STOCHASTIC APPROXIMATION WITH CORRELATED DATA
Technical Report

by

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SUMMARY

STOCHASTIC APPROXIMATION WITH CORRELATED DATA

New almost sure convergence results are developed for a special form of the multidimensional Robbins-Monro (RM) stochastic approximation procedure. The special form treated can be viewed as a stochastic approximation to the solution \( w = w_0 e^{RP} \) of the linear equations \( Rw = P \), where \( R \) is a pxp positive definite symmetric matrix. This special form commonly arises in adaptive signal processing applications. Essentially, previous convergence results for the RM procedure contain a common "conditional expectation condition" which is extremely difficult (if not impossible) to satisfy when the "training data" is a correlated sequence. In contrast, the new convergence results incorporate moment conditions and covariance function decay rate conditions. The ease with which these results can be applied in many cases is illustrated.
1. Introduction. Consider the set of linear equations \( Rw = P \),
where \( R \) is a \( p \times p \) symmetric positive definite matrix, and \( w \) and \( P \) are \( p \times 1 \) matrices. In case \( R \) and \( P \) are unknown, and the solution, \( w = w_o = R^{-1}P \), is desired, many techniques are available for finding an estimate of \( w_o \). In many adaptive signal processing applications, a recursive, computationally efficient procedure for estimating \( w_o \) is an important issue. A suitable multidimensional version of the Robbins-Monro (RM) stochastic approximation procedure (Robbins and Monro (1951)) for recursively estimating \( w_o \) is given by

\[
W_{n+1} = W_n + \mu_n (P_n - F_n W_n), \quad n \geq 1,
\]

where \( \{P_n\} \) is a sequence of random variables, \( P \in \mathbb{R}^p \), \( \{F_n\} \) is a sequence of random \( p \times p \) matrices, \( \{\mu_n\} \) is a sequence of positive constants, and \( W_1 \in \mathbb{R}^p \) is arbitrary. It is assumed that \( P_n \) and \( F_n \) are functionally independent of \( W_1, W_2, \ldots, W_n \). It is somewhat helpful to consider \( w_o \) to be the vector which minimizes \( \xi(w) = w'Rw - 2w'P \), where ' denotes matrix transpose. Interpreting \( F_n \) and \( P_n \) as "instantaneous estimates" of \( R \) and \( P \), respectively, the relationship between (1) and (deterministic) steepest descent procedures is obvious. Consequently, the family of algorithms represented by (1) has an interpretation as a family of "stochastic gradient-following" algorithms.

Algorithm (1) provides a suitable framework for the analytical treatment of many of the algorithms that have been proposed in the engineering literature for adaptive signal processing applications (e.g., see Sakrison (1966) or Farden (1975)). For such applications, any "conditional expectation assumption" is extremely difficult (if not impossible) to establish. Such conditions are commonly required by
existing convergence theorems treating the RM procedure. For example, see Schmetterer (1961, 1969), Sakrison (1966), or Farden (1975) for a discussion of existing results. The practical application of existing convergence results for the RM procedure to (1) essentially requires that \( \{P_n - F_w\} \) is an uncorrelated sequence for all fixed parameter \( w \in \mathbb{R}^p \). The special form of the RM procedure represented by (1) enables us to obtain convergence results which make use of no such conditional expectation requirements and have a decidedly different flavor than existing results for the RM procedure.

The contents and organization of this paper are as follows. Notation and basic assumptions are presented in Section 2. The framework presented in Assumption (2.1) establishes that the sequences \( \{P_n\} \) and \( \{F_n\} \) are such that the "time averages" of \( E(F_n) \) and \( E(P_n) \), respectively, are equal to \( R \) and \( P \). The generalization resulting from these definitions is applicable to cases where \( E(F_n) \) and \( E(P_n) \) are periodic, such as occurs in some adaptive digital communication applications. In Section 3, it is shown that Assumptions (2.1) through (2.7) are sufficient for the a.s. convergence of \( W_n \) to \( w_0 \). The proofs of Lemma (3.1) and Theorem (3.2) below are quite similar in spirit to the proofs of Theorem (6.1) and Theorem (6.3) of Albert and Gardner (1967), respectively. However, the seemingly less restrictive assumptions made in the present work, the simplification in proof resulting for symmetric \( F_n \), and the basic differences in the form of algorithms treated here are offered as justification for including the results of Section 3 in this paper. Furthermore, in contrast with the
assumptions made by Albert and Gardner (1967), the form of Assumptions (2.3) - (2.7) permits us to exploit the Borel-Cantelli Lemma and the results of Serfling (1970) to prove Corollary (4.5). Corollary (4.5) provides easily verified sufficient conditions for Assumptions (2.3)-(2.7), and hence, for the a.s. convergence of \( W_n \) to \( w_0 \). Several special cases of (1) are considered in Section 5 to illustrate the application of these results. In case \( F_n \) and \( P_n \) are strongly consistent estimates of \( R \) and \( P \), respectively, the much simpler convergence result of Section 6 is applicable.
2. Notation and basic assumptions. The norm of a pxp matrix A, denoted by \(|A|\), is defined here by \(|A| = \sup_{x \in \mathbb{R}^p} |Ax|\), where \(R^p\) denotes p-dimensional Euclidean space and \(|x|, x \in \mathbb{R}^p\) denotes the usual p-dimensional Euclidean norm. For A real and symmetric, \(|A| = \max_{i} \{\lambda_i(A)\}\), where \(\{\lambda_i(A)\}_{i=1}^{p}\) are the p eigenvalues of A. The minimum and maximum eigenvalues of a pxp matrix A are denoted by \(\lambda_{\min}(A)\) and \(\lambda_{\max}(A)\), respectively.

The element of a pxp matrix A occurring in the ith row and jth column of A is denoted by \((A)_{i,j}\). Similarly, the ith element of \(x \in \mathbb{R}^p\) is denoted by \((x)_i\). The trace of a pxp matrix A is denoted by \(\text{tr}(A) = \sum_{i=1}^{p} (A)_{i,i}\). The symbol 0 is used to denote the additive identity for \(\mathbb{R}, \mathbb{R}^n\), or to denote a pxp matrix of zeros. Square brackets \([\ ]\) are used to denote integer part. Finally, subscripted variables like \(v_k, n_k, k_1\), etc. are sometimes denoted by \(v(k), n(k), k(i)\), etc.

All random variables are assumed to be defined on a probability space \((\Omega, \mathcal{F}, P)\). All relationships between random variables are to be interpreted to hold with probability one.

It is worth emphasizing that Corollary (4.5) below establishes sufficient conditions for Assumptions (2.3)-(2.7).

(2.1) ASSUMPTION. The sequence \(\{W_n\}_{n=1}^{\infty}, W_n \in \mathbb{R}^p\), satisfies the recursion

\[
W_{n+1} = W_n + \mu_n (P_n - F_n W_n), \quad n \geq 1,
\]

where \(\{P_n\}\) is a sequence of random variables, \(P_n \in \mathbb{R}^p\), \(\{F_n\}\) is a sequence of real symmetric nonnegative definite (pxp) random matrices, \(\{\mu_n\}\) is a sequence of positive real numbers, and \(W_1 \in \mathbb{R}^p\) is arbitrary.

Define \(w_0 = \lim_{n \to \infty} \sum_{\ell=1}^{\infty} E(F_{\ell})\), \(P = \lim_{n \to \infty} \sum_{\ell=1}^{\infty} E(F_{\ell})\).

It is assumed that \(E(F_{\ell})\), \(E(F_{\ell})\) exist, that the above limits exist, and that \(R\) is positive definite. Further, it is assumed that \(\sum_{\ell=a+1}^{n} E(F_{\ell})\) converges uniformly to \(R\) as \(n \to \infty\) for all non-negative integers \(a\).
Defining $V_n = W_n - w_0$ and $C_n = P_n - F w_0$, we have

(2) $V_{n+1} = (I - \mu_n F_n)V_n + \mu_n C_n$.

For a sequence $\{A_i\}$ of $p \times p$ matrices, define

$$
A_i = \begin{cases} 
A_k A_{k-1} \ldots A_\ell, & \text{if } k \geq \ell \\
I, & \text{if } k < \ell.
\end{cases}
$$

Defining $Q_{k,m} = \prod_{j=\ell}^m (I - \mu_j F_j)$, $\Lambda_n = \sum_{k=1}^n Q_{k+1,n} \mu_k c_k$,

and iterating (2), one obtains

(4) $V_{n+1} = Q_{1,n} V_1 + \Lambda_n$.

(2.2) ASSUMPTION. The sequence $\{|\mu_n|\}$ is a nonincreasing sequence of positive constants $\mu_n = O(n^{-1})$, $0 < \lim_{n \to \infty} n \mu_n < \infty$.

(2.3) ASSUMPTION. Assumptions (2.1) and (2.2) hold and $u_n \|F_n\| \xrightarrow{a.s.} 0$ as $n \to \infty$.

(2.4) ASSUMPTION. Assumption (2.1) holds and $n^{-1} \sum_{k=a+1}^{a+n} F_k \xrightarrow{a.s.} 0$ as $n \to \infty$ for all positive integers $a$.

(2.5) ASSUMPTION. There exists a sequence of random integers $\{v_k\}$ with $1 = v_1 < v_2 < v_3 < \ldots$ such that, with $p_k = v_{k+1} - v_k$ and $J_k = \{v_k, v_k+1, \ldots, v_{k+1} - 1\}$ we have (i) $p_k = \xi[k]$, for some $\xi, 0 < \xi < 1$; (ii) $\min_{j \in J_k} (\sum_{j} F_j) \geq \delta > 0$, and (iii) $\max_{j \in J_k} (\sum_{j} F_j) \leq \gamma < \infty$.

The quantities $\xi$, $\delta$, and $\gamma$ are all random variables that are independent of $k$.

(2.6) ASSUMPTION. Assumptions (2.1) and (2.2) hold and there exists a random variable $S \in \mathbb{R}$ such that $S_n = \sum_{k=1}^n \mu_k c_k \xrightarrow{a.s.} S$ as $n \to \infty$.

(2.7) ASSUMPTION. Assumptions (2.1) and (2.6) hold and

$$
|F_n (S - S_{n-1})| \xrightarrow{a.s.} 0 \text{ as } n \to \infty.
$$

3. Almost sure convergence of $W_n$ to $w_0$. 

(3.1) LEMMA. If Assumptions (2.1)-(2.5) are satisfied, then

\[ |Q_{1,n}|^\alpha S \to 0. \]

**PROOF.** For any positive integer \( n \), let \( K = K(n) \) be such that

\[ n \in J_K. \]

Then \( Q_{1,n} = Q_{\nu(k),n} Q_{1,\nu(k)+1} \), and hence \( |Q_{1,n}| \leq |Q_{\nu(K),n}|. \)

**\[ |Q_{1,\nu(K)-1}|. \]** Defining \( \Gamma_k = Q_{\nu(k),\nu(k)+1} \), \( \Gamma_k \) may be expressed as

\[ \Gamma_k = I - \sum_{j \in J_K} \mu_j F_j + \sum_{q=2}^\infty (-1)^q \sum_{l_1 \geq 1} l_2 \ldots l_q \prod_{i=1}^q \mu_{k(1)} F_i(q-i+1), \]

where \( L_{k,q} = \{ l_1 \geq 1, l_2 \geq 1, \ldots, l_q \geq 1, i=1,2,\ldots,q \} \),

so that (for \( \delta \leq 1 \))

\[ |\Gamma_k|^2 \leq 1 - \mu_{\nu(k)+1} \Gamma_k^p + \sum_{q=2}^\infty (\mu_{\nu(k)p} \Gamma^q) - 1 - \mu_{\nu(k)+1} - 1 \Gamma_k^p, \]

where \( \alpha_k = \mu_{\nu(k)p} \Gamma^q \) and the last equality holds provided that \( \alpha_k \neq 1 \).

Assumptions (2.2) and (2.5) imply that there exists a random positive integer \( k_0 \) such that \( \alpha_k \leq 1/2 \) and \( |\Gamma_k|^2 \leq 1 - \beta_k \) for all \( k \geq k_0 \), where

\[ \beta_k = \frac{1}{2} \mu_{\nu(k)+1} - 1 \Gamma_k^p. \]

Similarly, for all \( K \geq k_0 \) and \( n \in J_K \), \( |Q_{\nu(K),n}| \leq 3/2 \).

It follows that there exists a random variable \( M \) such that for all \( K \geq k_0 \),

\[ |Q_{1,n}| \leq 3 M \pi \sum_{k=k_0}^K \mu_{\nu(k)+1} - 1 \Gamma_k, \]

since \( 1 - x \leq e^{-x} \) for all real \( x \). It is easily shown that Assumptions (2.2) and (2.5) imply that the above summation diverges to \( \infty \) as \( K \to \infty \).

Consequently, \( |Q_{1,n}|^{\alpha S} \to 0 \) as \( n \to \infty \).

(3.2) THEOREM. If Assumptions (2.1)-(2.7) are satisfied, then

\[ |\nu|^{\alpha S} \to 0 \text{ as } n \to \infty. \]

**PROOF.** From (2.1.4) and Lemma (3.1), it remains only to show that

\[ |A_n|^{\alpha S} \to 0 \text{ as } n \to \infty. \]

From Assumptions (2.1) and (2.6), with \( S_0 = 0 \) and \( Q_{n+1,n} = I \), we have
(1) \( \Lambda_n = \sum_{k=1}^{n} (Q_{k,n} - Q_{k+1,n}) S_{k-1} + S_n \)

\[ = - \sum_{k=1}^{n} Q_{k+1,n} v_{k,k} F_{k,k} S_{k-1} + S_n \]

Defining \( B_{k,n} = \sum_{k=1}^{n} Q_{k+1,n} v_{k,k} F_{k,k} (S - S_{k-1}) \),

and \( D_n = \sum_{k=1}^{n} Q_{k+1,n} v_{k,k} F_{k,k} \),

(1) may be expressed as \( \Lambda_n = B_{1,n} - D_n + S_n \).

From Lemma 1 of Albert and Gardner (1967, p. 189), \( D_n = (I - Q_{1,n})S \), so that \( \Lambda_n = B_{1,n} + (S - S) + Q_{1,n} S \). Since \( S_n \to S \) and \( ||Q_{1,n}|| \to 0 \) as \( n \to \infty \) (Assumption (2.6) and Lemma (3.1)), it remains to show that \( ||B_{1,n}|| \to 0 \) as \( n \to \infty \).

Using the same notation as in the proof of Lemma (3.1), \( |B_{1,n}| \) may be bounded as

\[ |B_{1,n}| \leq \left| \left| Q_{v(K),n} \right| \right| \cdot \left( \sum_{k=k_0}^{K-1} \frac{1}{k} \sum_{i=k+1}^{k} \sum_{j \in J_k} |Q_{j+1,v(k+1)-1}^i| F_{j,k} (S - S_{j-1}) \right) \]

\[ + \left| \left| Q_{v(K),n} \right| \right| \cdot (1 - \beta_\ell) \left[ \sum_{k=k_0}^{K-1} \frac{1}{k} \sum_{i=k+1}^{k} \sum_{j \in J_k} |Q_{j+1,v(k+1)-1}^i| F_{j,k} (S - S_{j-1}) \right] \]

\[ + \left| \sum_{k=k_0}^{K-1} Q_{k+1,n} v_{k,k} F_{k,k} (S - S_{k-1}) \right| \cdot \left( \sum_{k=k_0}^{K-1} \left| Q_{v(K),n} \right| \right) \cdot \left( \sum_{k=k_0}^{K-1} \frac{1}{k} \sum_{i=k+1}^{k} \sum_{j \in J_k} |Q_{j+1,v(k+1)-1}^i| F_{j,k} (S - S_{j-1}) \right) \]

Note that for all \( j, m \in J_k \) and \( k \geq k_0 \) we have \( ||Q_{j,m}|| \leq 1 + \sum_{q=2}^{p_k} \left( \mu_{v(K),p} \right)^q \frac{3}{2} \).

Consequently, defining \( d_k = \max_{j \in J_k} |F_j(S - S_{j-1})| \), there exists a random variable \( M \) such that
\[ |B_{ln}| \leq \frac{9}{4} \sum_{k=k_0}^{K-1} \frac{\pi}{i=k+1} (1-\beta_i) P_k \mu_v(k) d_k \]

\[ + \frac{3}{2} M_1 \sum_{i=k_0}^{K-1} (1-\beta_i) + \frac{3}{2} P_k \mu_v(k) d_k \]

Defining \( a_n,i = \frac{n}{(1-\beta_i)} \), it remains only to show that

\[ \sum_{k=k_0}^{K-1} a_{K-1,k} \beta_k^{-1} p_k \mu_v(k) d_k \to 0 \quad \text{as} \quad K \to \infty. \]

Clearly, for all fixed \( k \), \( a_{i,k} \to 0 \) as \( n \to \infty \). From Lemma 1 of Albert and Gardner (1967, p. 189), \( \sum_{i=k_0}^{n} |a_{n,i}| = 1 - \pi (1-\beta_i) \), which converges a.s. to 1 as \( n \to \infty \). Consequently, by the Toeplitz Lemma (e.g., see Knopp (1947, p. 75)),

\[ \lim_{K \to \infty} \sum_{k=k_0}^{K-1} a_{K-1,k} \beta_k^{-1} p_k \mu_v(k) d_k = \lim_{k \to \infty} \mu_v(k) d_k p_k \beta_k^{-1}. \]

By Assumptions (2.2) and (2.5), \( \mu_v(k) p_k \beta_k^{-1} = \frac{2}{\delta} \mu_v(k+1) \) is bounded; hence, \( |\beta_{1,n}| \to 0 \) as \( n \to \infty \).

(3.3) COROLLARY. If \( \{u_k\} \) satisfies Assumption (2.2) and \( ||F_k|| \) is a.s. uniformly bounded (in \( k \)), then Assumptions (2.3) and (2.7) may be deleted and Theorem (3.2) remains true.

PROOF. It suffices to consider \( u_k = k^{-1} \). Since

\[ |F_n(S-S_{n-1})| \leq ||F_n|| \cdot |S-S_{n-1}|, \]

Assumption (2.6) implies Assumption (2.7). The Borel-Cantelli Lemma and the Chebychev inequality can easily be applied to show that Assumption (2.3) is satisfied.

4. Sufficient conditions for Assumptions (2.3) - (2.7). Several auxiliary lemmas will be needed before the main result of this section, Corollary (4.5) may be stated and proved. The following lemma is a reasonably straightforward extension of Theorem A presented by

\[ \Sigma a_k, \Sigma b_k \to 0 \quad \text{as} \quad K \to \infty. \]
Serfling (1970). Consequently, the proof will be omitted.

(4.1) **Lemma.** Let \( \{x_i\} \) be a sequence of random variables, \( x_i \in \mathbb{R}^d \), having finite "variances" \( \sigma_i^2 = E[(x_i - E(x_i))^2] \). For integer \( n \geq 1 \) define \( X_{a,n} = (x_{a+1}, \ldots, x_{a+n}) \), \( S_{a,n} = \sum_{i=a+1}^{a+n} x_i \), and \( M_{a,n} = \max \{|S_{a,1}|, \ldots, |S_{a,n}|\} \). For \( n < 1 \) define \( X_{a,n} = (x_{a+n+1}, \ldots, x_a) \), \( S_{a,n} = \sum_{i=a+n+1}^{a} x_i \), and \( M_{a,n} = \max \{|S_{a,1}|, \ldots, |S_{a,n-1}|\} \). For \( |n| \geq 1 \) let \( F_{a,n} \) denote the distribution function for \( X_{a,n} \), and let \( g(F_{a,n}) \) be a functional depending on \( F_{a,n} \). Let \( a_0 \) be an arbitrary but fixed integer and let \( v \geq 2 \). Suppose \( g(F_{a,k}) + g(F_{a+k,l}) \leq g(F_{a+k+1,l}) \) for all \( a \geq a_0 \) and \( 1 \leq k < k + l \) or \( a_0 - a \leq k + l \leq -1 \) such that \( E(S_{a,n}^v) \leq \lambda^v(F_{a,n}) \) for all \( a \geq a_0 \) and \( n \geq 1 \) or \( a_0 - a \leq n \leq -1 \). Then \( E(M_{a,n}^v) \leq (\log 2|n|)^v \lambda^v(F_{a,n}) \) for all \( a \geq a_0 \) and \( n \geq 1 \) or \( a_0 - a \leq n \leq -1 \).

(4.2) **Lemma.** Let \( a_{k,l} = a_{k,l} \) and \( \rho(k,l) = \rho(k,l) \) be real-valued functions defined for all non-negative integers \( k,l \). For \( 1 \leq n < m \) define

\[
(1) \quad \gamma_{n,m} = \sum_{k=n}^{m-1} \sum_{l=n}^{m} \alpha_{k,l} \rho(k,l)
\]

\[
= 2 \sum_{u=1}^{m-n} \sum_{k=n}^{m-u} \alpha_{k,k+u} \rho(k,k+u) + \sum_{k=n}^{m-1} \alpha_{k,k} \rho(k,k).
\]

Suppose that \( |\rho(k,k+u)| = O(u^{-v}) \) uniformly for all positive integers \( k \). If \( a_{k,l} = 1 \) and \( 0 < v < 1 \), then for large \( m-n \),

\[
(2) \quad |\gamma_{n,m}| = O((m-n)^{2-v}).
\]

Finally, if \( a_{k,k} = \nu_k \nu_{k+1} \), \( \nu_k = O(k^{-1}) \), and \( 0 < v < 1 \), then

\[
(3) \quad |\gamma_{n,m}| = O(n^{-v/(v+1)}).
\]
PROOF. Suppose that \( \alpha_k, \xi = 1 \) and \( |p(k,k+u)| = O(u^{-\nu}), \ 0 < \nu < 1 \), uniformly for all positive integers \( k \). It suffices to consider
\[
|p(k,k+u)| = (|u|^{1+\nu}) \text{in which case (from (1))}
\]
\[
|\gamma_{n,m}| \leq 2(m-n+1) \sum_{u=0}^{m-n} (u+1)^{-\nu} 2(m-n+1)(1+\int_1^x x^{-\nu} dx).
\]
The result (2) follows easily by evaluating the above integral.

Suppose now that \( \alpha_k, \xi = \nu_k \), \( \nu_k = O(k^{-1}) \), and \( |p(k,k+u)| = O(u^{-\nu}), \ 0 < \nu < 1 \), uniformly for all positive integers \( k \). For this case, it suffices to consider \( \nu_k = k^{-1} \) and \( |p(k,k+u)| = (|u|+1)^{-\nu} \). Then from (1)
\[
|\gamma_{n,m}| \leq 2 \sum_{u=1}^{m-n} (u+1)^{-\nu} \int_1^{u+1} x^{-\nu} dx + (n-1)^{-1} - m^{-1}
\]
\[
= 2 \sum_{u=1}^{m-n} u^{-1} (u+1)^{-\nu} \beta_{n,m,u} + (n-1)^{-1} - m^{-1},
\]
where \( \beta_{n,m,u} = \ln ((m-u)(n+u-1)) - \ln (m(n-1)) \). If \( 1 \leq u \leq n < m \), then \( \beta_{n,m,u} \leq \ln (n+u-1) - \ln (n-1) \). If \( 1 \leq u < n < m-n \), then \( \beta_{n,m,u} \leq \ln 2 \). If \( n \leq u \leq m-n \), then \( \beta_{n,m,u} \leq \ln (u) \).
Consequently, for all \( 1 \leq \ell \leq n-2 < m-n-2 \),
\[
|\gamma_{n,m}| \leq 2 \ell \ln \left( \frac{n+\ell-1}{n-1} \right) + 2 \ell \ln 2 \sum_{u=\ell+1}^{n-1} u^{-1-\nu} u + (n-1)^{-1} - m^{-1}.
\]
Letting \( \ell = n^\beta \), \( 1 > \beta > 0 \), and using the fact that \( \ln (1+x) \leq x \) for all \( x > -1 \), it follows that there exist constants \( c_1, c_2, \text{ and } c_3 \) such that
\[
|\gamma_{n,m}| \leq c_1 n^{2\beta-1} + c_2 n^{-\beta\nu} + c_3 n^{-\nu/2}.
\]
Finally, if \( \beta = (\nu+2)^{-1} \), then \( |\gamma_{n,m}| = O(n^{-\nu/(\nu+2)}) \). \( \square \)
(4.3) LEMMA. Define

\[ \gamma_F(i, j, k, m) = \max_{|w| = 1} \omega \in \mathbb{R}^P \omega^T (F_i - \mathbb{E}(F_i)) \omega^T (F_j - \mathbb{E}(F_j)) \omega^T (F_k - \mathbb{E}(F_k)) \omega^T (F_m - \mathbb{E}(F_m)) \omega, \]

where \( \omega \in \mathbb{R}^P \). Define \( P_k(a) = a + [k^a] \), where \( a \) is a positive integer,

\( 0 \leq a < 1 \). Define the sequence \( \{v_k(a)\} \) by \( v_1(a) = 1, v_{k+1}(a) = v_k(a) + P_k(a) \), \( k = 1, 2, \ldots \), and define \( J_k(a) = \{v_k(a), v_k(a)+1, \ldots, v_{k+1}(a)-1\} \). If

Assumption (2.4) holds and

\[ \sum_{k=1}^{\infty} P_k(a) \sum_{i,j,k,m} \omega \in \mathbb{R}^P \] \[ \gamma_F(i, j, k, m) < \infty, \]

for some \( a, 0 \leq a < 1 \), and for some positive integer \( a \), then Assumption (2.5) is satisfied.

PROOF. Define \( C_k(a) = p_k^{-1}(a) \sum_{j \in J_k(a)} F_j \) and \( R_k(a) = p_k^{-1}(a) \sum_{j \in J_k(a)} E(F_j) \).

Let \( \varepsilon \) be given such that \( 0 < \varepsilon < \lambda_{\min}(R) \). Since

\[ \lambda_{\min}(C_k(a)) - \lambda_{\min}(G_k(a) - R) + \lambda_{\max}(R) \]

and \( \lambda_{\max}(G_k(a)) - \lambda_{\min}(G_k(a) - R) + \lambda_{\max}(R) \),

it is sufficient to show that there exists a random sequence \( \{v_k(\xi)\} \), \( \xi \)

an integer-valued random variable with \( \xi \) a.s. finite such that

\[ 0 < \lambda_{\min}(R) - \varepsilon < \lambda_{\min}(G_k(\xi) - R) + \lambda_{\min}(R) < \lambda_{\max}(G_k(\xi) - R) + \lambda_{\max}(R) \]

\[ < \lambda_{\max}(R) + \varepsilon \]

for all \( k \), or, equivalently, that

\[ \max_{|w| = 1} |w^T(C_k(\xi) - R)w| < \varepsilon \]

for \( k = 1, 2, \ldots \). Hence, it is also sufficient to require the stronger condition that

\[ \max_{|w| = 1} |w^T(G_k(\xi) - R(\xi))w| + \max_{|w| = 1} |w^T(R_k(\xi) - R)w| < \varepsilon \]

for all \( k \). By Assumption (2.1), there exists a positive integer \( q_1 \)

such that for all sequences \( \{v_k(a)\} \) with \( a \geq q_1 \) we have
max \(|w'(R_k(a) - R_k))w| < \epsilon/2\). It follows from Assumption (2.4) that for all \(n = 1, 2, \ldots\), there exists an a.s. finite random variable \(\xi_n^*\) such that

\[
\max_{1 \leq k \leq n} \max_{|w| = 1} |w'(G_k(\xi_n^*) - R_k(\xi_n^*)))w| < \epsilon/2.
\]

Let \(\xi_n\) be the smallest such \(\xi_n^*\) so that \(\xi_n \geq q_1\).

It follows that

\[
P(\xi_n > a > q_1) \leq P(\max_{|w| = 1} |w'(G_k(a) - R_k(a)))w| \geq \epsilon/2 \text{ for some } k, 1 \leq k \leq n)
\]

\[
\leq \sum_{k=1}^{n} P(\max_{|w| = 1} |w'(G_k(a) - R_k(a)))w| \geq \epsilon/2)
\]

\[
\leq \sum_{k=1}^{n} (\frac{\epsilon}{p_k(a)})^{-4} E(\max_{|w| = 1} |w'(G_k(a) - R_k(a)))w|^4)
\]

\[
= \left(\frac{\epsilon}{\epsilon}\right)^4 \sum_{k=1}^{n} (p_k(a))^{-4} \sum_{i,j,\ell \in \mathcal{J}_k(a)} \gamma_f(i,j,\ell,m).
\]

By hypothesis the above series converges as \(n \to \infty\). Applying the definition of \(p_k(a)\), it follows that for all \(\epsilon_1 > 0\), there exists a positive integer \(N(\epsilon_1)\) such that for all \(a \geq N(\epsilon_1)\) we have

\[
\lim_{n \to \infty} P(\xi_n > a > q_1) < \epsilon_1.
\]

Consequently, \(\xi = \sup_{n} \xi_n\) is a.s. finite and Assumption (2.5) is satisfied.

**Lemma.** If there exists a real-valued function \(f(k,k+u) = O(|u|^{-\beta})\), \(\beta > \frac{1}{2}\), uniformly for all positive integers \(k\) such that

\[
|Y_f(i,j,h,m)| \leq f(i,j) f(h,m) + f(i,h) f(i,m) f(j,m) f(i,m),
\]

then Assumption (2.5) is satisfied.

**Proof.**

It suffices to consider \(f(i,j) = (|i-j|+1)^\beta\).
Define the $ij$th element of a matrix $A$ to be $f(i,j)$. Then
\[
\sum_{i=a}^{b} \sum_{j=a}^{b} \sum_{\ell=a}^{\ell} \sum_{m=a}^{m} f(i,j)f(i,\ell)f(j,m)f(\ell,m) = \text{tr}(A^4) \leq (\text{tr}A^2)^2.
\]

The $ii$th element of $A^2$ is given by
\[
(A^2)_{ii} = \sum_{m=a}^{b} f(i,m+1)^{-2\beta} + \sum_{m=i+1}^{b} f(m,i+1)^{-2\beta} + 1
\leq c_1(i-a+1)^{1-2\beta} + c_2(b-i+1)^{1-2\beta},
\]
for $0<\beta<1/2$ and some positive constants $c_1$ and $c_2$. Consequently, for some constant $c_3$, $\text{tr}(A^2) \leq c_3(b-a+1)^{-2-2\beta}$, and hence $\text{tr}(A^4) = O((b-a)^4\beta)$. For this sequence, \((4.3.1)\) be satisfied if $\sum_{k=1}^{\infty} k^{-4\alpha} k^{-4\alpha-4\alpha\beta} < \infty$. It follows that Assumption (2.5) is satisfied for $1>\alpha>1/4\beta$ if $\beta > \frac{1}{4}$.

(4.5) COROLLARY. Suppose that Assumptions (2.1), (2.2), and (2.5) are satisfied. Define $\rho_F(k,l) = E(F_kF_l) - E(F_k)E(F_l)$, $\rho_C(k,l) = E(C_kC_l) - E(C_k)E(C_l)$, and $\rho_{FC}(k,l,n) = E((C_k - E(C_k))F_n(C_l - E(C_l)))$, where all expectations are assumed to exist. Suppose that for some $\nu>0$, \(\nu^\alpha \max \{||\rho_F(k_ku)||, ||\rho_C(k_ku)||, ||\rho_{FC}(k_ku, n)||\} \) is uniformly bounded for all non-negative integers $k, u, n$. Further suppose that $\gamma_n = \nu_n E(C_k) \in E(\nu_n^\beta)$ exists and that there exists a constant $\beta > 1$ such that $\sum_{n=1}^{\infty} \nu_n^\beta E(||\nu_n||^\beta) < \infty$. Then $\nu_n \rightarrow 0$ a.s. as $n \rightarrow \infty$.\]
PROOF. First consider Assumption (2.4). Define \( S_{a,n} = \sum_{k=a+1}^{a+n} (F_k - E(F_k))w \), where \( w \in R^P \). Clearly, Assumption (2.4) is satisfied iff
\[
-n^{-1} S_{a,n} a, S. 0
\]
for all positive integers \( a \) as \( n \to \infty \). Define
\[
M_{a,n} = \max \{ |S_{a,1}|, \ldots, |S_{a,n}| \}. \text{ Let } \{n_k\} \text{ be an increasing sequence}
\]
of positive integers such that \( n_k \to \infty \) as \( k \to \infty \). For all \( n_k \leq n \leq n_{k+1}-1 \),
\[
(1) \quad n^{-1} S_{a,n-k} S_{a,n(k)}^{-1} + n_{k-1}^{-1} M_{a+n(k)-1,n(k+1)-n(k)}^{-1}
\]
To apply Lemma (4.1), define \( g(F_{a,n}) \) as
\[
(2) \quad g(F_{a,n}) = |w|^2 \sum_{k=a+1}^{a+n} \sum_{\ell=a+1}^{a+n} \|p \|_{F(k,\ell)}\| = O(n^{2-v}), 0 < v < 1,
\]
uniformly for all positive integers \( a \), from Lemma (4.2). It is
easily seen that \( g(F_{a,n}) \) satisfies the needed conditions of Lemma
(4.1). Letting \( n_k = k^a, n_{k-2} E(|S_{a,n(k)}|^2) \) is summable for all \( a>v^{-1} \).
The Chebychev inequality and Borel-Cantelli Lemma thus imply that
\[
(3) \quad n_{k-2} E(\xi_k^2) = O(n_{k-2}^v (log_2 2(n_{k+1} - n_k)))^2,
\]
from (2) and Lemma (4.1). Substituting \( n_k = k^a \) into (3), the Borel-
Cantelli Lemma and the Chebychev inequality imply that \( n_{k-1}^{-1} \xi_k \to S. 0 \)
as \( k \to \infty \) for \( a>v^{-1} \). Consequently, from (1), Assumption (2.4) is satisfied.

Now consider Assumption (2.6). Define \( S_{a,m} = \sum_{k=a+m+1}^{a} \nu_k (C_k - E(C_k)) \),
\[
M_{a,m} = \max \{ |S_{a,m}|, \ldots, |S_{a,m-1}| \}, \text{ and } g(F_{a,m}) = \sum_{k=a+m+1}^{a} \sum_{\ell=a+1}^{a+m} \|p \|_{C(k,\ell)}\|
\]
It is easily seen that \( g(F_{a,m}) \) satisfies the needed conditions of
Lemma (4.1). From
Lemma (4.2), $E(|S_{a,n-a-1}|^2) = O(n^{-\nu/(\nu+2)})$ for all $a > n$. With $n_k = k^\alpha$, the Chebychev inequality and the Borel-Cantelli Lemma thus imply that

$$S_n(k) - S_n(k-1) = \sum_{i=n_k}^n (C_i - E(C_i)) = O(n_k^{\alpha}),$$

as $k \to \infty$ for $\alpha > \nu^{-1}(\nu+2)$.

Consequently, $S_n(k) - S_n(k-1) = O(n_k^{\alpha})$, $\alpha > \nu^{-1}(\nu+2)$, as $k \to \infty$. For all $n_k - n \leq n_{k+1} - 1$,

$$|S_{n-1} - S_{n_{k+1} - 1}| \leq \xi_k + b_k,$$

where $\xi_k = M_{n(k+1)-1,n(k)-n(k+1)}$ and $b_k$ is a sequence of positive constants converging to zero as $k \to \infty$. From Lemmas (4.1) and (4.2),

$$E(\xi_k^2) = O((\log_2(2(n_{k+1}-n_k)))^2 n_k^{-\nu/(\nu+2)}).$$

Substituting $n_k = k^\alpha$ into (5), the Chebychev inequality and the Borel-Cantelli Lemma imply that $\xi_k \to 0$ as $k \to \infty$ for all $\alpha > \nu^{-1}(\nu+2)$.

Consequently, Assumption (2.6) is satisfied.

Finally, consider Assumption (2.7). Define $Z_n(k) = F_n(C_k - E(C_k))$,

$$S_{a,m} = \sum_{k=a+m+1}^n \mu_k Z_{k,a+m+1}, \quad a \leq m \leq a \leq n,$$

and $g(F_{a,m}) = \sum_{k=a+1}^n \sum_{\ell=a+m+1}^{n_k} \mu_k \mu_\ell |\rho_{FC}(k,\ell,a+m+1)|$.

Proceeding as above, $E(|S_{a,n-a-1}|^2) = O(n^{-\nu/(\nu+2)})$ for all $a > n$.

With $n_k = k^\alpha$, $\alpha > \nu^{-1}(\nu+2)$, we have

$$F_n(k) (S_{n(k)} - S_{n(k)-1}) = F_n(k) g_n(k) = \sum_{i=n_k}^\infty \mu_i Z_{i,n(k)} = n_k^{\alpha},$$

as $k \to \infty$. Since $\sum_{n=1}^\infty |g_n| = E(||F_n||)^{\nu}$, the Markov inequality and the Borel-Cantelli Lemma imply that $F_n g_n \to 0$ as $n \to \infty$. Consequently, for all $n_k - n < n_{k+1} - 1$,

$$|F_n(S_{n-1}) - F_n(S_{n-1})| \leq |F_n(k+1) (S_{n(k+1)-1}) + \xi_k + \psi_k|,$$

where $\xi_k = M_{n(k+1)-1,n(k)-n(k+1)}$ and $\psi_k \to 0$ as $k \to \infty$. Since $\xi_k$ satisfies (5), $\xi_k \to 0$ as $k \to \infty$. Consequently, Assumption (2.7) is satisfied.
Finally, $\beta > 1$ and $\mu_n = O(n^{-1})$ easily provide that $V_n \alpha_s |F_n|^s \rightarrow 0$ by the Markov inequality and the Borel-Cantelli Lemma, and hence, Assumption (2.3) is satisfied.

Consequently, from Theorem (3.2), $V_n \alpha_s \rightarrow 0$ as $n \rightarrow \infty$. ⊠

(4.6) REMARK. An important ancillary result contained in Corollary (4.5) is that sufficient conditions for the strong consistency of the usual sample covariance function are provided. For example, let $\{x_k\}_{k=-\infty}^{\infty}$ be a zero mean wide-sense stationary real-valued normal random process and define $\rho_{x}(v) = E(x_k x_{k+v})$. Consider

$F_k = F_{k}(u) = x_k x_{k+u}$. It is easily shown that for this case, $\rho_{x}(k,k+v) = |\rho_{x}(v) + \rho_{x}(v+u)\rho_{x}(v-u)|$. The proof of Corollary (4.5) shows that if $\rho_{x}(u) = O(u^{-1})$ for $v > 0$, then $n^{-1} \sum_{k=1}^{n} x_k x_{k+u} \rightarrow \rho_{x}(v)$ as $n \rightarrow \infty$. A similar result, presented as Theorem 8B of Parzen (1961), states that $n^{-1} \sum_{k=1}^{n} x_k x_{k+u} \rightarrow \rho_{x}(v)$ as $n \rightarrow \infty$ provided that there exist positive constants $c, q$ such that $n^{-1} \sum_{u=0}^{n} \rho_{x}^2(u) \leq cn^{-q}$ for all positive integers $n$. The conditions of Parzen (1961) are clearly satisfied by $\rho_{x}(u) = O(u^{-1})$ and $v > 0$.

(4.7) REMARKS. Recall that sufficient conditions for Assumption (2.5) have been presented as Lemmas (4.3) and (4.4). Lemma (4.4) is useful for several specific choices of $\{F_{n}\}$, as shown below in Section 5. In case $||F_{n}||$ is bounded, or if $\{F_{n}\}$ and/or $\{P_{n}\}$ are deterministic, then the conditions of Corollary (4.5) are simplified, as shown below in Corollaries (4.8)-(4.11).

(4.8) COROLLARY. Let $\rho_{F}, \rho_{C},$ and $g_{n}$ be as in Corollary (4.5). Suppose that Assumptions (2.1), (2.2), and (2.5) are satisfied, and that $||F_{n}||$ is bounded. If there exists a $v > 0$ such that $u^v \max \{||\rho_{F}(k,k+u)||, ||\rho_{C}(k,k+u)||\}$ is uniformly bounded for all non-negative
integers \( k \) and \( u \), and there exists a \( \beta > 0 \) such that \( \sum_{n=1}^{\infty} |g_n|^\beta < \infty \), then \( |V_n|^\alpha \rightarrow 0 \) as \( n \to \infty \).

**Proof.** Simply apply Corollary (3.3) to Corollary (4.5).

(4.9) **COROLLARY.** Suppose that \( F_n \) is deterministic, and that Assumptions (2.1) and (2.2) are satisfied. Define \( \rho_P(k,\ell) = E(F_k F_\ell) - E(F_k)E(F_\ell) \), and \( g_n = \sum_{k=1}^{\infty} \mu_k (E(F_k) - F_k) \). If there exists a \( \nu > 0 \) such that \( \nu \max \{|\rho_P(k,\ell)|, |\rho_P(k,\ell,n)|\} \) is uniformly bounded for all non-negative integers \( k \) and \( \ell \), and if \( g_n \rightarrow 0 \) as \( n \to \infty \) then \( |V_n|^\alpha \rightarrow 0 \) as \( n \to \infty \).

**Proof.** An obvious consequence of Corollary (4.5).

(4.10) **COROLLARY.** Suppose that \( \{P_n\} \) is deterministic, and that Assumptions (2.1), (2.2), and (2.5) are satisfied. Define \( \rho_{PP}(k,\ell,\nu) = E((F_k - E(F_k))^2) - E((F_k - E(F_k))^2) \), and \( g_n = \sum_{k=1}^{\infty} \mu_k (E(F_k) - F_k) \). If there exists a \( \nu > 0 \) such that \( \nu \max \{||\rho_P(k,\ell)|, |\rho_{PP}(k,\ell,\nu)|\} \) is uniformly bounded for all non-negative integers \( k \), \( \ell \), and \( \nu \), and there exists a \( \beta > 1 \) such that \( \sum_{n=1}^{\infty} |g_n|^\beta \leq \sum_{n=1}^{\infty} |g_n|^\beta < \infty \), then \( |V_n|^\alpha \rightarrow 0 \) as \( n \to \infty \).

**Proof.** Follows directly from Corollary (4.5).

(4.11) **COROLLARY.** Suppose that both \( \{F_n\} \) and \( \{P_n\} \) are deterministic, and that Assumptions (2.1) and (2.2) are satisfied. If \( \sum_{k=1}^{\infty} \mu_k (P_k - E(F_k)) \) exists, then \( |V_n| \rightarrow 0 \) as \( n \to \infty \).

**Proof.** Trivial case of Corollary (4.5).

5. **Application of Corollary **(4.5).

(5.1) **Special families of \( F_n \) and \( P_n \).** Let \( \{X_{ij}\}^{\infty}_{i=-\infty} \) be a sequence of \( R^p \)-valued zero-mean random variables and let \( \{s_{ij}\}^{\infty}_{i=-\infty} \) be a sequence of real-valued zero-mean random variables. Define \( R_{xx}(k,\ell) = E(X_k X_{\ell}^\prime) \),
\[ P_{\mathbf{S}}(k,\ell) = E(s_kX_{\ell_k}) \text{ and } P_{\mathbf{S}}(k,k) = E(s_k s_{\ell_k}). \] Suppose that \( R_{\mathbf{XX}}(k,k+u) \), \( P_{\mathbf{S}}(k,k+u) \), and \( \rho_{\mathbf{S}}(k,k+u) \) are periodic in \( k \) with period \( N \). Define \( R = N^{-1} \sum_{k=1}^{N} R_{\mathbf{XX}}(k,k) \) and \( P = N^{-1} \sum_{k=1}^{N} P_{\mathbf{S}}(k,k) \). It will become apparent in what follows that \( R \) and \( P \) satisfy Assumption (2.1).

Suppose that it is desired to choose \( w \in \mathbb{R}^p \) to minimize \( \xi(w) = \sum_{k=1}^{N} E((s_k - w \cdot x_k)^2) \). Such problems arise frequently in adaptive transversal filter channel equalization in digital communications. When \( N = 1 \), the problem reduces to the use of jointly wide-sense stationary sequences \( \{s_j\} \) and \( \{x_j\} \). Assuming that \( R \) is positive definite, the desired solution, \( w_o \), is given by \( w_o = R^{-1}P \). Assume now that \( R \) and/or \( P \) are unknown, and that it is desired to use algorithm (2.1.1), with \( F_n \) and \( P_n \) functions of the observed time series, \( \{x_j\} \) and \( \{s_j\} \), to recursively estimate \( w_o \). Obvious candidates for \( F_n \) and \( P_n \) are

1. \( F_n = K_n^{-1} \sum_{j=n-K+1}^{n} x_j x_j^\prime \),
2. \( P_n = K_n^{-1} \sum_{j=n-K+1}^{n} s_j x_j \),

where \( K_n \) is a positive integer; e.g. \( K_n = 1, K, \) or \( n \). In fact, algorithms represented here by (2.1.1) with \( F_n \) and \( P_n \) given by (1) and (2) have frequently appeared in the engineering literature for consideration in a wide range of applications.

Note that if \( K_n = K \) (a constant), then \( E(F_n) \), \( E(P_n) \), and \( E(C_n) \) are all periodic (in \( n \)) with period \( N \). Furthermore, for any \( n > 0 \), \( N^{-1} \sum_{k=1}^{n+N} E(C_k) = 0 \). If, in addition \( \{u_k\} \) satisfies either \( u_k = a(k+b)^{-1} \) or \( u_k = a(k+b)^{-1} \), where \( a > 0 \) and \( b > 0 \), then it can be
shown that \[ |g_n| = \sum_{k=n}^{\infty} \mu_k E(C_k) = O(n^{-1}) \]. If, for example, \( E(||F_n||^2) \) is bounded (in \( n \)), then \( \sum_{n=1}^{\infty} |g_n|^2 E(||F_n||^2) < \infty \), thus satisfying the conditions on \( \{g_n\} \) stated in Corollaries (4.5) and (4.8)-(4.10). Of course, many other choices of \( \{\mu_k\} \) are permissible. Finally, if \( K_n = N \), then \( E(C_n) \leq 0 \), and the conditions on \( \{g_n\} \) are a fortiori satisfied.

The remaining conditions on \( \{F_n\} \) and \( \{P_n\} \) stated in the preceding corollaries are quite mild "asymptotic covariance decay rate" conditions. The strongest of these decay rate conditions is that imposed on \( \gamma_F \) via Lemma (4.3) in order to satisfy Assumption (2.5).

Regarding the covariance decay rate conditions in the corollaries of Section 4, it may be helpful to note that

\[ \rho_C(k,\ell) = \rho_F(k,\ell) + \omega_0 \rho_{PF}(k,\ell) \omega_0 - \rho_{PF}(\ell,k) \omega_0 - \rho_{PF}(\ell,k) \omega_0, \]

where \( \rho_F \) and \( \rho_F \) are defined in Corollaries (4.9) and (4.5), respectively, and \( \rho_{PF}(k,\ell) = E(P_k^* F_{\ell}) - E(P_k^*) E(F_{\ell}). \) Furthermore, \( \rho_{FC} \) can be expressed as

\[ \rho_{FC}(k,\ell,n) = \rho_{PF}(k,\ell,n) + \omega_0 \rho_{PF}(k,\ell,n) \omega_0 - \rho_{PF}(\ell,k,n) \omega_0 - \rho_{PF}(\ell,k,n) \omega_0, \]

where \( \rho_{PF} \) is defined in Corollary (4.10) and \( \rho_{PF}(k,\ell,n) = E(P_k^* - E(P_k^*)) \cdot F_n^2(P_{k,\ell} - E(P_{k,\ell})), \) and \( \rho_{PF}(k,\ell,n) = E((P_k^* - E(P_k^*)) F_n^2(P_{k,\ell} - E(F_{k,\ell}))). \) The results (3) and (4) simply express the conditions on \( \{F_n\} \) and \( \{C_n\} \) stated in the corollaries of Section 4 in terms of similar conditions on \( \{F_n\} \) and \( \{P_n\} \).

In view of the widespread consideration of algorithms fitting the framework of (2.1.1), (1), and (2), it is of interest to reduce the moment conditions on \( \{F_n\} \) and \( \{P_n\} \) stated in the corollaries of Section 4 to moment conditions on \( \{X_j\} \) and \( \{s_j\} \). In case \( \{X_j\} \) and \( \{s_j\} \) are normally distributed, greatly simplified conditions can be established for this family of algorithms, as shown in Section (5.2).
(5.2) The normal case. Assume the same notation and structure as in Section 5.1. Further, assume that all random variables involved are jointly normally distributed. Then all of the "covariance decay rate conditions" for \( \{F_n\} \) and \( \{L_n\} \) stated in the corollaries of Section 4 and Section 5.1 can be expressed in terms involving the covariance functions of elements of \( \{X_j\} \) and \( \{s_j\} \). In order to accomplish this, some properties of joint moments of normally distributed random variables are required.

Let \( z_1, z_2, \ldots, z_8 \) be zero-mean jointly normally distributed with \( E(z_iz_j) = \sigma(i,j) \). Using the properties of the characteristic function of \( z_1, z_2, \ldots, z_8 \), it is straightforward to show that

\[
(1) \quad E(\prod_{i=1}^{4} (z_{2i-1}z_{2i} - \sigma(2i-1, 2i))) = \sum_{\ell_1, \ell_2, \ell_3, \ell_4} \sigma(\ell_1, \ell_2)\sigma(\ell_3, \ell_4) \cdots \sigma(\ell_7, \ell_8),
\]

where the summation is over all possible ways of combining \( \ell_1, \ell_2, \ldots, \ell_8 \in\{1, 2, \ldots, 8\}, \ell_i \neq \ell_j \) for \( i \neq j \), into four distinct unordered pairs, no pair of which is \( (1, 2), (3, 4), (5, 6), \) or \( (7, 8) \). There are sixty terms in the sum.

With \( F_n \) given by (5.1.1) and \( K_n = K \), and defining \( f(k, \ell) = \max_{1 \leq i, j \leq p} |(R_{xx}(k, \ell))_{i,j}| \), (1) can be applied to show that the conditions of Lemma (4.4) are satisfied if for some \( \nu > 1/4 \),

\[
\nu \max_{1 \leq i, j \leq p} |(R_{xx}(k, k+u))_{i,j}|
\]

is uniformly bounded for all non-negative integers \( k \) and \( u \). Similarly, if for some \( \nu_1 > 0 \),
$$u^1 \max_{1 \leq i \leq p} \{ |(P_s(k,k+u))_{k}|, |\rho_s(k,k+u)| \}$$
is uniformly bounded for all non-negative integers $k$ and $u$, then the
decay rate conditions stated in Corollary (4.5) on $\rho_P$, $\rho_C$, and $\rho_{PC}$ are
all satisfied. Consequently, each member of the family of algorithms
represented by (2.1.1) with $F_n$ and $P_n$ given by (5.1.1) and (5.1.2),
and $K_n = K$ converge a.s. to $w_0$ provided that, in addition, $\{u_k\}$ and $K$
are chosen as discussed in Section 5.1.

For example, let $\{x_k : k \in \{0, \pm 1, \ldots \}\}$ be a real-valued zero-mean wide-sense station-
ary autoregressive-moving average (ARMA) normal random process. With
$X_k = (x_k, x_{k-1}, \ldots, x_{k-p+1})'$ and $s_k = x_{k+a}$ (integer $a$), the required
decay rates of (5.1.3) are easily established. Hence, each member of
the family of algorithms represented by (2.1.1) with $F_n$ and $P_n$ given
by (5.1.1) and (5.1.2), $K_n = K$, and $\{u_n\}$ satisfying Assumption (2.2)
converge a.s. to $w_0$.

6. A simple almost sure convergence result.

(6.1) THEOREM. Suppose that Assumption (2.1) holds and that there
exist sequences $\{a_n\}_{n=1}^{\infty}$ and $\{b_n\}_{n=1}^{\infty}$ of non-negative real numbers (possibly
random) satisfying $\|F_n - R\| \leq a_n$ and $\|F_n - P_n\| \leq b_n$. Furthermore, suppose
that there exists a positive integer $n_o$ (possibly random) such that for
all $n > n_o$, $0 < u_n \lambda_{\min}^o (R) < 1$. Then for all $n > n_o$,

(1) $|V_{n+1}| \leq |V_n| + \frac{1}{n} \max_{k=n_o}^n \left( 1 - u_k d_k \right) + \frac{1}{n} \sum_{j=n_o}^n \left( 1 - u_j d_j \right) \left( 1 - u_j d_j \right)$,

where $d_k = \lambda_{\min}^o (R) - a_k$. Furthermore, if $\sum_{k=1}^n u_k d_k \rightarrow^s \infty$ and $b d_n \rightarrow^s 0$ as
$n \rightarrow \infty$, then $|V_n| \rightarrow^s 0$ as $n \rightarrow \infty$. 
PROOF. From Assumption (2.1),

\[ V_{n+1} = V_n - \mu RV_n - \mu (F V + F \omega - P - RV), \]

so that for all \( n \geq n_0 \),

\[ |V_{n+1}| \leq (1 - \mu \min(k)) |V_n| + \mu |F_n - R| \cdot |V_n| + \mu |F_n \omega - P| \]

\[ \leq (1 - \mu \frac{d}{n}) |V_n| + \mu \frac{b_n}{n}. \]

Iterating (2), for all \( n \geq n_0 \),

\[ |V_{n+1}| \leq \sum_{k=n_0}^{n} (1 - \mu \frac{d_k}{k}) + \frac{\mu}{k} \sum_{j=k+1}^{n} (1 - \mu \frac{d_j}{j}) \frac{d_k}{k} \frac{d_j}{j} (b_j d_j^{-1}). \]

Since all terms appearing in the sum in (3) are a.s. non-negative, (3) follows immediately from (3) with the aid of Lemma 1 of Albert and Gardner (1967, p. 189). Furthermore, if \( \sum_{k=1}^{n} \mu_k d_k \frac{a}{s} \frac{\infty}{\infty} \) and \( b_n d_n^{-1} \frac{a}{s} \frac{\infty}{\infty} \) as \( n \to \infty \), then (3) and the Toeplitz Lemma show that \( |V_n| \to a.s. 0 \) as \( n \to \infty \).

(6.2) Remark. Theorem (6.1) is applicable to the case discussed in Section 5 with \( K = k \).
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New almost sure convergence results are developed for a special form of the multidimensional Robbins-Monro (RM) stochastic approximation procedure. The special form treated can be viewed as a stochastic approximation to the solution \( w = w_{RM}^{RP} \) of the linear equations \( R w = P \), where \( R \) is a \( p \times p \) positive definite symmetric matrix. This special form commonly arises in adaptive signal processing applications. Essentially, previous convergence results for the RM procedure contain a common conditional expectation condition which is
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extremely difficult (if not impossible) to satisfy when the "training data" is a correlated sequence. In contrast, the new convergence results incorporate moment conditions and covariance function decay rate conditions. The ease with which these results can be applied in many cases is illustrated.