# Stopping Times for Stochastic Approximation Procedures* 

Robert L. Sielken, Jr.

## 1. Introduction

The pioneering paper in the field of stochastic approximation was published in 1951 by Robbins and Monro [6]. That paper dealt with the following situation. Suppose that, for every point $x$ belonging to the real line, a random variable $Y(x)$ can be observed. The distribution function of $Y(x)$ and the expected value of $Y(x)$, denoted by $M(x)$ and assumed to exist, are both unknown. Assuming that the equation $M(x)=\alpha$ has a unique root, denoted by $x=\theta$, it is desired to estimate $\theta$ by making observations on $Y$ at points $x_{1}, x_{2}, x_{3}, \ldots$ which are generated sequentially in accordance with some definite experimental procedure in such a way that $x_{n} \rightarrow \theta$ in probability as $n \rightarrow \infty$.

The Robbins-Monro procedure (RM) for generating the sequence $\left\{x_{n}\right\}$ is to take $x_{1}$ to be any constant and define $x_{2}, x_{3}, \ldots$ in accordance with the recurrence relation

$$
\begin{equation*}
x_{n+1}=x_{n}+a_{n}\left(\alpha-y_{n}\right) \tag{1.1}
\end{equation*}
$$

where $y_{n}$ is a random variable whose conditional distribution for given $x_{n}$ coincides with the distribution of the random variable $Y\left(x_{n}\right)$ and is independent of $x_{1}, \ldots, x_{n-1}$ and the sequence $\left\{a_{n}\right\}$ is a sequence of positive constants which converge to zero as $n \rightarrow \infty$.

Robbins and Monro [6] established sufficient conditions for $x_{n} \rightarrow \theta$ in probability as $n \rightarrow \infty$. Later researchers gave results on the rate of convergence, convergence with probability one, convergence in mean square, and asymptotic normality of the sequence $x_{n}$ in both RM and various generalizations of RM-see $[5,8]$, or [10] for specific results and references.

In many practical situations it is desirable to terminate the successive approximation procedure when $x_{n}$ is sufficiently close to $\theta$ with high probability. This paper indicates stopping times $N_{d, \gamma}$ and $T_{d, \gamma}$ which terminate RM in such a way that, for any given $\gamma$ in the open interval $\left(0, \frac{1}{2}\right)$,

$$
\lim _{d \rightarrow 0} P\left(\left|x_{N_{d, \gamma}+1}-\theta\right| \leqq d\right)=1-2 \gamma
$$

and

$$
\lim _{d \rightarrow 0} P\left[\left|M\left(x_{T_{a, \gamma}+1}\right)-\alpha\right| \leqq d\right]=1-2 \gamma .
$$

The empirical behaviors of some of these stopping times have been investigated by the author [8] in a Monte Carlo study.

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## 2. Notation

The following conditions with $Z(x) \equiv Y(x)-M(x)$ are used in all of the theorems in this paper:
$\mathrm{C} 1: \gamma$ is a positive constant less than $\frac{1}{2}$.
C 2 : The sequence $\left\{c_{n}\right\}$ is a sequence of positive constants such that $c_{n} n^{\lambda} \rightarrow c$ as $n \rightarrow \infty, c>0$, and $0<\lambda<\frac{1}{2}$.

C3: The sequence $\left\{a_{n}\right\}$ has the form $\{A / n\}$ where $A$ is a constant such that $2 A \alpha_{1}>1$.

C4: $M$ is a Borel measurable function.
C 5: For each $\varepsilon>0, \inf _{\varepsilon<x-\theta<\varepsilon^{-1}} M(x)-\alpha>0$ and $\sup _{\varepsilon<\theta-x<\varepsilon^{-1}} M(x)-\alpha<0$.
C6: For some constants $K_{1}$ and $K_{2},|M(x)-\alpha| \leqq K_{1}+K_{2}|x-\theta|$ for all $x$.
C7: $\sup _{x} E\left[|Z(x)|^{2}\right] \equiv w<\infty$.
C8: $\lim _{x \rightarrow \theta}^{x} E\left[|Z(x)|^{2}\right]=E\left[|Z(\theta)|^{2}\right] \equiv \sigma^{2}>0$.
C9: $\lim _{R \rightarrow \infty} \lim _{\varepsilon \rightarrow 0^{+}} \sup _{|x-\theta| \leqq e_{i \mid Z(x)}} \int_{\mid>R\}}|Z(x)|^{2} d P=0$.
C10: For some positive constants $g$ and $\alpha_{1}$, if $|x-\theta|<g$, then

$$
M(x)=\alpha+\alpha_{1}(x-\theta)+\delta(x)
$$

where $\delta(x)=o(|x-\theta|)$ as $|x-\theta| \rightarrow 0$.
C11: The distribution function of $Y(x)$, denoted by $H(\cdot \mid x)$, is such that, for every $y, \mathrm{H}\left(\left.y\right|^{\cdot}\right)$ is Borel measurable.

C12: There exists $\varepsilon>0$ such that, for every positive integer $r$,

$$
\sup _{|x-\theta|<\varepsilon} E\left[|Z(x)|^{r}\right]<\infty .
$$

The majority of this paper is concerned with results about the behavior of the sequence $\left\{x_{n}\right\}$. All of these results hold for any initial value $x_{1}$.

The procedures in the next section apply directly to the case in which the random variables $Y(x)$ are such that the $\alpha_{1}$ referred to in condition C10 is positive. If the experimenter is actually observing random variables $Y^{*}(x)$ with a corresponding $\alpha_{1}$ which is negative and he is seeking the value of $x$ such that $E\left[Y^{*}(x)\right]=\alpha^{*}$, then when he carries out the procedures to follow he should let $Y(x)=-Y^{*}(x)$ and $\alpha=-\alpha^{*}$.

## 3. Sequentially Determined Bounded Length Confidence Intervals for RM

Blum [2] gave sufficient conditions, slightly weaker than C3-C7, for the sequence $\left\{x_{n}\right\}$ in RM to be such that

$$
\begin{equation*}
x_{n} \rightarrow \theta \quad \text { as } n \rightarrow \infty \quad \text { w.p. } 1 \tag{3.1}
\end{equation*}
$$

Then Sacks [7] showed that conditions C3-C10 are sufficient for
$n^{\frac{1}{2}}\left(x_{n+1}-\theta\right)$ in RM to be asymptotically normally distributed with mean zero and variance $A^{2} \sigma^{2} /\left(2 A \alpha_{1}-1\right)$.
Furthermore, Burkholder [3] proposed estimators of $\alpha_{1}$ and $\sigma^{2}$ and obtained sufficient conditions for these estimators to converge to $\alpha_{1}$ and $\sigma^{2}$ respectively
with probability one. The form of Burkholder's estimator of $\alpha_{1}$ requires that at the $n$-th step in RM an observation is taken not only on $Y\left(x_{n}\right)$ but also on $Y\left(x_{n}+c_{n}\right)$ where the sequence $\left\{c_{n}\right\}$ is a sequence of positive constants such that $c_{n} n^{\lambda} \rightarrow c$ as $n \rightarrow \infty, c>0$, and $0<\lambda<\frac{1}{2}$. Let the sequence $\left\{y_{n}^{\prime \prime}\right\}$ be a sequence of random variables such that the conditional distribution of $y_{n}^{\prime \prime}$ for given $x_{n}$ coincides with the distribution of $Y\left(x_{n}+c_{n}\right)$ and is independent of $x_{1}, x_{2}, \ldots, x_{n-1}, y_{1}, y_{2}, \ldots, y_{n}, y_{1}^{\prime \prime}, y_{2}^{\prime \prime}, \ldots, y_{n-1}^{\prime \prime}$. Burkholder's results imply that, if conditions C2-C12 are satisfied, then as $n \rightarrow \infty$

$$
\begin{equation*}
t_{n} \equiv \max \left[1 / 2 A, \sum_{j=1}^{n}\left(y_{j}^{\prime \prime}-y_{j}\right) / c_{j} n\right] \rightarrow \alpha_{1} \text { w.p. } 1 \tag{3.3}
\end{equation*}
$$

and

$$
\begin{equation*}
s_{n}^{2} \equiv\left(\frac{1}{2}\right)\left\{\left[\sum_{j=1}^{n}\left(y_{j}^{\prime \prime}-\alpha\right)^{2} / n\right]+\left[\sum_{j=1}^{n}\left(y_{j}-\alpha\right)^{2} / n\right]\right\} \rightarrow \sigma^{2} \text { w.p.1. } \tag{3.4}
\end{equation*}
$$

If, for each $n, n^{\frac{1}{2}}\left(x_{n+1}-\theta\right)$ is normally distributed with mean zero and variance $A^{2} \sigma^{2} /\left(2 A \alpha_{1}-1\right)$ and if $K_{\gamma}$ is the $100(1-\gamma)$-th percentile of a standard normal random variable, then

$$
P\left\{x_{n+1}-K_{\gamma}\left[A^{2} \sigma^{2} / n\left(2 A \alpha_{1}-1\right)\right]^{\frac{1}{2}} \leqq \theta \leqq x_{n+1}+K_{\gamma}\left[A^{2} \sigma^{2} / n\left(2 A \alpha_{1}-1\right)\right]^{\frac{1}{2}}\right\}=1-2 \gamma
$$

Thus, if a $100(1-2 \gamma) \%$ confidence interval on $\theta$ of length $2 d$ is desired, $n$ cold be chosen as the smallest integer such that

$$
d \geqq K_{\gamma}\left[A^{2} \sigma^{2} / n\left(2 A \alpha_{1}-1\right)\right]^{\frac{1}{2}}
$$

or, equivalently,

$$
\begin{equation*}
n \geqq K_{\gamma}^{2} A^{2} \sigma^{2} /\left(2 A \alpha_{1}-1\right) d^{2} \tag{3.5}
\end{equation*}
$$

Thus, when a $100(1-2 \gamma) \%$ confidence interval on $\theta$ of length $2 d$ is desired, the proposed stopping time for RM is $N_{d, \gamma}$ defined as the smallest positive integer $n$ such that

$$
\begin{equation*}
n \geqq K_{\gamma}^{2} A^{2} s_{n}^{2} /\left(2 A t_{n}-1\right) d^{2} \tag{3.6}
\end{equation*}
$$

The principle results concerning this stopping time are:
Theorem 1. If conditions C1-C12 are satisfied, then

$$
\begin{equation*}
\lim _{d \rightarrow 0} N_{d, \gamma} /\left[K_{\gamma}^{2} A^{2} \sigma^{2} /\left(2 A \alpha_{1}-1\right) d^{2}\right]=1 \text { w.p.1. } \tag{3.7}
\end{equation*}
$$

Theorem 2. If conditions C1-C12 are satisfied, then

$$
\begin{equation*}
\lim _{d \rightarrow 0} P\left(\left|x_{N_{d, \gamma}+1}-\theta\right| \leqq d\right)=1-2 \gamma \tag{3.8}
\end{equation*}
$$

If $n_{d, \gamma}=K_{\gamma}^{2} A^{2} \sigma^{2} /\left[\left(2 A \alpha_{1}-1\right) d^{2}\right]$, then $n_{d, \gamma}$ would be the number of observations required to determine a $100(1-2 \gamma) \%$ confidence interval on $\theta$ of length $2 d$ if, for each $n, n^{\frac{1}{2}}\left(x_{n}-\theta\right)$ were exactly normally distributed with mean zero and known variance $A^{2} \sigma^{2} /\left(2 A \alpha_{1}-1\right)$. Thus, Theorem 1 implies that $N_{d, \gamma}$ is an asymptotically efficient stopping time in the sense that $N_{d, \gamma} / n_{d, \gamma} \rightarrow 1$ as $d \rightarrow 0$ with probability one.

If $Y(x)$ is not a continuous random variable for all $x$, then the sequence $\left\{s_{n}^{2}\right\}$ in the definition of $N_{d, \gamma}$ should be replaced by the sequence $\left\{s_{n}^{2^{*}}\right\}$ where

$$
s_{n}^{2^{*}}=\max \left(k_{n}, s_{n}^{2}\right)
$$

and the sequence $\left\{k_{n}\right\}$ is a sequence of positive constants which converge to zero as $n \rightarrow \infty$. This substitution would not affect the results which follow and would only need to be done in order to insure that RM would not be stopped because the estimator of $\sigma^{2}$ was zero.

The experimenter might want to be able to determine stopping times $T_{d, \gamma}$ with the property that, for any given $\gamma$ in $\left(0, \frac{1}{2}\right)$,

$$
\begin{equation*}
\lim _{d \rightarrow 0} P\left[\left|M\left(x_{T_{a, \gamma}+1}\right)-\alpha\right| \leqq d\right]=1-2 \gamma \tag{3.9}
\end{equation*}
$$

Since

$$
n^{\frac{1}{2}}\left[M\left(x_{n}\right)-\alpha\right]=n^{\frac{1}{2}}\left(x_{n}-\theta\right)\left\{\left[M\left(x_{n}\right)-\alpha\right] /\left(x_{n}-\theta\right)\right\}
$$

and condition C10 implies

$$
[M(x)-\alpha] /(x-\theta) \rightarrow \alpha_{1} \quad \text { as } x \rightarrow \theta,
$$

(3.1) and (3.2) imply that, if conditions C3-C10 are satisfied, then $n^{\frac{1}{2}}\left[M\left(x_{n}\right)-\alpha\right]$ is asymptotically normally distributed with mean zero and variance

$$
\alpha_{1}^{2} A^{2} \sigma^{2} /\left(2 A \alpha_{1}-1\right)
$$

Thus, since conditions C2-C12 are sufficient for $t_{n} \rightarrow \alpha_{1}$ and $s_{n}^{2} \rightarrow \sigma^{2}$ as $n \rightarrow \infty$ with probability one, an intuitively appealing stopping time is $T_{d, \gamma}$ defined as the smallest positive integer $n$ such that

$$
K_{\gamma}\left[t_{n}^{2} A^{2} s_{n}^{2} /\left(2 A t_{n}-1\right) n\right]^{\frac{1}{2}} \leqq d
$$

or, equivalently,

$$
\begin{equation*}
n \geqq K_{\gamma}^{2} A^{2} s_{n}^{2} /\left(2 A t_{n}-1\right)\left(d / t_{n}\right)^{2} \tag{3.10}
\end{equation*}
$$

In fact, under conditions C1-C12

$$
\begin{align*}
\lim _{d \rightarrow 0} P & {\left[\left|M\left(x_{T_{d, \gamma}+1}\right)-\alpha\right| \leqq d\right] } \\
& =\lim _{d \rightarrow 0} P\left[\left|\alpha_{1}\left(x_{T_{d, \gamma}+1}-\theta\right)+o\left(\left|x_{T_{d, \gamma}+1}-\theta\right|\right)\right| \leqq d\right]  \tag{3.11}\\
& =\lim _{d \rightarrow 0} P\left(\left|x_{T_{d, \gamma}+1}-\theta\right| \leqq \frac{d}{\alpha_{1}}\right)=1-2 \gamma
\end{align*}
$$

and with probability one

$$
\begin{equation*}
\lim _{d \rightarrow 0} T_{d, \gamma} /\left[K_{\gamma}^{2} \alpha_{1}^{2} \sigma^{2} A^{2} /\left(2 A \alpha_{1}-1\right) d^{2}\right]=1 \tag{3.12}
\end{equation*}
$$

The proofs of (3.11) and (3.12) are similar to the proofs of Theorem1 and Theorem 2.

## 4. Proofs of Theorems 1 and 2

The proof of Theorem 1 follows easily from the definition of $N_{d, \gamma}$ and the convergence of $s_{n}^{2}$ and $t_{n}$ to $\sigma^{2}$ and $\alpha_{1}$ respectively.

The key aspect of the martingale theory which is used in proving Theorem 2 is a generalization of Kolmogorov's inequality for martingales which was proved by Csőrgő [4]. Csőrgö's result is Lemma 1.

Lemma 1. Let the sequence $\left\{W_{k}\right\}$ be a sequence of random variables such that $E\left(W_{1}\right)=0$ and, for $n \geqq 2, E\left(W_{n} \mid W_{1}, \ldots, W_{n-1}\right)=0$. Let $S_{n}=W_{1}+W_{2}+\cdots+W_{n}$ for
each $n \geqq 1$. If $E\left(S_{n}^{2}\right)<\infty$ for each $n \geqq 1$ and the sequence $\left\{g_{n}\right\}$ is a nonincreasing sequence of positive constants, then, for any positive integers $m$ and $n$ with $m<n$ and arbitrary $\varepsilon>0$,

$$
\begin{aligned}
P\left(\max _{m \leqq k \leqq n} g_{k}\left|S_{k}\right| \geqq \varepsilon\right) & \leqq \varepsilon^{-2}\left[\sum_{k=m}^{n-1}\left(g_{k}^{2}-g_{k+1}^{2}\right) E\left(S_{k}^{2}\right)+g_{n}^{2} E\left(S_{n}^{2}\right)\right] \\
& =\varepsilon^{-2}\left[g_{m}^{2} \sum_{k=1}^{m} E\left(W_{k}^{2}\right)+\sum_{k=m+1}^{n} g_{k}^{2} E\left(W_{k}^{2}\right)\right]
\end{aligned}
$$

The proof of Theorem 2 is as follows: Let $N=N_{d, \gamma}$. Since the conditions of Theorems 1 and 2 are the same, it follows from Theorem 1 that

$$
\begin{equation*}
\lim _{d \rightarrow 0} d /\left[K_{\gamma}^{2} A^{2} \sigma^{2} /\left(2 A \alpha_{1}-1\right) N\right]^{\frac{1}{2}}=1 \text { w.p. } 1 \tag{4.1}
\end{equation*}
$$

The definition of $K_{\gamma}$ and (4.1) imply that to prove Theorem 2 it is sufficient to show that, for all $x$,

$$
\begin{equation*}
\lim _{d \rightarrow 0} P\left[N^{\frac{1}{2}}\left(2 A \alpha_{1}-1\right)^{\frac{1}{2}} A^{-1} \sigma^{-1}\left(x_{N+1}-\theta\right) \leqq x\right]=\Phi(x) \tag{4.2}
\end{equation*}
$$

where $\Phi$ is the distribution function of a standard normal random variable.
Since Sacks [7] has shown that, for all $x$,

$$
\begin{equation*}
\lim _{n \rightarrow \infty} P\left[n^{\frac{1}{2}}\left(2 A \alpha_{1}-1\right)^{\frac{1}{2}} A^{-1} \sigma^{-1}\left(x_{n+1}-\theta\right) \leqq x\right]=\Phi(x), \tag{4.3}
\end{equation*}
$$

it need only be shown that this conclusion is valid when $n$ is replaced by $N$.
Let $M_{n}=M\left(x_{n}\right), Z_{n}=Y\left(x_{n}\right)-M\left(x_{n}\right)$, and $\delta_{n}=M_{n}-\alpha-\alpha_{1}\left(x_{n}-\theta\right)$. Thus

$$
\begin{align*}
x_{n+1}-\theta & =\left(x_{n}-\theta\right)-(A / n)\left(Z_{n}+M_{n}-\alpha\right) \\
& =\left[1-\left(A \alpha_{1} / n\right)\right]\left(x_{n}-\theta\right)-(A / n) Z_{n}-(A / n) \delta_{n} \tag{4.4}
\end{align*}
$$

If

$$
\begin{aligned}
B_{m n} & =\prod_{j=m+1}^{n}\left(1-A \alpha_{1} j^{-1}\right) & & \text { for } 0 \leqq m<n \\
& =1 & & \text { for } m=n
\end{aligned}
$$

then the iteration of (4.4) back to $n=1$ implies

$$
\begin{equation*}
x_{n+1}-\theta=B_{0 n}\left(x_{1}-\theta\right)-A \sum_{m=1}^{n} B_{m n} m^{-1} \delta_{m}-\left(\sigma / \alpha_{1}\right) S_{n} \tag{4.5}
\end{equation*}
$$

where

$$
S_{n}=\left(A \alpha_{1} / \sigma\right) \sum_{m=1}^{n} B_{m n} m^{-1} Z_{m}
$$

Let

$$
h_{n}=\left[\sum_{m=1}^{n}\left(A \alpha_{1}\right)^{2} m^{-2} B_{m n}^{2}\right]^{-\frac{1}{2}}
$$

It is easily shown that there exists a sequence of constants $\left\{\varepsilon_{m}\right\}$ such that $\varepsilon_{m} \rightarrow 0$ as $m \rightarrow \infty$ and, for all $n \geqq m$,

$$
\begin{equation*}
\left(1-\varepsilon_{m}\right) m^{A \alpha_{1}} n^{-A \alpha_{1}} \leqq B_{m n} \leqq\left(1+\varepsilon_{m}\right) m^{A \alpha_{1}} n^{-A \alpha_{1}} \tag{4.6}
\end{equation*}
$$

Thus

$$
\begin{equation*}
\lim _{n \rightarrow \infty} h_{n} / n^{\frac{1}{2}}\left(2 A \alpha_{1}-1\right)^{\frac{1}{2}}\left(A \alpha_{1}\right)^{-1}=1 \tag{4.7}
\end{equation*}
$$

Sacks establishes (4.3) by showing that

$$
\begin{align*}
h_{n} B_{0_{n}}\left(x_{1}-\theta\right) \rightarrow 0 & \text { as } n \rightarrow \infty,  \tag{4.8}\\
h_{n} \sum_{m=1}^{n} B_{m n} m^{-1}\left|\delta_{m}\right| \rightarrow 0 & \text { in probability as } n \rightarrow \infty, \tag{4.9}
\end{align*}
$$

and

$$
\begin{equation*}
P\left(h_{n} S_{n} \leqq x\right) \rightarrow \Phi(x) \quad \text { as } n \rightarrow \infty \text { for all } x \tag{4.10}
\end{equation*}
$$

Since (3.7) implies

$$
\begin{equation*}
N \rightarrow \infty \quad \text { as } \quad d \rightarrow 0 \text { w.p. } 1 \tag{4.11}
\end{equation*}
$$

(4.5) and (4.8) imply that (4.2) will hold if

$$
\begin{equation*}
h_{N} \sum_{m=1}^{N} B_{m N} m^{-1}\left|\delta_{m}\right| \rightarrow 0 \quad \text { in probability as } d \rightarrow 0 \tag{4.12}
\end{equation*}
$$

and, for each $x$,

$$
\begin{equation*}
P\left(h_{N} S_{N} \leqq x\right) \rightarrow \Phi(x) \quad \text { as } d \rightarrow 0 \tag{4.13}
\end{equation*}
$$

If $n_{d}=K_{\gamma}^{2} A^{2} \sigma^{2} /\left[\left(2 A \alpha_{1}-1\right) d^{2}\right]$, then (3.7) implies

$$
\begin{equation*}
\lim _{d \rightarrow 0} N / n_{d}=1 \quad \text { w.p. } 1 \tag{4.14}
\end{equation*}
$$

Thus, (4.6) and (4.7) imply (4.12) will hold if

$$
\begin{equation*}
n_{\frac{1}{d}-A \alpha_{1}}^{n} \sum_{m=1}^{N} m^{A \alpha_{1}-1}\left|\delta_{m}\right| \rightarrow 0 \quad \text { in probability as } d \rightarrow 0 \tag{4.15}
\end{equation*}
$$

Now (4.6), (4.7), and (4.9) imply

$$
\begin{equation*}
n^{\frac{1}{2}-A \alpha_{1}} \sum_{m=1}^{N} m^{A \alpha_{1}-1}\left|\delta_{m}\right| \rightarrow 0 \quad \text { in probability as } n \rightarrow \infty \tag{4.16}
\end{equation*}
$$

Thus, since for any $\varepsilon>0$

$$
\begin{aligned}
& P\left(n_{d .}^{\frac{1}{2}-A \alpha_{1}} \sum_{m=1}^{N} m^{A \alpha_{1}-1}\left|\delta_{m}\right|>\varepsilon\right) \\
& \quad \leqq P\left(n_{d .}^{\frac{1}{2}-A \alpha_{1}} \sum_{m=1}^{2 n_{d}} m^{A \alpha_{1}-1}\left|\delta_{m}\right|>\varepsilon \text { and } N / n_{d} \leqq 2\right)+P\left(N / n_{d}>2\right) \\
& \quad \leqq P\left[\left(2 n_{d}\right)^{\frac{1}{2}-A \alpha_{1}} \cdot \sum_{m=1}^{2 n_{d}} m^{A \alpha_{1}-1}\left|\delta_{m}\right|>\varepsilon 2^{\frac{1}{2}-A \alpha_{1}}\right]+P\left(N / n_{d}>2\right),
\end{aligned}
$$

(4.14) and (4.16) imply (4.15) and consequently (4.12).

Anscombe's central limit theorem [1] will imply (4.13) if the sequence $\left\{S_{n}\right\}$ is uniformly continuous in probability; that is, given any small positive $\varepsilon$ and $\eta$,
there are positive constants $v$ and $\rho$ such that $n>v$ implies

$$
\begin{equation*}
P\left(\max _{i:|i-n| \leqq n \rho}\left|S_{i}-S_{n}\right| \geqq \varepsilon h_{n}^{-1}\right)<\eta \tag{4.17}
\end{equation*}
$$

If $\varepsilon^{\prime}=\varepsilon \sigma / h_{n} A \alpha_{1}$ and $m_{1}$ and $m_{2}$ are integers such that $n(1-\rho) \leqq m_{1}<n(1-\rho)+1$ and $n(1+\rho)-1<m_{2} \leqq n(1+\rho)$, then

$$
\begin{align*}
& P\left(\max _{i:|i-n| \leqq n \rho}\left|S_{i}-S_{n}\right| \geqq \varepsilon h_{n}^{-1}\right) \leqq P\left(\max _{m_{1} \leqq i \leqq n-1}\left|\sum_{m=1}^{i} m^{-1} B_{m i} Z_{m}-\sum_{m=1}^{n} m^{-1} B_{m n} Z_{m}\right| \geqq \varepsilon^{\prime}\right) \\
& \quad+P\left(\max _{n+1 \leqq i \leqq m_{2}}\left|\sum_{m=1}^{i} m^{-1} B_{m i} Z_{m}-\sum_{m=1}^{n} m^{-1} B_{m n} Z_{m}\right| \geqq \varepsilon^{\prime}\right) \\
& \quad \leqq \\
& \quad P\left[\max _{m 1}\left|\sum_{m=1}^{i}\left(B_{m i}-B_{m n}\right) m^{-1} Z_{m}\right| \geqq \varepsilon^{\prime} / 2\right]  \tag{4.18}\\
& \quad+P\left(\max _{m_{1} \leqq i \leqq n-1}\left|\sum_{m=i+1}^{n} m^{-1} B_{m n} Z_{m}\right| \geqq \varepsilon^{\prime} / 2\right) \\
& \quad+P\left[\max _{n+1 \leqq i \leqq m_{2}}\left|\sum_{m=1}^{n}\left(B_{m i}-B_{m n}\right) m^{-1} Z_{m}\right| \geqq \varepsilon^{\prime} / 2\right] \\
& \quad+P\left(\max _{n+1 \leqq i \leqq m_{2}}\left|\sum_{m=n+1}^{i} m^{-1} B_{m i} Z_{m}\right| \geqq \varepsilon^{\prime} / 2\right) .
\end{align*}
$$

The establishement of (4.17), which will complete the proof of Theorem 2, is accomplished by proving that, for $n$ sufficiently large, each of the last four probabilities in (4.18) converge to zero as $\rho \rightarrow 0$. For example, in the first of these last four probabilities, if

$$
\begin{array}{rlrl}
D_{m m_{1}} & =B_{m m_{1}} & \text { for } m \leqq m_{1} \\
& =B_{m_{1} m}^{-1} & & \text { for } m>m_{1},
\end{array}
$$

then it can be readily verified that

$$
B_{m i}-B_{m n}=\left(1-B_{i n}\right) B_{m_{1} i} D_{m m_{1}}
$$

so that

$$
\begin{align*}
& P\left[\max _{m_{1} \leqq i \leqq n-1}\left|\sum_{m=1}^{i}\left(B_{m i}-B_{m n}\right) m^{-1} Z_{m}\right| \geqq \varepsilon^{\prime} / 2\right] \\
& \quad=P\left[\max _{m_{1} \leqq i \leqq n-1}\left(1-B_{i n}\right) B_{m_{1} i}\left|\sum_{m=1}^{i} D_{m m_{1}} m^{-1} Z_{m}\right| \geqq \varepsilon^{\prime} / 2\right]  \tag{4.19}\\
& \quad \leqq\left(\varepsilon^{\prime} / 2\right)^{-2}\left[\left(1-B_{m_{1} n}\right)^{2} B_{m_{1} m_{1}}^{2} \sum_{m=1}^{m_{1}} m^{-2} D_{m m_{1}}^{2} E\left(Z_{m}^{2}\right)\right. \\
& \left.\quad+\sum_{m=m_{1}+1}^{n-1}\left(1-B_{m n}\right)^{2} B_{m_{1} m}^{2} D_{m m_{1}}^{2} m^{-2} E\left(Z_{m}^{2}\right)\right]
\end{align*}
$$

where the inequality in (4.19) follows from Lemma 1 . The convergence of the upper bound in (4.19) to zero as $\rho \rightarrow 0$ is a straightforward consequence of condition C7, (4.6), and (4.7).

## 5. Two Modifications of RM

An unfortunate characteristic of the use of $N_{d, \gamma}$ as a stopping time for RM is that, although two observations must be taken at each step in order to obtain convenient information of $\alpha_{1}$, only one of those two observations is actually used to generate the next estimate of $\theta$. One way to avoid this shortcoming is to generate the sequence $\left\{x_{n}\right\}$ by the following procedure which is herein called the modified Robbins-Monro procedure: Let $x_{1}$ be any constant, and let

$$
x_{n+1}=x_{n}+a_{n}\left[\alpha-\left(\frac{1}{2}\right)\left(y_{n}^{\prime}+y_{n}^{\prime \prime}\right)\right]
$$

in which $y_{n}^{\prime}$ and $y_{n}^{\prime \prime}$ are random variables whose conditional distributions for given $x_{n}$ are independent, coincide with the distributions of $Y\left(x_{n}-c_{n}\right)$ and $Y\left(x_{n}+c_{n}\right)$ respectively, and are independent of $x_{1}, x_{2}, \ldots, x_{n-1}, y_{1}^{\prime}, y_{2}^{\prime}, \ldots, y_{n-1}^{\prime}, y_{1}^{\prime \prime}, y_{2}^{\prime \prime}, \ldots, y_{n-1}^{\prime}$.

Venter [9] actually went one step further than the modified Robbins-Monro procedure. He proposed a successive approximation procedure which at the $n$-th step incorporates the cumulative information on $\alpha_{1}$ as well as the observations on $Y\left(x_{n}+c_{n}\right)$ and $Y\left(x_{n}-c_{n}\right)$ into the recurrence relation which generates $x_{n+1}$. If the conditions C3-C10 are satisfied, the asymptotic variance of the error in RM, $A^{2} \sigma^{2} /\left(2 A \alpha_{1}-1\right) n$, is minimized with respect to $A$ if $A=\alpha_{1}^{-1}$. Unfortunately, $\alpha_{1}$ will usually be unknown. Venter recognized these facts and proposed that the sequence $\left\{x_{n}\right\}$ be generated according to the recurrence relation

$$
x_{n+1}=x_{n}+\left(1 / n t_{n}^{*}\right) \cdot\left[\alpha-\left(\frac{1}{2}\right)\left(y_{n}^{\prime \prime}+y_{n}^{\prime}\right)\right]
$$

where $t_{n}^{*}$ is an estimator of $\alpha_{1}, x_{1}$ is any constant, and the form of the sequences $\left\{y_{n}^{\prime \prime}\right\}$ and $\left\{y_{n}^{\prime}\right\}$ is the same as in the modified Robbins-Monro procedure.

Stopping times analogous to $N_{d, \gamma}$ and $T_{d, \gamma}$ have been proposed in [8] for both the modified Robbins-Monro procedure and Venter's procedure. Statements analogous to Theorems 1 and 2, (3.11), and (3.12) are given there for the stopping times of these two additional procedures. Also, in [8] are the results of a fairly extensive Monte Carlo study of the empirical behaviors of the stopping times for all three procedures.

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Robert L.Sielken, Jr. Institute of Statistes Texas A \& M University College Station, Texas 77843 USA

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