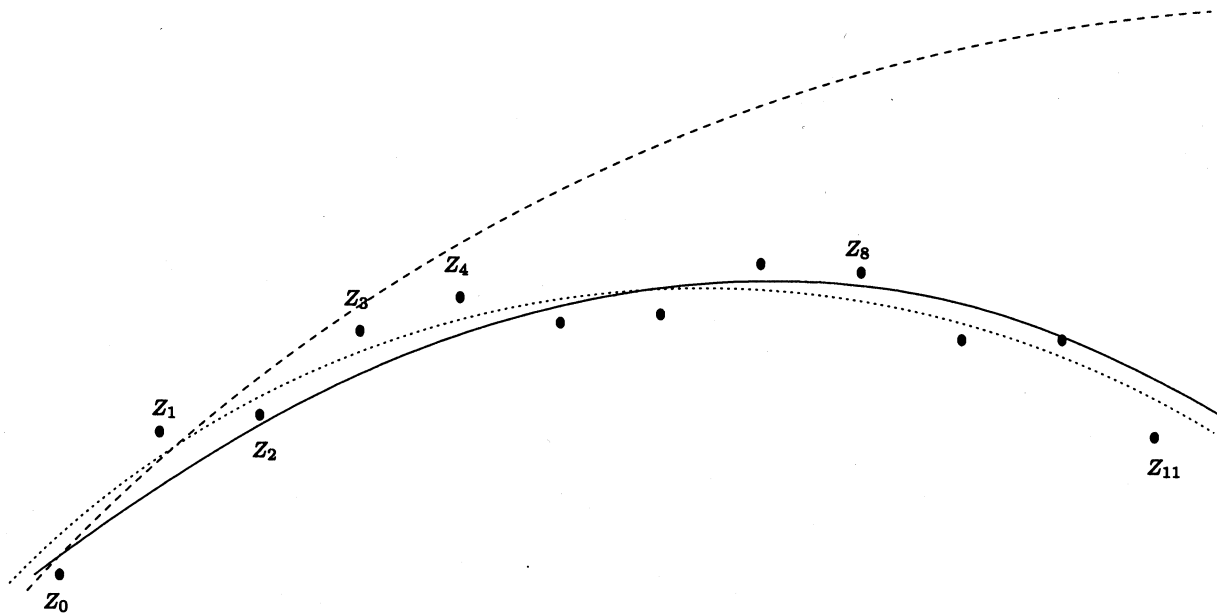


Fig. 2. Regression curve resulting from the natural search order (Z_0, Z_1, Z_2, Z_3) (the dash curve), and that due to a preferred uniformly contracting Z search order (Z_0, Z_4, Z_8, Z_{11}) (the dot curve). The actual values are indicated by the solid curve.

i.e., $\mathbf{z}^{(0)} = (Z_0, Z_{\lfloor L/2 \rfloor}, Z_{L-1})^T$, where $\lfloor x \rfloor$ = the integer part of x , will give a crude but probably acceptable approximation. Using two additional middle points to form the new subsequence $\mathbf{z}^{(1)} = (Z_0, Z_{\lfloor L/4 \rfloor}, Z_{\lfloor L/2 \rfloor}, Z_{\lfloor 3L/4 \rfloor}, Z_{L-1})^T$, we obtain a better approximation by the projection $\mathbf{V}\mathbf{z}^{(1)}$. If, instead, we use the first three points $(Z_0, Z_1, Z_2)^T$ or the first five points, $(Z_0, \dots, Z_4)^T$ for LS fitting, we are likely to obtain curves that are far from the final solution; see Fig. 2 for an example.

The above geometric property of LS-fitted solutions has two important implications:

- 1) the convergence speed of our tree-search algorithm, i.e., the number of nodes visited before the “distance” between the corresponding LS-fitted curve and the final solution is within an acceptable bound, does depend on the order of our search;
- 2) in most cases, a good approximation of the final LS-fitted curve can be obtained by using only a few properly selected components of \mathbf{z} .

1) suggests that a preferred search order is the contracting zigzag pattern $(0, L, \lfloor L/2 \rfloor, \lfloor L/4 \rfloor, \lfloor 3L/4 \rfloor, \dots)$, where the integers denote the data locations. For example, if the sequence length $L = 10$, a preferred order is $(0, 9, 5, 2, 7, 1, 3, 6, 4, 8)$. On the other hand, 2) indicates that we can reduce the search complexity by estimating the CRs based on the received samples at some selected locations to begin with, then use the estimated CRs for detecting data at the remaining locations. Often, we achieve significant complexity reduction with negligible performance degradation by using the preferred search order and selective detection. For example, if $L = 12$, we can use, say Z_0, Z_4, Z_8, Z_{11} to estimate X_0, X_4, X_8, X_{11} and H_0, H_4, H_8, H_{11} via the proposed tree-searching algorithm. Fig. 2 indicates that when this much-shortened search is

finished, already the implicit regression curve comes very close to the desired one. Given the CR estimates $\hat{H}_i, i = 0, 1, \dots, 11$ derived from the curve, we then make hard decisions on $Y_i/\hat{H}_i, i \notin \{0, 4, 8, 11\}$ to determine the remaining X_i 's.

C. Partially Known Sequences and Semiblind Detection

In some situations, a few elements of the sequence $\mathbf{x} = (X_0, X_1, \dots, X_{L-1})^T$ in (13) are known or have been detected. These known symbols may be pilots, or members of the overlapped part of the current and the previous blocks. Denote the detected or known symbols in \mathbf{x} by $(X_{\alpha_0}, X_{\alpha_1}, \dots, X_{\alpha_{l_1-1}})^T$. We can arrange the search order so that it begins with these known symbols, followed by the undetected symbols $(X_{\alpha_{l_1}}, \dots, X_{\alpha_{L-1}})^T$. Since the first l_1 symbols are known, the tree degenerates to one with only one branch for the beginning l_1 nodes.

Semiblind schemes refer to those that use sparse pilot symbols. Besides resolving the phase ambiguity, these few pilot symbols can, as mentioned before, be used to reduce the tree size and provide an initial reference solution and reference metric. We can apply any conventional pilot-assisted method to obtain the reference sequence and metric.

Let us consider a semiblind scheme where two pilots are located at both ends of a block. We apply the blind algorithm with the preferred search order that starts at the pilot positions 0 and $L-1$. An initial reference solution \mathbf{z}_0 is obtained by using the two pilots only. To detect the remaining data, we conduct the tree search to level L_0 , where we can choose $L_0 \ll L-1$ to reduce the complexity. This scheme will be referred to as the low-complexity semiblind (LCSB) detector.

Note that the pilot-assisted scheme proposed in [5] corresponds to the case when the search order is such that the resulting tree is partially degenerated with the search depth equal to the number of pilots.

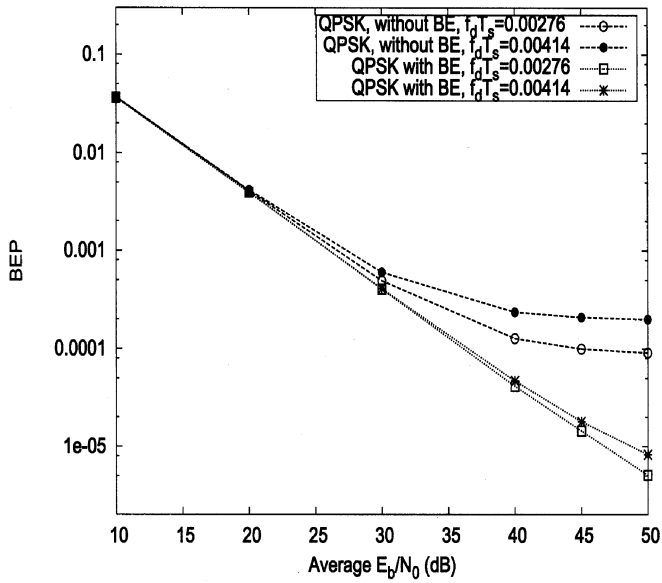


Fig. 3. BEP performance of the differential QPSK-OFDM system with or without blind channel estimation.

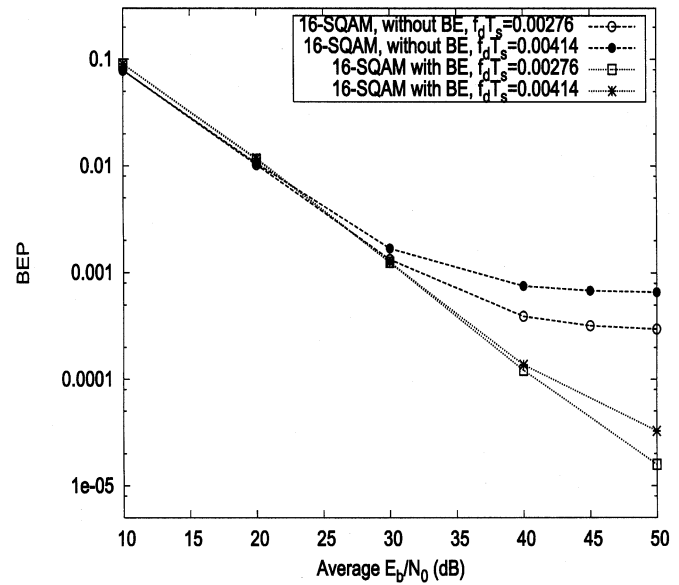


Fig. 4. BEP performance of the differential 16 STAR-QAM-OFDM system with or without blind channel estimation.

V. NUMERICAL RESULTS AND DISCUSSIONS

Numerical results reported in this section assume an OFDM system that has a total bandwidth of 1 MHz, 32 equally spaced mutual-orthogonal subchannels, and a symbol time (included guard interval) of 34.5 μ s. A four-ray Jakes' model [15] is used to simulate a time-varying Rayleigh fading channel with relative power strengths 0.5938, 0.7305, 0.3175, 0.1137, and delays 0, 0.1, 0.5, 1 μ s, respectively. As mentioned before, the algorithms presented in the previous two sections are 1-D solutions. Except for Section V-C, we assume that a 1-D algorithm is carried out on each subchannel for detecting OFDM signals.

A. Differential Encoding Schemes

The first case we consider is the differential scheme that employs a quaternary phase-shift keying (QPSK) or 16 STAR-QAM [13] signal. The (time) block length used is $L = 12$ symbols. Figs. 3 and 4 plot the bit-error probability (BEP) performance curves of differential QPSK (DQPSK) and 16 STAR-QAM systems using the algorithm presented in Section III-D. The performance of the uncompensated systems, i.e., those labeled by "without BE" (blind equalization), is also given in those two figures. It can be seen that the proposed blind detection scheme does enhance the performance of both systems and yields no obvious error floor.

B. Semiblind Scheme (I)

For the semiblind scheme, very few pilots are available, and they are used only for eliminating the phase ambiguity and providing the initial channel estimation. In each subchannel, we insert one pilot symbol for every D_p blocks with a block length of $L = 12$, or equivalently, every $L \cdot D_p$ symbols. The pilot density is much sparser than that of the conventional pilot-assisted schemes for time-varying channels [5]. Fig. 5 plots the BEP performance for the square 16QAM-OFDM system with $D_p = 1, 2, 3$ and normalized Doppler shift 0.00414. We notice a small increase in BEP occurs when D_p increases, especially

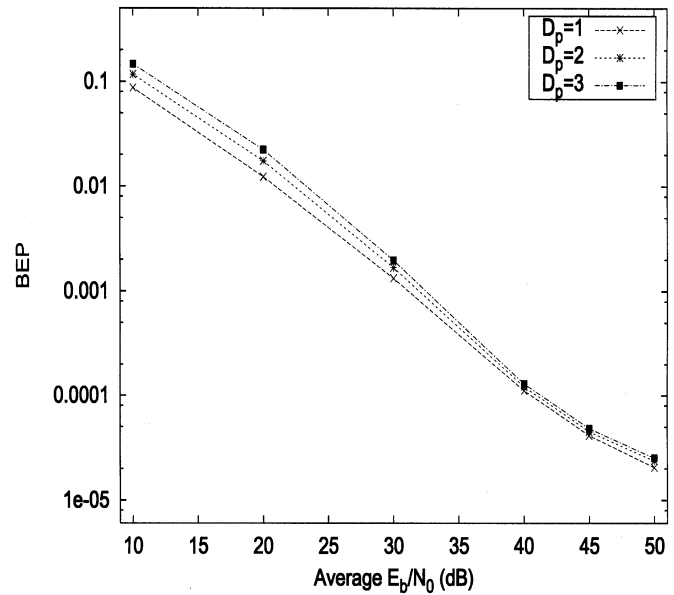


Fig. 5. BEPs of the 16QAM-OFDM system using the semiblind detector with full searching depth; $f_d T_s = .00414$.

at high SNRs. There are no clear error floors within the range of interest, though.

C. Semiblind Scheme (II) (Low-Complexity Scheme)

The performance of the LCSB algorithm described in Section IV-C is given in Fig. 6, with $L = 21$ and search depth $L_d = 3$ ($L_d = 4$ yields almost identical performance). The two pilot symbols within a block are at positions 0 and 20, respectively. We adopt the contracting- Z search pattern that begins with the pilot positions 0 and 20. For $L_d = 3$, the search order is 0, 20, 10, and then 5. The normalized Doppler shift is $f_d T_s = 0.01$, which corresponds to a fast time-varying channel. We also plot the BEP curves of the ideal receiver that has perfect channel estimate and the conventional pilot-assisted scheme. The performance of the ideal receiver provides the absolute

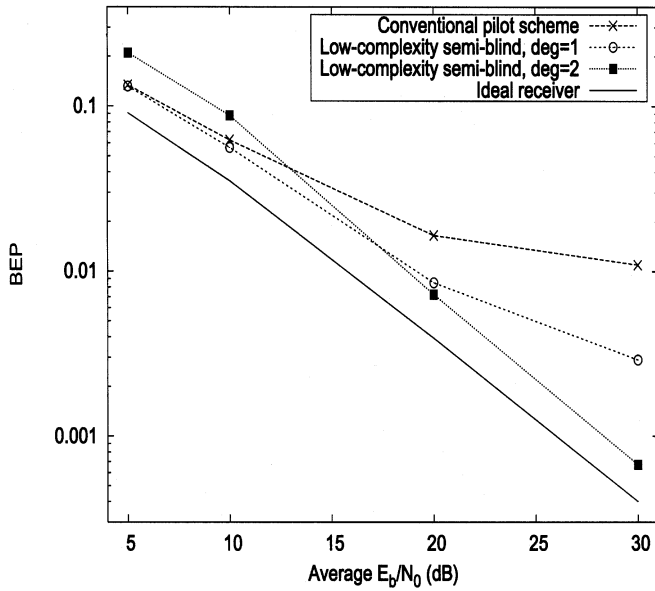


Fig. 6. BEPs of the 16QAM-OFDM system using the LCSB detector with a single pilot at both ends of a block; $f_d T_s = 0.01$. The search depth is $L_d = 3$ for both degree 1 and 2 models.

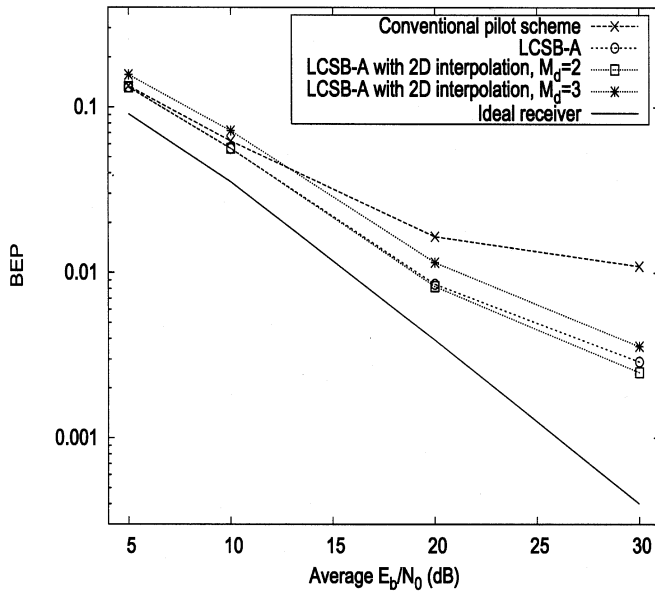


Fig. 7. BEPs of the 16QAM-OFDM system using the LCSB two-stage algorithm with a 2-D regression model. The system and channel parameters are the same as Fig. 6. The curve labeled by LCSB-A is the same as the the curve labeled by “Low complexity semiblind, deg = 1” in Fig. 6.

lower bound for comparison. Since the channel is fast fading and there are only two pilots in a block, the pilot-assisted scheme cannot give a reliable channel estimate.

The LCSB estimate has more impressive relative performance gain at high SNRs. For SNR less than 10 dB, however, using a degree-2 polynomial model may degrade the performance, though it greatly improves the performance at SNR higher than 30 dB. In most cases, the performance is only a few decibels away from that of the ideal receiver. Such a behavior has been expected, as a higher order model involves more parameters, and the channel estimation error which dictates

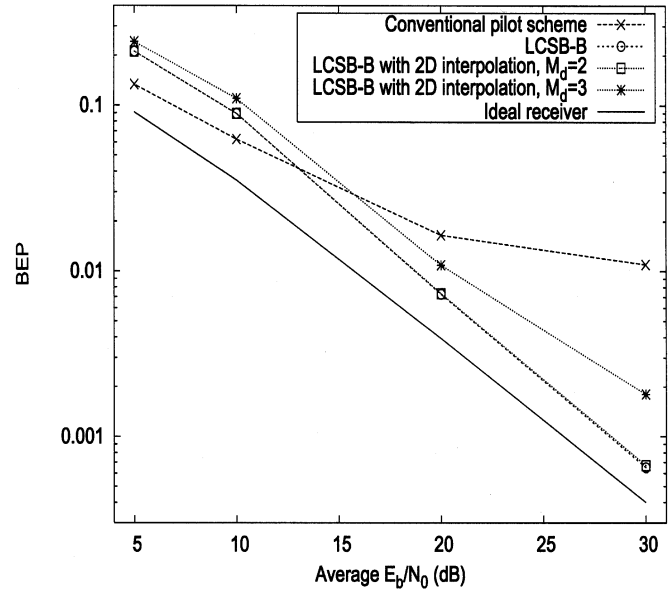


Fig. 8. BEPs of the 16QAM-OFDM system using the LCSB two-stage algorithm with a 2-D regression model. The system and channel parameters are the same as Fig. 6. The curve labeled by LCSB-B is the same as the the curve labeled by “Low complexity semiblind, deg = 2” in Fig. 6.

the BEP performance is mainly caused by thermal noise at low SNRs, while the modeling error dominates at high SNRs.

The performance of the two-stage algorithm described in Section II-A is plotted in Fig. 7. The semiblind algorithm is employed on every M_d subchannels, then the detected symbols are treated as pilots, so that the 2-D model-based algorithm of [5] can be used to detect the remaining symbols.

For comparison, Fig. 7 presents the performance of the detector that uses the LCSB scheme with degree-1 polynomial model (labeled by LCSB-A) on all subchannels under the same system and channel conditions. As the 2-D model takes into account the correlations among subchannels, the two-stage scheme outperforms LCSB-A for the case $M_d = 2$. For $M_d = 3$, the correlation between two detected subchannels becomes less obvious, and the resulting modeling error degrades the performance of the two-stage algorithm. Fig. 8 provides a similar comparison for the LCSB scheme with a degree-2 polynomial model (labeled by LCSB-B). Similar to the 1-D case shown in Fig. 8, a higher degree model yields smaller modeling error, and thus better BEP performance at high SNRs.

D. Complexity

Fig. 9 shows the average number of visited nodes for detecting QPSK and 16QAM signals with different search depths (L_d). The average E_b/N_0 or SNR per block in this figure denotes the “locally” average SNR within a block of OFDM symbols as our algorithms operate in a block-by-block fashion. Note that even if the average E_b/N_0 is 10 dB, the received signal may suffer from severe fading in a few consecutive blocks, where the average E_b/N_0 per block is far less than 0 dB. It is more appropriate to present the complexity performance in terms of average E_b/N_0 per block.

We see that for larger average SNRs, the visited number is very small for all cases, while if the average SNR is smaller than a threshold, which depends on the signal constellation size

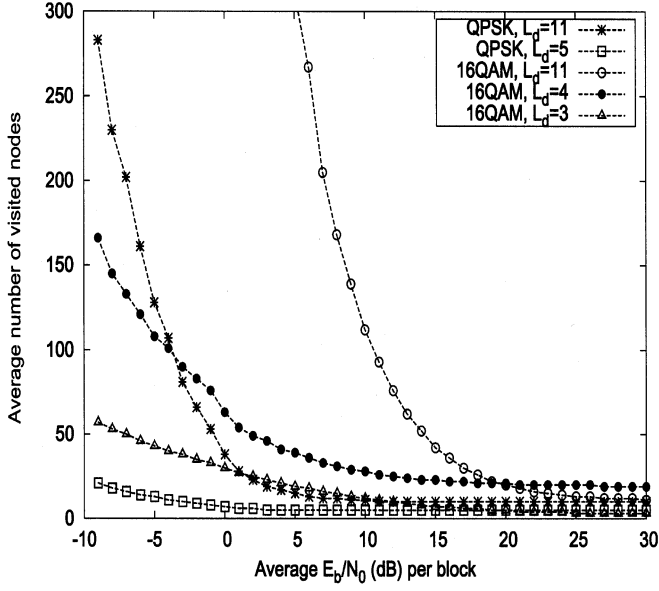


Fig. 9. Average number of visited nodes within a block for detecting QPSK-OFDM and 16QAM-OFDM signals with different search depths L_d .

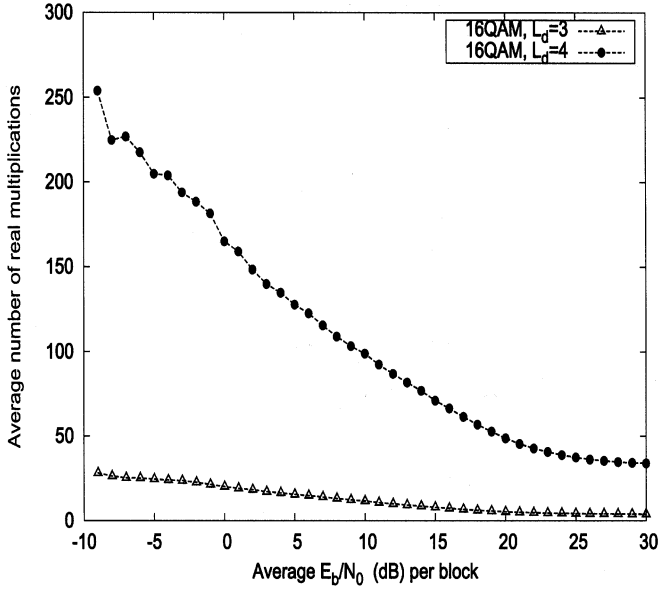


Fig. 10. Average number of real multiplications per symbol in the tree-search process for the LCSB scheme with $L = 21$.

and the search depth, the visited node numbers increase rapidly. On the other hand, adopting the partial search algorithm ($L_d < L - 1$) greatly reduces the number of visited nodes, especially at low SNRs. The curves $L_d = 3$ and $L_d = 4$ (with the search order 0, 20, 10, 5, and 15) correspond to the LCSB schemes of Fig. 6. We also show in Fig. 10 the average number of real multiplications per symbol in a block of length $L_d = 21$, which corresponds to the curves $L_d = 3$ and $L_d = 4$ in Fig. 9. We see that for the case $L_d = 3$, we need about 30 real multiplications for each symbol per subchannel when the average SNR within a block is -9 dB.

Using the complete blind algorithm of Section III-D for every subchannel is expected to require high computing complexity at low and medium SNRs. Therefore, we have proposed various reduced-complexity derivative algorithms, and have just shown that, if a preferred search order together with a reduced

search-depth scheme are used, the resulting complexity is only moderate. The availability of few pilots (semiblind) helps in reducing the complexity. The numbers in the above two figures are obtained on a per-subchannel basis, hence, the two-stage algorithm will bring further complexity reduction. However, the complexity of our algorithms is lower bounded by that of [5], as both use the same model-based approach.

VI. CONCLUSION

In this paper, we present new joint detection and estimation methods for blindly and semiblindly detecting OFDM signals in a fading environment, where the CR may vary from symbol to symbol. Based on a regression model, we formulate the blind detection problem as an integer programming problem. We then apply an elementary inequality and the branch-and-bound principle to find the optimal solution. A recursive formula is derived to simplify the computation effort. We propose low-complexity variations that use the concepts of preferred search order and selective detection. A two-stage algorithm is suggested to reduce the overall detection complexity, and numerical examples are given to examine the performance and complexity of the proposed blind and semiblind detection schemes. Our algorithms are shown to be effective in combating fast channel variations and in providing satisfactory system performance.

APPENDIX

DERIVATION OF A RECURSIVE FORMULA

Let v_{k-1} and v_k be the node values of a node of depth $k - 1$ and one of its immediate descendent nodes of depth k , respectively. From (25) and (27), we have

$$v_{k-1} = \mathbf{z}_a^{(k-1)H} \left[\left(\mathbf{I}_a^{(k-1)} - \mathbf{V}_a^{(k-1)} \right) - \mathbf{V}_c^{(k-1)} \left(\mathbf{I}_b^{(k-1)} - \mathbf{V}_b^{(k-1)} \right)^{-1} \mathbf{V}_c^{(k-1)T} \right] \mathbf{z}_a^{(k-1)} \quad (\text{A.1})$$

$$v_k = \mathbf{z}_a^{(k)H} \left[\left(\mathbf{I}_a^{(k)} - \mathbf{V}_a^{(k)} \right) - \mathbf{V}_c^{(k)} \left(\mathbf{I}_b^{(k)} - \mathbf{V}_b^{(k)} \right)^{-1} \mathbf{V}_c^{(k)T} \right] \mathbf{z}_a^{(k)} \quad (\text{A.2})$$

where $\mathbf{z}_a^{(k-1)} = (Z_0, \dots, Z_{k-1})^T$, and $\mathbf{z}_a^{(k)} = (Z_0, \dots, Z_{k-1}, Z_k)^T$. The corresponding unconstrained sequence for $\mathbf{z}_a^{(k-1)}$ is

$$\mathbf{z}_b^{(k-1)} = \mathbf{A}^{(k-1)} \mathbf{z}_a^{(k-1)} \triangleq (z_k, z_{k+1}, \dots, z_{L-1})^T. \quad (\text{A.3})$$

Let

$$\begin{aligned} \mathbf{z}'_a^{(k)} &= (Z_0, \dots, Z_{k-1}, z_k)^T \\ \mathbf{z}'_b^{(k)} &= (z_{k+1}, \dots, z_{L-1})^T. \end{aligned}$$

Since

$$\mathbf{V} = \begin{bmatrix} \mathbf{V}_a^{(k-1)} & \mathbf{V}_c^{(k-1)} \\ \mathbf{V}_c^{(k-1)T} & \mathbf{V}_b^{(k-1)} \end{bmatrix} = \begin{bmatrix} \mathbf{V}_a^{(k)} & \mathbf{V}_c^{(k)} \\ \mathbf{V}_c^{(k)T} & \mathbf{V}_b^{(k)} \end{bmatrix} \quad (\text{A.4})$$

and

$$\begin{aligned} [\mathbf{z}_a^{(k-1)T} \mathbf{z}_b^{(k-1)T}] &= [\mathbf{z}'_a^{(k)T} \mathbf{z}'_b^{(k)T}] \\ &= (Z_0, \dots, Z_{k-1}, z_k, z_{k+1}, \dots, z_{L-1}) \end{aligned}$$

we obtain

$$\begin{aligned} & \begin{bmatrix} (\mathbf{I}_a^{(k-1)} - \mathbf{V}_a^{(k-1)})\mathbf{z}_a^{(k-1)} - \mathbf{V}_c^{(k-1)}\mathbf{z}_b^{(k-1)} \\ -\mathbf{V}_c^{(k-1)\top}\mathbf{z}_a^{(k-1)} + (\mathbf{I}_b^{(k-1)} - \mathbf{V}_b^{(k-1)})\mathbf{z}_b^{(k-1)} \end{bmatrix} \\ &= \begin{bmatrix} (\mathbf{I}_a^{(k)} - \mathbf{V}_a^{(k)})\mathbf{z}'_a^{(k)} - \mathbf{V}_c^{(k)}\mathbf{z}'_b^{(k)} \\ -\mathbf{V}_c^{(k)\top}\mathbf{z}'_a^{(k)} + (\mathbf{I}_b^{(k)} - \mathbf{V}_b^{(k)})\mathbf{z}'_b^{(k)} \end{bmatrix}. \quad (\text{A.5}) \end{aligned}$$

Invoking (24) and (26), we obtain

$$\left(\mathbf{I}_b^{(k-1)} - \mathbf{V}_b^{(k-1)}\right)\mathbf{z}_b^{(k-1)} = \mathbf{V}_c^{(k-1)\top}\mathbf{z}_a^{(k-1)} \quad (\text{A.6})$$

i.e., the elements in the lower part of the vector in the left-hand side of (A.5), from the $(k+1)$ th to the L th elements, are all zero. Hence, the corresponding elements in the right-hand side of (A.5) are also zero, and we have

$$\left[\left(\mathbf{I}_a^{(k)} - \mathbf{V}_a^{(k)}\right)\mathbf{z}'_a^{(k)} - \mathbf{V}_c^{(k)}\mathbf{z}'_b^{(k)}\right]_{k+1} = 0 \quad (\text{A.7})$$

$$-\mathbf{V}_c^{(k)\top}\mathbf{z}'_a^{(k)} + \left(\mathbf{I}_b^{(k)} - \mathbf{V}_b^{(k)}\right)\mathbf{z}'_b^{(k)} = 0 \quad (\text{A.8})$$

where $[\mathbf{v}]_k$ denotes the k th elements of the vector \mathbf{v} . Multiplying the right-hand side of (A.5) by $\begin{bmatrix} \mathbf{z}'_a^{(k)\text{H}} & \mathbf{z}'_b^{(k)\text{H}} \end{bmatrix}$ and using (A.8), we obtain

$$\begin{aligned} v_{k-1} &= \mathbf{z}'_a^{(k)\text{H}} \left[\left(\mathbf{I}_a^{(k)} - \mathbf{V}_a^{(k)}\right) \right. \\ &\quad \left. - \mathbf{V}_c^{(k)} \left(\mathbf{I}_b^{(k)} - \mathbf{V}_b^{(k)}\right)^{-1} \mathbf{V}_c^{(k)\top} \right] \mathbf{z}'_a^{(k)}. \quad (\text{A.9}) \end{aligned}$$

Comparing (A.9) with (A.2) and noticing that $\mathbf{z}_a^{(k)} = \mathbf{z}'_a^{(k)} + (0, \dots, 0, \Delta z_k)^\top$, where $\Delta z_k = (Z_k - z_k)$, we obtain

$$\begin{aligned} v_k - v_{k-1} &= (\Delta z_k)^* \left[\left(\mathbf{I}_a^{(k)} - \mathbf{V}_a^{(k)}\right) \mathbf{z}'_a^{(k)} \right. \\ &\quad \left. - \mathbf{V}_c^{(k)} \left(\mathbf{I}_b^{(k)} - \mathbf{V}_b^{(k)}\right)^{-1} \mathbf{V}_c^{(k)\top} \mathbf{z}'_a^{(k)} \right]_{k+1} \\ &\quad + \left[\mathbf{z}'_a^{(k)\text{H}} \left(\mathbf{I}_a^{(k)} - \mathbf{V}_a^{(k)}\right) \right. \\ &\quad \left. - \mathbf{z}'_a^{(k)\text{H}} \mathbf{V}_c^{(k)} \left(\mathbf{I}_b^{(k)} - \mathbf{V}_b^{(k)}\right)^{-1} \mathbf{V}_c^{(k)\top} \right]_{k+1} \Delta z_k \\ &\quad + (\Delta z_k)^* \left[\left(\mathbf{I}_a^{(k)} - \mathbf{V}_a^{(k)}\right) \right. \\ &\quad \left. - \mathbf{V}_c^{(k)} \left(\mathbf{I}_b^{(k)} - \mathbf{V}_b^{(k)}\right)^{-1} \mathbf{V}_c^{(k)\top} \right]_{(k+1)(k+1)} \Delta z_k. \quad (\text{A.10}) \end{aligned}$$

Substituting (A.7) and (A.8) into (A.10) and invoking (27), we obtain the desired recursive formula

$$\begin{aligned} v_k - v_{k-1} &= \left[\left(\mathbf{I}_a^{(k)} - \mathbf{V}_a^{(k)}\right) \right. \\ &\quad \left. - \mathbf{V}_c^{(k)} \left(\mathbf{I}_b^{(k)} - \mathbf{V}_b^{(k)}\right)^{-1} \mathbf{V}_c^{(k)\top} \right]_{(k+1)(k+1)} |\Delta z_k|^2 \\ &= \left[\mathbf{B}^{(k)} \right]_{(k+1)(k+1)} |\Delta z_k|^2 \\ &\triangleq \mu_k |\Delta z_k|^2 \quad (\text{A.11}) \end{aligned}$$

where μ_k is the $(k+1, k+1)$ th element of $\mathbf{B}^{(k)}$.

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