# Gradient algorithms for quadratic optimization <br> with fast convergence rates 

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#### Abstract

We propose a family of gradient algorithms for minimizing a quadratic function $f(x)=$ $(A x, x) / 2-(x, y)$ in $\mathbb{R}^{d}$ or a Hilbert space, with simple rules for choosing the step-size at each iteration. We show that when the step-sizes are generated by a dynamical system with ergodic distribution having the arcsine density on a subinterval of the spectrum of $A$, the asymptotic rate of convergence of the algorithm can approach the (tight) bound on the rate of convergence of a conjugate gradient algorithm stopped before $d$ iterations, with $d \leq \infty$ the space dimension.


Key words: Chebyshev polynomials, conjugate gradient, Krylov space, logistic map, quadratic operator, steepest descent.

## 1 Introduction

Consider the problem of minimizing a quadratic function $f(\cdot)$ defined either on $\mathbb{R}^{d}$ or a Hilbert space by

$$
\begin{equation*}
f(x)=\frac{1}{2}(A x, x)-(x, y), \tag{1}
\end{equation*}
$$

where $(\cdot, \cdot)$ denotes the inner product. We assume that $A$ is either a symmetric positive-definite matrix or a self-adjoint operator, with

$$
0<m=\inf _{(x, x)=1}(A x, x)<M=\sup _{(x, x)=1}(A x, x)<\infty
$$

If $A$ is a matrix, then $m$ and $M$ are the smallest and largest eigenvalues of $A$, respectively.
Consider a general gradient algorithm with iterations of the form

$$
\begin{equation*}
x_{k+1}=x_{k}-\gamma_{k} g_{k}, k=0,1,2 \ldots \tag{2}
\end{equation*}
$$

where $g_{k}=\nabla f\left(x_{k}\right)$ is the gradient of the objective function $f(\cdot)$ at point $x_{k}$. For the objective function (1), $\nabla f(x)=A x-y$. The iteration (2) can be rewritten in terms of the gradients as

$$
\begin{equation*}
g_{k+1}=g_{k}-\gamma_{k} A g_{k} \tag{3}
\end{equation*}
$$

In a series of papers [10, 11, 12] and the monograph [9] many gradient algorithms have been shown to be equivalent to special algorithms for updating measures on the interval $[m, M]$. The central idea is that of renormalization applied to the gradient. For simplicity the presentation is made for the finite dimensional case where $A$ is a matrix, which can be assumed, without loss of generality, to be diagonal

[^0]$A=\operatorname{diag}\left(\lambda_{1}, \ldots, \lambda_{d}\right)$ with eigenvalues $m=\lambda_{1} \leq \lambda_{2} \leq \cdots \leq \lambda_{d}=M$. Extension to the Hilbert-space case will be considered in Section 5.

Write $z_{k}=g_{k} / \sqrt{\left(g_{k}, g_{k}\right)}$ for the normalized gradient at $x_{k}$ and define

$$
p_{i}^{(k)}=\left\{z_{k}\right\}_{i}^{2}=\frac{\left\{g_{k}\right\}_{i}^{2}}{\sum_{j=1}^{d}\left\{g_{k}\right\}_{j}^{2}}, \quad i=1, \ldots, d
$$

as the $i$-th probability corresponding to vector $z_{k}$, where $\{v\}_{i}$ denotes the $i$-th component of vector $v$. Let $\nu_{k}$ denote the probability measure on the spectrum of $A$ defined by the $p_{i}^{(k)}$ 's, that is, $\nu_{k}\left(\lambda_{i}\right)=p_{i}^{(k)}$. The probability measure $\nu_{k+1}$ is defined by

$$
p_{i}^{(k+1)}=\frac{\left\{g_{k+1}\right\}_{i}^{2}}{\left(g_{k+1}, g_{k+1}\right)} \quad \text { for } i=1, \ldots, d
$$

Note that (3) gives

$$
\begin{equation*}
\left(g_{k+1}, g_{k+1}\right)=\left(g_{k}, g_{k}\right)-2 \gamma_{k}\left(A g_{k}, g_{k}\right)+\gamma_{k}^{2}\left(A^{2} g_{k}, g_{k}\right), \tag{4}
\end{equation*}
$$

so that

$$
\begin{equation*}
p_{i}^{(k+1)}=\frac{\left(1-\gamma_{k} \lambda_{i}\right)^{2}}{\left(g_{k}, g_{k}\right)-2 \gamma_{k}\left(A g_{k}, g_{k}\right)+\gamma_{k}^{2}\left(A^{2} g_{k}, g_{k}\right)}\left\{g_{k}\right\}_{i}^{2}=\frac{\left(1-\gamma_{k} \lambda_{i}\right)^{2}}{1-2 \gamma_{k} \mu_{1}^{(k)}+\gamma_{k}^{2} \mu_{2}^{(k)}} p_{i}^{(k)}, \tag{5}
\end{equation*}
$$

where

$$
\begin{equation*}
\mu_{\alpha}^{(k)}=\mu_{\alpha}\left(\nu_{k}\right)=\frac{\left(A^{\alpha} g_{k}, g_{k}\right)}{\left(g_{k}, g_{k}\right)} \tag{6}
\end{equation*}
$$

is the $\alpha$-th moment of the measure $\nu_{k}$. When two eigenvalues of $A$ are equal, say $\lambda_{j}=\lambda_{j+1}$, the updating rules for $p_{j}^{(k)}$ and $p_{j+1}^{(k)}$ are identical so that the analysis of the behaviour of the algorithm remains the same when $p_{j}^{(k)}$ and $p_{j+1}^{(k)}$ are confounded. We may thus assume that all eigenvalues of $A$ are distinct. Also, a zero weight remains equal to zero at all subsequent iterations, we thus assume that $\nu_{0}\left(\lambda_{i}\right)>0$ for all $i$.

A common definition for the rate of convergence of the algorithm at iteration $k$ is $r_{k}=\left(g_{k+1}, g_{k+1}\right) /\left(g_{k}, g_{k}\right)$. The rate for $n$ iterations is

$$
\prod_{k=0}^{n-1} r_{k}=\frac{\left(g_{n}, g_{n}\right)}{\left(g_{0}, g_{0}\right)}
$$

therefore, the asymptotic rate of the algorithm can naturally be defined as

$$
\begin{equation*}
R=\lim _{n \rightarrow \infty} R_{n}, \quad \text { with } \quad R_{n}=\left(\prod_{k=0}^{n-1} r_{k}\right)^{1 / n} \tag{7}
\end{equation*}
$$

Of course, this rate may depend on the initial point $x_{0}$ or, equivalently, on $g_{0}$. Other rates which are asymptotically equivalent to $\left\{r_{k}\right\}$ can be considered as well, see [12] and Remark 6.

The most familiar gradient algorithm is the steepest-descent algorithm, for which the step-size $\gamma_{k}$ at iteration $k$ is chosen so as to minimize $f\left(x_{k}-\gamma g_{k}\right)$ with respect to $\gamma$, which gives $\gamma_{k}=\left(g_{k}, g_{k}\right) /\left(A g_{k}, g_{k}\right)=$ $1 / \mu_{1}^{(k)}$. Its asymptotic behaviour is well-known, see [1, 10]. In particular, its convergence is slow: the asymptotic rate $R$ depends on the starting point but is never far from its worst value given by the Kantorovich bound

$$
R_{\max }=\left(\frac{\rho-1}{\rho+1}\right)^{2}
$$

where $\rho=M / m$, the condition number of $A$. The asymptotic behaviour of the family of algorithms defined by $\gamma_{k}=\mu_{\alpha}^{(k)} / \mu_{\alpha+1}^{(k)}$ (which includes the method of minimum residues for $\alpha=1$ ) is shown in [12] to be similar.

Obtaining a faster asymptotic rate of convergence for gradient algorithms requires to extend the possible choices for the step-size $\gamma_{k}$. Rewrite the updating rule (5) in terms of iteration applied to the probability measure $\nu_{k}$,

$$
\begin{equation*}
\nu_{k+1}(\lambda)=\frac{\left(1-\gamma_{k} \lambda\right)^{2}}{1-2 \gamma_{k} \mu_{1}^{(k)}+\gamma_{k}^{2} \mu_{2}^{(k)}} \nu_{k}(\lambda)=\frac{\left(\lambda-\beta_{k}\right)^{2}}{\beta_{k}^{2}-2 \beta_{k} \mu_{1}^{(k)}+\mu_{2}^{(k)}} \nu_{k}(\lambda) \tag{8}
\end{equation*}
$$

where $\beta_{k}=1 / \gamma_{k}$ and $\nu_{k}(\lambda)$ is the weight assigned by the measure $\nu_{k}$ to the point $\lambda$. The roots $\beta_{k}$ in (8) are the key control variables for a gradient algorithm. Different strategies for choosing $\beta_{k}$ give different families of algorithms. Note that the only information about $\nu_{k}$ one has access to corresponds to its moments $\mu_{\alpha}^{(k)}, \alpha=1,2 \ldots$ Many of the examples of algorithms presented in [6], with $\beta_{k}$ a function of $\mu_{1}^{(k)}$ and $\mu_{2}^{(k)}$, exhibit a much faster asymptotic rate of convergence than $R_{\max }$ (it seems that allowing $\beta_{k}$ to depend on more moments $\mu_{\alpha}^{(k)}$ does not yield further improvement in the rate of convergence). Fast convergence (small $R$ ) is observed for algorithms that exhibit a chaotic-type behaviour in $\mathbb{R}^{d}$, which makes their theoretical study difficult. The same is true for some algorithms for which $\beta_{k}$ is allowed to depend on moments of several previous measures $\nu_{k-i}, i=1, \ldots, u$. For instance, in the Barzilai-Borwein algorithm [2], $\beta_{k}$ is either $\mu_{1}^{(k-1)}$ or $\mu_{2}^{(k-1)} / \mu_{1}^{(k-1)}$.

Conjugate gradient, $s$-step optimal, MINRES and other algorithms based on Krylov spaces do not use gradient directions for their successive iterations, see, e.g., [8]. However, when analyzing their behaviour, one can construct an equivalent sequence of iterations following the gradient directions with control variables $\beta_{k}$ depending on $k$ and on moments of previous measures $\nu_{k-i}, i=0,1,2 \ldots$ The conjugate gradient algorithm in $\mathbb{R}^{d}$ converges in $d$ iterations. When $d$ is large, preserving the conjugacy of successive directions is difficult and restarting the algorithm after each sequence of $s$ iterations is recommended. This corresponds to the $s$-step optimal gradient algorithm, see [5, 13], which does not have finite convergence but whose guaranteed asymptotic rate of convergence is

$$
\begin{equation*}
R_{s}^{*}=\left(\frac{R_{\infty}^{s / 2}+R_{\infty}^{-s / 2}}{2}\right)^{-2 / s}=T_{s}^{-2 / s}\left(\frac{\rho+1}{\rho-1}\right) \tag{9}
\end{equation*}
$$

where

$$
R_{\infty}=\lim _{s \rightarrow \infty} R_{s}^{*}=\left(\frac{\sqrt{\rho}-1}{\sqrt{\rho}+1}\right)^{2}
$$

and $T_{s}(\cdot)$ is the $s$-th Chebyshev polynomial:

$$
T_{s}(t)=\cos [s \arccos (t)]=\frac{\left(t+\sqrt{t^{2}-1}\right)^{s}+\left(t-\sqrt{t^{2}-1}\right)^{s}}{2}
$$

In this paper we propose a family of gradient algorithms based on simple rules for choosing the sequence of control variables $\beta_{k}$. The main idea is to force $\nu_{k}\left(\lambda_{j}\right), j=2, \ldots, d-1$, to tend to zero as $k \rightarrow \infty$. The measure $\nu_{k}$, which summarizes the state of the iterates at step $k$, is then almost fully characterized by $\nu_{k}(m)$, which facilitates the analysis of the asymptotic behaviour. Furthermore, we show that the sequence $\left\{\beta_{k}\right\}$ can be chosen independently of $\left\{\nu_{k}\right\}$ while ensuring that the asymptotic rate of convergence is arbitrarily close to $R_{\infty}$. This independence of $\left\{\beta_{k}\right\}$ on $\left\{\nu_{k}\right\}$ makes the algorithms at the same time simple and robust with respect to the precision of calculations. Also, the step-sizes $\gamma_{k}=1 / \beta_{k}$, $k=1,2 \ldots$ are simpler to calculate than those of the steepest-descent algorithm. Convergence rates close
to $R_{\infty}$ are obtained when the $\beta_{k}$ 's are constructed so that their asymptotic distribution is close to a distribution with the arcsine density.

The worst-case rate $R_{s}^{*}$ can be reached for the $s$-step optimal gradient when $d>s$, in the sense that there exist eigenvalues $\lambda_{i}$ and initial point $x_{0}$ for the algorithm such that the rate of convergence after $s$ iterations is exactly $R_{s}^{*}$ (and the behavior in terms of renormalized gradient $z_{k}$ is then periodic with period $s$ ), see [5, 13]. The same is true for the conjugate gradient algorithm: for $s<d$ there exist eigenvalues $\lambda_{i}$ and a starting point $x_{0}$ such that the convergence rate after $s$ iterations is exactly $R_{s}^{*}$.

If $d$ is large (relative to the total number of iterations), $s$ is not very large and the eigenvalues of $A$ are well-spread in the spectral interval $[m, M]$, then the actual rates (per one matrix-vector multiplication) of the MINRES and other optimal methods based on the use of $s$-dimensional Krylov spaces are very close to $R_{s}^{*}$ and are often larger than $R_{\infty}$. Bearing in mind that the asymptotic rates of the algorithms suggested below can be arbitrarily close to $R_{\infty}$ and these algorithms are extremely simple and robust, these algorithms may be preferable to MINRES and other Krylov space based methods for large-scale quadratic optimization problems.

The paper is organized as follows. In Section 2 we show that for a suitable choice of the sequence $\left\{\beta_{k}\right\}$ the algorithm attracts to the plane spanned by the eigenvectors associated with $\lambda_{1}=m$ and $\lambda_{d}=M$. In Section 3, we assume that the values of $m$ and $M$ are known and give the expression of the asymptotic rate of convergence of the algorithm in the case where the $\beta_{k}$ 's are generated by pairs symmetric with respect to $(m+M) / 2$. Several examples are presented, some with a rate arbitrarily close to $R_{\infty}$. The case where $m$ and $M$ are unknown is considered in Section 4 where a practical algorithm is suggested and some simulation results are presented. Finally, the infinite dimensional situation where $f(\cdot)$ is defined on a Hilbert space is considered Section 5.

## 2 Attraction of the sequence $\left\{\nu_{k}\right\}$ to the set of measures supported at $m$ and $M$

Theorem 1 Assume that $\beta_{k}>0, \beta_{k} \notin\{m, M\}$ for all $k$ and that the sequence $\left\{\beta_{k}\right\}$ has asymptotic distribution function $F(\beta)$ which is supported on an interval $\left[m^{\prime}, M^{\prime}\right]$ with $0<m^{\prime} \leq M^{\prime}<\infty$. Suppose, moreover, that the limiting distribution satisfies

$$
\begin{equation*}
\int \log (\beta-\lambda)^{2} d F(\beta)<\max \left\{\int \log (M-\beta)^{2} d F(\beta), \int \log (\beta-m)^{2} d F(\beta)\right\}, \quad \forall \lambda \in\left\{\lambda_{2}, \ldots, \lambda_{d-1}\right\} \tag{10}
\end{equation*}
$$

Then, the gradient algorithm associated with the sequence $\left\{\beta_{k}\right\}$ is such that $\lim _{k \rightarrow \infty} \nu_{k}\left(\lambda_{i}\right)=0$ for all $i=2, \ldots, d-1$. Furthermore, there exist constants $C>0, k_{0}>0$ and $0 \leq \theta<1$ such that

$$
\begin{equation*}
\sum_{i=2}^{d-1} \nu_{k}\left(\lambda_{i}\right) \leq C \theta^{k} \text { for } k>k_{0} \tag{11}
\end{equation*}
$$

Proof. The fact that the sequence $\left\{\beta_{k}\right\}$ has asymptotic distribution function $F(\beta)$ implies

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \frac{1}{k} \sum_{j=0}^{k-1} h\left(\beta_{j}\right)=\int h(\beta) d F(\beta) \tag{12}
\end{equation*}
$$

for any continuous function $h(\cdot)$ such that $\int|h(\beta)| d F(\beta)<\infty$, see [7]. Define

$$
\begin{equation*}
H_{k}(\lambda)=C_{k}\left(\lambda-\beta_{0}\right)^{2}\left(\lambda-\beta_{1}\right)^{2} \cdots\left(\lambda-\beta_{k-1}\right)^{2}, \tag{13}
\end{equation*}
$$

with $C_{k}$ a normalizing constant such that $\nu_{k}(\lambda)=H_{k}(\lambda) \nu_{0}(\lambda)$ in (8), and assume that

$$
\begin{equation*}
\int \log (M-\beta)^{2} d F(\beta) \leq \int \log (\beta-m)^{2} d F(\beta) \tag{14}
\end{equation*}
$$

(if this inequality is not met, $m$ should be replaced with $M$ in all considerations below). Define the sum

$$
\begin{equation*}
S_{k}(\lambda, m)=\frac{1}{k} \log \frac{H_{k}(\lambda)}{H_{k}(m)}=-\frac{1}{k} \sum_{j=0}^{k-1} \log \left(\beta_{j}-m\right)^{2}+\frac{1}{k} \sum_{j=0}^{k-1} \log \left(\lambda-\beta_{j}\right)^{2} \tag{15}
\end{equation*}
$$

and consider the first sum $I_{k}(m)=(1 / k) \sum_{j=0}^{k-1} \log \left(\beta_{j}-m\right)^{2}$ in the right-hand side of (15) and the related integral $I(m)=\int \log (\beta-m)^{2} d F(\beta)$. Since the c.d.f. $F(\cdot)$ is supported on a bounded interval $\left[m^{\prime}, M^{\prime}\right]$ we have $I(m)<\infty$. The assumptions (10) and (14) imply $I(m)>-\infty$ and the property (12) then gives the convergence $I_{k}(m) \rightarrow I(m)$ as $k \rightarrow \infty$.

Consider now the second sum $I_{k}(\lambda)=(1 / k) \sum_{j=0}^{k-1} \log \left(\beta_{j}-\lambda\right)^{2}$ in the right-hand side of (15) and the related integral $I(\lambda)=\int \log (\beta-\lambda)^{2} d F(\beta)$. Since the c.d.f. $F(\cdot)$ is supported on a bounded interval, the integral $I(\lambda)$ is properly defined but may equal $-\infty$ (for example, if the c.d.f. $F(\cdot)$ has a discontinuity at the point $\lambda$ ). If $I(\lambda)=-\infty$ then as $k \rightarrow \infty$ the sum $I_{k}(\lambda)$ tends to $-\infty$ too. If $I(\lambda)>-\infty$ then either $I_{k}(\lambda)=-\infty$ for all $k$ large enough (when at least one $\beta_{j}$ is equal to $\lambda$ ) or (12) implies that $I_{k}(\lambda)$ tends to $I(\lambda)$ as $k \rightarrow \infty$.

Therefore, from (10), $S_{k}(\lambda, m)$ tends to a negative value (possibly $-\infty$ ) as $k \rightarrow \infty$. This implies that there exists $k_{0} \geq 0$ and $\delta>0$ such that for all $k \geq k_{0}$ and $\lambda \in\left\{\lambda_{2}, \ldots, \lambda_{d-1}\right\}$

$$
\begin{equation*}
S_{k}(\lambda, m)=\frac{1}{k} \log \frac{H_{k}(\lambda)}{H_{k}(m)} \leq-\delta \tag{16}
\end{equation*}
$$

that is, $H_{k}(\lambda) / H_{k}(m) \leq \theta^{k}$, where $\theta=\exp (-\delta)<1$. This yields $\sum_{i=2}^{d-1} \nu_{k}\left(\lambda_{i}\right) \leq \theta^{k}\left(\sum_{i=2}^{d-1} \nu_{0}\left(\lambda_{i}\right)\right) / \nu_{0}(m)$ for $k>k_{0}$, hence (11). The result $\lim _{k \rightarrow \infty} \nu_{k}\left(\lambda_{i}\right)=0$ for $i=2, \ldots, d-1$ obviously follows from (11).

Remark 1 The sequence $\left\{\beta_{k}\right\}$ can be assumed random, for instance formed by independent and identically distributed random variables. In this case, all the statements are true with probability one. When the $\beta_{k}$ 's are simply independent, with $\left\{F_{k}\right\}$ the sequence of corresponding distribution functions and $(1 / k) \sum_{j=0}^{k-1} F_{j}$ converging weakly to $F$ as $k$ tends to infinity, one may refer to [3, Th. 2.5.3, p. 36] for a property similar to (12).

Remark 2 Typically, the spectrum of $A$ is unknown. In that case, the condition (10) can be replaced with the more restrictive one

$$
\begin{equation*}
\int \log (\beta-\lambda)^{2} d F(\beta)<\max \left\{\int \log (M-\beta)^{2} d F(\beta), \int \log (\beta-m)^{2} d F(\beta)\right\}, \quad \forall \lambda \in(m, M) \tag{17}
\end{equation*}
$$

Remark 3 If the distribution with c.d.f. $F(\cdot)$ is symmetric with respect to $(m+M) / 2$, then we have $\int \log (M-\beta)^{2} d F(\beta)=\int \log (\beta-m)^{2} d F(\beta)$ and therefore the condition (17) simplifies to

$$
\begin{equation*}
\int \log (\beta-\lambda)^{2} d F(\beta)<\int \log (\beta-m)^{2} d F(\beta), \quad \forall \lambda \in(m, M) \tag{18}
\end{equation*}
$$

Remark 4 Note that the support $\left[m^{\prime}, M^{\prime}\right]$ of the distribution with c.d.f. $F(\cdot)$ could be different from $[m, M]$ and does not have to be a subset of $[m, M]$.

Remark 5 The results of Theorem 1 also apply when $\beta_{k}$ depends on the moments of previous measures $\nu_{k-i}, i=0,1,2 \ldots$

Example 1 For the steepest-descent algorithm with $\beta_{k}=\mu_{1}^{(k)}$, the limiting measure for $\left\{\beta_{k}\right\}$ is the two-point measure assigning weights $1 / 2$ at $z$ and $m+M-z$ for some $z \in(m, M)$. The condition (17) then simply expresses the property that two successive iterations (8) of the algorithm asymptotically give a larger increase of the weights at the endpoints $m$ and $M$ than at any other point in the interval $(m, M)$; that is,

$$
\begin{equation*}
(z-m)^{2}(M-z)^{2}>(z-\lambda)^{2}(m+M-z-\lambda)^{2}, \forall \lambda \in(m, M) \tag{19}
\end{equation*}
$$

Since for all $z$ the only maximum of $(z-\lambda)^{2}(m+M-z-\lambda)^{2}$ with respect to $\lambda \in(m, M)$ is at $\lambda^{*}=(m+M) / 2$, the inequality (19) can be rewritten as $(z-m)^{2}(M-z)^{2}>\left(z-\lambda^{*}\right)^{2}\left(m+M-z-\lambda^{*}\right)^{2}$, which gives

$$
\begin{equation*}
z \in\left(\frac{1}{2}(m+M)-\frac{1}{2 \sqrt{2}}(M-m), \frac{1}{2}(m+M)+\frac{1}{2 \sqrt{2}}(M-m)\right) . \tag{20}
\end{equation*}
$$

This corresponds to the definition of the stability interval for the attractor in [10, 12]. A similar result holds for all gradient-type algorithms from the family considered in [12].

Example 2 If we choose $\beta_{k}=\sqrt{\mu_{2}^{(k)}}$, then the limiting measure for $\left\{\beta_{k}\right\}$ is the delta-measure concentrated at the point $\lambda^{*}=(m+M) / 2$; as a consequence, the asymptotic rate for the related gradient algorithm is $R_{\max }$. Proof of these facts can be found in [4] and [6], Sect. 2.7.

## 3 Asymptotic rate for symmetrically placed control variables

### 3.1 Main result

Theorem 2 Assume that the conditions of Theorem 1 are satisfied and that, moreover, the control variables $\beta_{k}$ are generated by symmetric pairs for large $k$; that is, $\beta_{2 j+1}=M+m-\beta_{2 j}$ for all $j \geq j_{0}$, with $\beta_{2 j} \in[m+\varepsilon, M-\varepsilon]$ for some $\varepsilon \in(0,(M-m) / 2)$. Then, the asymptotic rate $R$ satisfies

$$
\begin{equation*}
\log R=\int \log \left|\frac{(M-\beta)(\beta-m)}{\beta(m+M-\beta)}\right| d F(\beta)=\int \log \frac{(\beta-m)^{2}}{\beta^{2}} d F(\beta) \tag{21}
\end{equation*}
$$

Proof. First note that dividing (4) through by $\left(g_{k}, g_{k}\right)$ gives the following expression for the rate $r_{k}$,

$$
\begin{equation*}
r_{k}=1-2 \gamma_{k} \frac{\left(A g_{k}, g_{k}\right)}{\left(g_{k}, g_{k}\right)}+\gamma_{k}^{2} \frac{\left(A^{2} g_{k}, g_{k}\right)}{\left(g_{k}, g_{k}\right)}=1-2 \mu_{1}^{(k)} / \beta_{k}+\mu_{2}^{(k)} / \beta_{k}^{2} \tag{22}
\end{equation*}
$$

Consider a measure $\nu$ with weights $p$ and $1-p$ at $m$ and $M$ respectively, $0<p<1$. Apply two successive iterations (8) with control parameters $\beta$ and $\beta^{\prime}=m+M-\beta$ to this measure. The product of the two successive rates does not depend on $p$ and is equal to $R_{2}^{2}(\beta)=(M-\beta)^{2}(\beta-m)^{2} /[\beta(m+M-\beta)]^{2}$.

According to Theorem $1, \nu_{k}$ tends to be supported at $m$ and $M$ and the rate of convergence is exponential. We thus obtain for two successive iterations with control variables $\beta_{2 j}$ and $\beta_{2 j+1}=m+$ $M-\beta_{2 j}$

$$
R_{2}^{2}\left(\beta_{2 j}\right)\left[1-\frac{A \theta^{2 j}}{R_{2}^{2}\left(\beta_{2 j}\right)}\right]<r_{2 j} r_{2 j+1}<R_{2}^{2}\left(\beta_{2 j}\right)\left[1+\frac{A \theta^{2 j}}{R_{2}^{2}\left(\beta_{2 j}\right)}\right]
$$

for some $A>0$ and $j>k_{0} / 2$, see Theorem 1. Since $\beta_{2 j} \in[m+\varepsilon, M-\varepsilon]$, we have $R_{2}^{2}\left(\beta_{2 j}\right) \geq R_{2}^{2}(m+\varepsilon)=$ $\varepsilon(M-m-\varepsilon) /[(m+\varepsilon)(M-\varepsilon)]>0$. Therefore,

$$
\log R_{2}\left(\beta_{2 j}\right)-B \theta^{2 j}<\log \sqrt{r_{2 j} r_{2 j+1}}<\log R_{2}\left(\beta_{2 j}\right)+B \theta^{2 j}
$$



Figure 1: $R_{2}(\beta)$ for $m=1, M=4$
with $B=A / R_{2}^{2}(m+\varepsilon)$, for $j$ large enough. Since $\sum_{j=0}^{\infty} \theta^{2 j}=1 /\left(1-\theta^{2}\right)<\infty$, we obtain from (12),

$$
\log R=\lim _{k \rightarrow \infty} \frac{1}{k} \sum_{j=0}^{k-1} \log \sqrt{r_{2 j} r_{2 j+1}}=\int \log R_{2}(\beta) d F(\beta),
$$

hence the first expression in (21). The second expression follows from the fact that the c.d.f. $F(\cdot)$ is symmetric with respect to $(m+M) / 2$.

Example 3 Uniform density. Let the distribution with c.d.f. $F(\cdot)$ be uniform with density $p(\beta)=$ $1 /\left(M^{\prime}-m^{\prime}\right), \beta \in\left[m^{\prime}, M^{\prime}\right]$, with $m^{\prime}=m+\varepsilon, M^{\prime}=M-\varepsilon$ and $0<\varepsilon<(M-m) / 2$. Then the asymptotic rate of convergence is

$$
\begin{equation*}
R_{\text {uniform }, \varepsilon}=\exp \left\{\frac{1}{M^{\prime}-m^{\prime}} \int_{m^{\prime}}^{M^{\prime}} \log \frac{\left(\beta-m^{\prime}\right)^{2}}{\beta^{2}} d \beta\right\}=\left(M^{\prime}-m^{\prime}\right)^{2} \exp \left\{-2 \frac{M^{\prime} \log M^{\prime}-m^{\prime} \log m^{\prime}}{M^{\prime}-m^{\prime}}\right\} \tag{23}
\end{equation*}
$$

Remark 6 One can easily check that the result stated in Theorem 2 holds for other definitions for the rate of convergence, see, e.g., [12, Th. 6]. For instance, the rate

$$
r_{k}^{\prime}=\frac{f\left(x_{k+1}\right)-f^{*}}{f\left(x_{k}\right)-f^{*}}=\frac{\left(A^{-1} g_{k+1}, g_{k+1}\right)}{\left(A^{-1} g_{k}, g_{k}\right)}
$$

where $f^{*}=\min _{x} f(x)$, can be written as

$$
\begin{equation*}
r_{k}^{\prime}=1-2 /\left(\mu_{-1}^{(k)} \beta_{k}\right)+\mu_{1}^{(k)} /\left(\mu_{-1}^{(k)} \beta_{k}^{2}\right) \tag{24}
\end{equation*}
$$

and the corresponding asymptotic rate $R^{\prime}=\lim _{n \rightarrow \infty}\left(\prod_{k=0}^{n-1} r_{k}^{\prime}\right)^{1 / n}$ is equal to $R$ which can be computed by (21).

Remark 7 The shape of $R_{2}(\beta)$ as a function of $\beta$ shows that fast convergence is obtained for $\beta$ close to $m$ or $M$, see Figure 1, hence the interest of taking $\varepsilon$ small in Theorem 2.

Remark 8 When $\nu_{k}$ is a two-point measure supported at $m$ and $M$, two iterations of (8) with $\beta_{k+1}=$ $M+m-\beta_{k}$ give $\nu_{k+2}=\nu_{k}$. Under the conditions of Theorem 2 the measure $\nu_{k}$ thus converges to a
measure $\bar{\nu}_{k}=p_{k} \delta_{m}+\left(1-p_{k}\right) \delta_{M}$ supported at $m$ and $M$, with $p_{2 j}$ tending to a constant $p_{\infty}$ as $j$ tends to infinity. The limiting distribution of the sequence $\left\{p_{2 j+1}\right\}$ depends on $F(\cdot)$ and $p_{\infty}$, while the value of $p_{\infty}$ depends on the initial measure $\nu_{0}$ and the spectrum of $A$.

### 3.2 Finite collection of control variables

Assume that the points $\beta_{0}, \beta_{1} \ldots$ are generated in repeated groups $B=\left\{\beta_{0}, \ldots, \beta_{N}\right\}$ of $N+1$ points in $(m, M), N \geq 0$. Additionally, the points in $B$ are symmetric with respect to $(m+M) / 2$. We may always assume that $\beta_{0} \leq \ldots \leq \beta_{N}$. In this case, if $N$ is even then $\beta_{N / 2}=(m+M) / 2$. The condition (18) now becomes

$$
\begin{equation*}
\sum_{j=0}^{N} \log \left(\beta_{j}-\lambda\right)^{2}<\sum_{j=0}^{N} \log \left(\beta_{j}-m\right)^{2}, \quad \forall \lambda \in(m, M) \tag{25}
\end{equation*}
$$

If this condition is met then the asymptotic rate is

$$
\begin{equation*}
R=R_{N}=\left[\prod_{j=0}^{N} \frac{\left(\beta_{j}-m\right)^{2}}{\beta_{j}^{2}}\right]^{1 /(N+1)} \tag{26}
\end{equation*}
$$

Example 4 Uniform grid. Assume that for some integer $N \geq 0$,

$$
\begin{equation*}
B=\left\{\beta_{0}, \ldots, \beta_{N}\right\} \text { with } \beta_{i}=m+\frac{i+\frac{1}{2}}{N+1}(M-m), \quad i=0,1, \ldots, N \tag{27}
\end{equation*}
$$

It is easy to see that the condition (25) is met. The rate $R_{N}$ computed by (26) is given by

$$
R_{N}=\left(\frac{\Gamma^{2}(N+3 / 2) \Gamma^{2}\left(\frac{m+M+2 N m}{2(M-m)}\right)}{\pi \Gamma^{2}\left(\frac{2 N M+3 M-m}{2(M-m)}\right)}\right)^{1 /(N+1)}
$$

where $\Gamma(\cdot)$ is the gamma-function. The value of $R_{N}$ for $m=1, M=4$ is plotted in Figure 2 as a function of $N$. Asymptotically, as $N \rightarrow \infty, R_{N}$ approaches $R_{\text {uniform, } 0}$ defined in (23) (with $R_{\text {uniform, } 0} \simeq 0.2232$ for $m=1, M=4$ ). Instead of using the $\beta_{i}$ 's according to (27) for large $N$, one can generate the sequence $\left\{\beta_{i}\right\}$ using, for example, the Bernoulli shift:

$$
\begin{equation*}
H_{B}(t)=2 t[\bmod 1], t \in(0,1), \tag{28}
\end{equation*}
$$

with $\beta_{0}$ randomly chosen in $\left(m^{\prime}, M^{\prime}\right)$, and for all $j=0,1,2 \ldots$

$$
\beta_{2 j+1}=M^{\prime}+m^{\prime}-\beta_{2 j}, \beta_{2 j+2}=m^{\prime}+\left(M^{\prime}-m^{\prime}\right) H_{B}\left(\frac{\beta_{2 j}-m^{\prime}}{M^{\prime}-m^{\prime}}\right)
$$

with $m^{\prime}=m+\varepsilon, M^{\prime}=M-\varepsilon$ and $0<\varepsilon<(M-m) / 2$.
Example 5 Nearly optimal $N+1$ points. Consider first the case $N=1$. When the condition (25) is satisfied, the asymptotic rate is $R_{2}=|M-\beta||\beta-m| /[\beta|m+M-\beta|]$, see the proof of Theorem 2 . For $\beta \in[m, M], R_{2}$ improves when $|\beta-(m+M) / 2|$ increases and reaches its minimum value, zero, at $\beta \in\{m, M\}$, see Remark 7. Condition (25) imposes that $\beta$ belongs to the interval (20), by choosing $\beta$ sufficiently close to $(m+M) / 2 \pm(M-m) /(2 \sqrt{2})$ one makes the rate arbitrarily close to $R_{2}^{*}$, with $R_{s}^{*}$ defined by (9).


Figure 2: Asymptotic rate of convergence (26) for $m=1$ and $M=4$ when the $\beta_{j}$ 's are on the uniform grid (stars) and when they correspond to Chebyshev points (triangles, $\varepsilon=10^{-6}$ )

Take now $N=2$, with $\beta_{0}=\beta, \beta_{1}=(m+M) / 2$ and $\beta_{2}=m+M-\beta$. Similarly to the previous case, condition (25) imposes that $\beta$ belongs to the interval $((m+M) / 2-\sqrt{3}(M-m) / 4,(m+M) / 2+$ $\sqrt{3}(M-m) / 4)$, with the rate $R_{3}$ getting close to $R_{3}^{*}$ for $\beta$ close to $(m+M) / 2 \pm \sqrt{3}(M-m) / 4$.

By induction, one can show that the rate $R_{N}$ can be made arbitrarily close to the value $R_{s}^{*}$ defined by (9), with $s=N+1$, when the $N+1$ points $\beta_{i}$ are suitably chosen and are constructed from the roots of Chebyshev polynomials. This construction is considered in the next example. (Note that the fact that $R_{N}$ can be made arbitrarily close to $R_{N+1}^{*}$ is not a coincidence: the worst case analysis of the $s$-step optimal gradient algorithm, which yields the rate $R_{s}^{*}$, corresponds to the situation where the $\beta_{i}$ 's are rescaled roots of the $s$-th order Chebyshev polynomial, see [5].)

Example 6 Chebyshev points. Chebyshev points are defined by

$$
t_{k}=\cos \left(\frac{\pi}{2} \frac{2 k+1}{N+1}\right), \quad k=0, \ldots, N
$$

and correspond to the roots of $T_{N+1}(x)=\cos ((N+1) \arccos x)$, the Chebyshev polynomials of the first kind. These points are symmetric on $(-1,1)$. The asymptotic density of the points $\left\{t_{k}\right\}_{0}^{N}$, as $N \rightarrow \infty$, is $p(t)=1 /\left(\pi \sqrt{1-t^{2}}\right), t \in(-1,1)$.

Define

$$
\beta_{k}=\frac{m+M}{2}+\frac{M-m-2 \varepsilon}{2} t_{k}, k=0, \ldots, N
$$

where $0<\varepsilon<(M-m) / 2$. These points belong to the interval $(m+\varepsilon, M-\varepsilon)$ and are symmetric with respect to $(m+M) / 2$. As $\varepsilon>0$, the condition (25) holds. The rate $R_{N}$ computed by (26) is plotted in Figure 2 as a function of $N$ for $m=1, M=4$ and $\varepsilon=10^{-6}$. Asymptotically, as $N \rightarrow \infty, R_{N}$ approaches $R_{\text {arcsin, } \varepsilon}$ defined below in (31).

### 3.3 Control variables with arcsine density on a subinterval of $[m, M]$

Let us assume that the distribution with c.d.f. $F(\cdot)$ has the density

$$
\begin{equation*}
p_{\varepsilon}(\beta)=\frac{1}{\pi \sqrt{\left(\beta-m^{\prime}\right)\left(M^{\prime}-\beta\right)}}, \quad m^{\prime} \leq \beta \leq M^{\prime} \tag{29}
\end{equation*}
$$

where $m^{\prime}=m+\varepsilon, M^{\prime}=M-\varepsilon$ and $0<\varepsilon<(M-m) / 2$. The density $(29)$ is called the arcsine density on the interval $\left[m^{\prime}, M^{\prime}\right]$.

The sequence of points $\left\{\beta_{i}\right\}$ can be generated using, for example, the logistic map

$$
\begin{equation*}
H_{L}(x)=4 x(1-x), x \in(0,1), \tag{30}
\end{equation*}
$$

with $\beta_{0}$ randomly chosen in $\left(m^{\prime}, M^{\prime}\right)$, and for all $j=0,1,2 \ldots$

$$
\beta_{2 j+1}=M^{\prime}+m^{\prime}-\beta_{2 j}, \beta_{2 j+2}=m^{\prime}+\left(M^{\prime}-m^{\prime}\right) H_{L}\left(\frac{\beta_{2 j}-m^{\prime}}{M^{\prime}-m^{\prime}}\right) .
$$

Note that the control variables $\beta_{j}$ are placed symmetrically in the interval $[m, M]$. We show below that the condition (17) holds for each $\varepsilon>0$. According to (21), the asymptotic rate of convergence is then

$$
\begin{equation*}
R_{\arcsin , \varepsilon}=\exp \left\{\int_{m^{\prime}}^{M^{\prime}} \log \frac{(\beta-m)^{2}}{\beta^{2}} p_{\varepsilon}(\beta) d \beta\right\} \tag{31}
\end{equation*}
$$

and we show below that

$$
\begin{equation*}
R_{\arcsin , \varepsilon}=\left(\frac{M-m+2 \sqrt{\varepsilon(M-m-\varepsilon)}}{M+m+2 \sqrt{(M-\varepsilon)(m+\varepsilon)}}\right)^{2} \tag{32}
\end{equation*}
$$

For $\varepsilon=0$ this gives $R_{\arcsin , 0}=R_{\infty}=(\sqrt{\rho}-1)^{2} /(\sqrt{\rho}+1)^{2}$ where $\rho=M / m$. However, we cannot choose $\varepsilon=0$ as the condition (17) does not hold (we also show below that $I(\lambda)=\int \log (\beta-\lambda)^{2} d F(\beta)=$ $2 \log (M-m)-4 \log 2$ for $\lambda \in[m, M])$. Since the condition does hold for any $\varepsilon>0$, the rate of the algorithm can be made arbitrarily close to $R_{\infty}$ : for small $\varepsilon>0$, we have

$$
R_{\arcsin , \varepsilon}=R_{\infty}(1+4 \sqrt{\varepsilon(M-m)})+O(\varepsilon), \quad \varepsilon \rightarrow 0
$$

The rest of this section is devoted to the verification of (17) for $\varepsilon>0$ and to the derivation of the formula (32) for the rate $R_{\arcsin , \varepsilon}$. Define the integral

$$
\begin{equation*}
J\left(z, m^{\prime}, M^{\prime}\right)=\int_{m^{\prime}}^{M^{\prime}} \frac{\log (\beta-z)^{2}}{\pi \sqrt{\left(\beta-m^{\prime}\right)\left(M^{\prime}-\beta\right)}} d \beta \tag{33}
\end{equation*}
$$

where $-\infty<z<\infty$. The changes of variables $t=-1+2\left(\beta-m^{\prime}\right) /\left(M^{\prime}-m^{\prime}\right)$ and $x=-1+2(z-$ $\left.m^{\prime}\right) /\left(M^{\prime}-m^{\prime}\right)$ in the integral (33) give

$$
\begin{equation*}
J\left(z, m^{\prime}, M^{\prime}\right)=2 \log \frac{M^{\prime}-m^{\prime}}{2}+\frac{1}{\pi} I_{x}, \text { where } I_{x}=\int_{-1}^{1} \frac{\log (t-x)^{2}}{\sqrt{1-t^{2}}} d t \tag{34}
\end{equation*}
$$

Assume first that $|x| \leq 1$. By changing the variable $t=\cos \phi$ in the integral $I_{x}$, we obtain

$$
I_{x}=\int_{0}^{\pi} \frac{\log (\cos \phi-x)^{2}}{\sin \phi} \sin \phi d \phi=\int_{0}^{\pi} \log (\cos \phi-x)^{2} d \phi
$$

As $\cos (\phi)=\cos (2 \pi-\phi) \forall \phi$, we have $\int_{0}^{\pi} \log (\cos \phi-x)^{2} d \phi=\int_{\pi}^{2 \pi} \log (\cos \phi-x)^{2} d \phi$, which implies $I_{x}=\frac{1}{2} \int_{0}^{2 \pi} \log (\cos \phi-x)^{2} d \phi$. As we assume $-1 \leq x \leq 1$ we can set $\psi=\arccos x$ (so that $x=\cos \psi$ ). Using now the identity $\cos \phi-\cos \psi=2 \sin \frac{\psi-\phi}{2} \sin \frac{\phi+\psi}{2}$, we obtain

$$
\begin{aligned}
I_{x} & =\frac{1}{2} \int_{0}^{2 \pi} \log (\cos \phi-\cos \psi)^{2} d \phi=\frac{1}{2} \int_{0}^{2 \pi} \log \left(2 \sin \frac{\phi-\psi}{2} \sin \frac{\phi+\psi}{2}\right)^{2} d \phi \\
& =\frac{1}{2}\left[\int_{0}^{2 \pi} 2 \log 2 d \phi+\int_{0}^{2 \pi} \log \left(\sin \frac{\phi-\psi}{2}\right)^{2} d \phi+\int_{0}^{2 \pi} \log \left(\sin \frac{\phi+\psi}{2}\right)^{2} d \phi\right] \\
& =2 \pi \log 2+\left[\int_{0}^{\pi} \log \left(\sin ^{2}(\phi-\psi / 2)\right) d \phi+\int_{0}^{\pi} \log \left(\sin ^{2}(\phi+\psi / 2)\right) d \phi\right]
\end{aligned}
$$

The function $t \rightarrow \sin ^{2} t$ is $\pi$-periodic and therefore for any $\psi^{\prime}$ we get

$$
\int_{0}^{\pi} \log \left(\sin ^{2}\left(\phi+\psi^{\prime}\right)\right) d \phi=\int_{0}^{\pi} \log \left(\sin ^{2}(\phi)\right) d \phi=2 \int_{0}^{\pi} \log (\sin \phi) d \phi
$$

This implies

$$
\begin{equation*}
I_{x}=2 \pi \log 2+4 \int_{0}^{\pi} \log (\sin \phi) d \phi=2 \pi \log 2-4 \pi \log 2=-2 \pi \log 2, \quad \forall x \in[-1,1] \tag{35}
\end{equation*}
$$

Assume now that $|x| \geq 1$. From (35) we have $I_{1}=-2 \pi \log 2$ and differentiating $I_{x}$ we get

$$
I_{x}^{\prime}=\left(\int_{-1}^{1} \frac{\log (x-t)^{2}}{\sqrt{1-t^{2}}} d t\right)^{\prime}=\frac{2 \pi}{\sqrt{x^{2}-1}}
$$

Therefore, for $x>1$,

$$
\begin{equation*}
I_{-x}=I_{x}=I_{1}+\int_{1}^{x} I_{z}^{\prime} d z=-2 \pi \log 2+\int_{1}^{x} \frac{2 \pi}{\sqrt{z^{2}-1}} d z=-2 \pi \log 2+2 \pi \log \left(\frac{x+\sqrt{x^{2}-1}}{2}\right) \tag{36}
\end{equation*}
$$

Combining (35) and (36) we obtain

$$
I_{x}=\int_{-1}^{1} \frac{\log (t-x)^{2}}{\sqrt{1-t^{2}}} d t= \begin{cases}-2 \pi \log 2 & \text { if }|x| \leq 1 \\ 2 \pi \log \left(|x|+\sqrt{x^{2}-1}\right)-2 \pi \log 2 & \text { if }|x| \geq 1\end{cases}
$$

together with (34), it gives

$$
J\left(z, m^{\prime}, M^{\prime}\right)= \begin{cases}2 \log \left(M^{\prime}-m^{\prime}\right)-4 \log 2 & \text { if } m^{\prime} \leq z \leq M^{\prime}  \tag{37}\\ 2 \log \left(M^{\prime}-m^{\prime}\right)+2 \log \left(\left|t_{z}\right|+\sqrt{t_{z}^{2}-1}\right)-4 \log 2 & \text { if } z<m^{\prime} \text { or } z>M^{\prime}\end{cases}
$$

where $t_{z}=-1+2\left(z-m^{\prime}\right) /\left(M^{\prime}-m^{\prime}\right)$. Therefore, $J\left(\lambda, m^{\prime}, M^{\prime}\right)<J\left(m, m^{\prime}, M^{\prime}\right)=J\left(M, m^{\prime}, M^{\prime}\right)$ for all $\lambda$ in $(m, M)$ and (17) is satisfied. The expression (32) for the rate $R_{\text {arcsin, } \varepsilon}$ easily follows from (37) and the representation $R_{\text {arcsin }, \varepsilon}=\exp \left[J\left(m, m^{\prime}, M^{\prime}\right)-J\left(0, m^{\prime}, M^{\prime}\right)\right]$ with $m^{\prime}=m+\varepsilon$ and $M^{\prime}=M-\varepsilon$.

## 4 Estimation of $m, M$ and a practical algorithm

### 4.1 Estimation of $m, M$ and asymptotic behavior in the non symmetric case

The values of $m$ and $M$ can be easily estimated in the first iterations of the algorithm (3), for instance by computing the first moment $\mu_{1}^{(j)}$ for several values of $j=0,1,2 \ldots$ and taking

$$
\begin{equation*}
m_{k}=\min \left\{\mu_{1}^{(j)}, j=0, \ldots, k\right\}, \quad M_{k}=\max \left\{\mu_{1}^{(j)}, j=0, \ldots, k\right\} \tag{38}
\end{equation*}
$$

as estimates. We then necessarily have $m<m_{k}<M_{k}<M$ for $k \geq 1$.
Suppose that the estimation is stopped at some $k_{0}$, that is, $m_{k}=m_{k_{0}}$ and $M_{k}=M_{k_{0}}$ for all $k>k_{0}$. Then, under the conditions of Theorem 1 with $m^{\prime}=m_{k_{0}}$ and $M^{\prime}=M_{k_{0}}$ we have $\sum_{i=2}^{d-1} \nu_{k}\left(\lambda_{i}\right) \leq C \theta^{k}$ for $k$ larger than some $k_{1}$ and constants $C>0$ and $0 \leq \theta<1$. Suppose that the control variables are generated by pairs for $k>k_{0}$, as in Theorem 2 , but with $\beta_{2 k+1}=M_{k_{0}}+m_{k_{0}}-\beta_{2 k}$, for all $k>k_{0}$.

If $M_{k_{0}}+m_{k_{0}}=M+m$, Theorem 2 applies and the asymptotic rate $R$ satisfies (21). For instance, if the $\beta_{k}$ 's are generated as in Section 3.3, and have the arcsine density on $\left[m_{k_{0}}, M_{k_{0}}\right]$, the asymptotic rate is $R_{\text {arcsin, } \varepsilon}$ with $\varepsilon=m_{k_{0}}-m=M-M_{k_{0}}$. Consider now the standard situation where $M_{k_{0}}+m_{k_{0}} \neq M+m$
and suppose that $M-M_{k_{0}}>m_{k_{0}}-m$. The asymptotic distribution of the $\beta_{k}$ 's, symmetric in $\left[m_{k_{0}}, M_{k_{0}}\right]$, is then biased towards $m$ and $\nu_{k}(m)$ tends to zero when $k \rightarrow \infty$. Following the same line as in the proof of Theorem 2, we obtain that the product of rates at two successive iterations for the delta measure at $M$, with control parameters respectively $\beta$ and $\beta^{\prime}=M_{k_{0}}+m_{k_{0}}-\beta$, is $R_{2}^{2}=(M-\beta)^{2}\left(M-\beta^{\prime}\right)^{2} /\left(\beta \beta^{\prime}\right)^{2}$. The asymptotic rate then satisfies

$$
\log R=\int \log \left|\frac{(M-\beta)\left(M+\beta-M_{k_{0}}-m_{k_{0}}\right)}{\beta\left(M_{k_{0}}+m_{k_{0}}-\beta\right)}\right| d F(\beta) .
$$

Similarly, supposing that $M-M_{k_{0}}<m_{k_{0}}-m$ gives an asymptotic distribution of the $\beta_{k}$ 's biased towards $M$, so that $\nu_{k}(m)$ tends to 1 as $k \rightarrow \infty$, and the asymptotic rate satisfies

$$
\log R=\int \log \left|\frac{(\beta-m)\left(M_{k_{0}}+m_{k_{0}}-\beta-m\right)}{\beta\left(M_{k_{0}}+m_{k_{0}}-\beta\right)}\right| d F(\beta)
$$

Now, note that $\nu_{k}(m) \rightarrow 0$ implies that $\mu_{1}^{(k)} \rightarrow M$ and $\nu_{k}(m) \rightarrow 1$ implies that $\mu_{1}^{(k)} \rightarrow m, k \rightarrow \infty$, so that maintaining the adaptation of the estimation of $m_{k}$ and $M_{k}$ by (38) ensures that $m_{k} \rightarrow m$ and $M_{k} \rightarrow M$ as $k \rightarrow \infty$. This permits to recover the same asymptotic rates as Section 3.3, even in situations where $m$ and $M$ are unknown. Since the estimated values $m_{k}$ and $M_{k}$ quickly converge to $m$ and $M$, see for instance Figure 3, we need to generate the control variable $\beta_{k}$ in $\left[m_{k}+\varepsilon, M_{k}-\varepsilon\right]$ at iteration $k$. A practical algorithm is given below.

### 4.2 An algorithm based on the arcsine density

A possible algorithm is then as follows.

- Choose $\tau$ as a small positive number (e.g., $\tau=10^{-6}$ ), set $z_{0}=0$;
- for $k=0,1$, set $\beta_{k}=\mu_{1}^{(k)}$ (steepest-descent) and set $m_{1}=\min \left\{\mu_{1}^{(0)}, \mu_{1}^{(1)}\right\}, M_{1}=\max \left\{\mu_{1}^{(0)}, \mu_{1}^{(1)}\right\}$;
- for $k>1$, set $\varepsilon_{k}=\tau\left(M_{k-1}-m_{k-1}\right)$ and generate the $\beta_{k}$ 's by pairs:
- for $k=2 j$, set $z_{j}=\left\{\varphi+z_{j-1}\right\}$ and $\beta_{2 j}=m_{k}+\varepsilon_{k}+\left(\cos \left(\pi z_{j}\right)+1\right)\left(M_{k}-m_{k}-2 \varepsilon_{k}\right) / 2$, where $\{t\}$ denotes the fractional part of $t$ and $\varphi=(\sqrt{5}-1) / 2 \simeq 0.61803$;
- for $k=2 j+1$, set $\beta_{2 j+1}=M_{k}+m_{k}-\beta_{2 j}$;
set $m_{k}=\min \left\{m_{k-1}, \mu_{1}^{(k)}\right\}, M_{k}=\max \left\{M_{k-1}, \mu_{1}^{(k)}\right\}$.
The sequence $z_{1}, z_{2} \ldots$ is such that $z_{j}=\{j \varphi\}$ so that the sequence is asymptotically uniform on $[0,1]$, see, e.g., [7]. This implies that the asymptotic distribution of the sequence $\beta_{k}$ has the arcsine density on $[m+\varepsilon, M-\varepsilon]$ where $\varepsilon=\tau(M-m)$. From (32), the rate of the algorithm satisfies

$$
\lim _{n \rightarrow \infty} R_{n}=R_{\arcsin , \tau(M-m)}=R_{\infty}(1+4 \sqrt{\tau})+\mathcal{O}(\tau), \tau \rightarrow 0
$$

The dynamical system $z_{j}=\{j \varphi\}$ generates a sequence in $[0,1]$ with much better uniformity characteristics than sequences generated by the Bernoulli shift (28). Since the logistic map (30) corresponds to a transformation of the Bernoulli shift, the construction above, based on $z_{j}=\{j \varphi\}$, produces a sequence of control variables $\beta_{k}$ with better distribution characteristics than sequences generated with (30).

Figures 3, 4 and 5 illustrate the typical behavior of the algorithm above in a large-dimensional badly conditioned problem. In the example presented, $d=1000, m=1, M=\rho=1000$ and the eigenvalues $\lambda_{i}$


Figure 3: Convergence of the estimates $m_{n}$ and $M_{n}$ as functions of $n(m=1, M=1000, d=1000)$


Figure 4: Values $\nu_{n}(m)$ as functions of $n$; circles for $n=2 j$, dots for $n=2 j+1(\rho=1000, d=1000)$


Figure 5: Rate of convergence $R_{n}$, see (7), as a function of $n$; the limiting value $R_{\arcsin , \tau(M-m)}$ is indicated by the dashed line ( $m=1, M=1000, d=1000, \tau=10^{-6}$ )
are random and uniformly distributed on the interval $[m, M]$ (one would obtain exactly the same plots if the eigenvalues were equally-spaced on $[m, M]$ ).

In terms of complexity of calculations, only the multiplications of $d$-dimensional vectors by the $d \times d$ matrix $A$ are expensive. The steepest descent algorithm requires the calculation of $\beta_{k}=\mu_{2}^{(k)} / \mu_{1}^{k)}=$ $\left(A g_{k}, A g_{k}\right) /\left(A g_{k}, g_{k}\right)$ at iteration $k$. Having computed $g_{k}$ and $A g_{k}$, one may notice that next gradient $g_{k+1}$ can be obtained as $g_{k+1}=g_{k}-\left(1 / \beta_{k}\right) A g_{k}$, so that only the computation of $A g_{k+1}$ is expensive at iteration $k+1$. However, a long sequence of iterations of this type may produce an accumulation of rounding errors, and it is rather recommended to recalculate $g_{k+1}$ from $x_{k+1}$ by $g_{k+1}=A x_{k+1}-y$, see (1). This then requires two multiplications by $A$ at each steepest-descent iteration.

In the algorithm above, iteration $k$ only requires the calculation of the gradient $g_{k}=A x_{k}-y$, and thus only one multiplication by $A$. Notice that the estimation of $m_{k}$ and $M_{k}$ through the moments $\mu_{1}^{(j)}=\left(A g_{j}, g_{j}\right) /\left(g_{j}, g_{j}\right)$, see (38), does not require the calculation of $A g_{j}$ at step $k$. Indeed, allowing a delay of one step in the estimation, we have $\left(g_{j}, g_{j+1}\right)=\left(g_{j}, g_{j}-\left(1 / \beta_{j}\right) A g_{j}\right)$ so that $\mu_{1}^{(j)}$ is obtained at next step from

$$
\mu_{1}^{(j)}=\beta_{j}\left[1-\frac{\left(g_{j}, g_{j+1}\right)}{\left(g_{j}, g_{j}\right)}\right]
$$

Also, one may observe in Figure 3 that the convergence of $m_{n}$ and $M_{n}$ to $m$ and $M$ respectively is very fast, so that the estimation can be stopped after a few iterations. On the whole, it makes iterations with the algorithm above about twice simpler than steepest-descent iterations (even when $m$ and $M$ are estimated), with much faster convergence.

## 5 Hilbert space case

In the Hilbert-space case, $A$ is a self-adjoint operator and its spectrum $\mathcal{S}_{A}$ is a closed subset of the interval [ $m, M$ ] of the real line, with $m, M \in \mathcal{S}_{A}$. Let $E_{\lambda}$ be the spectral family associated with $A$ and define the measure $\nu_{k}=d\left(E_{\lambda} z_{k}, z_{k}\right), m \leq \lambda \leq M$, with $z_{k}=g_{k} / \sqrt{\left(g_{k}, g_{k}\right)}$ the normalized gradient at $x_{k}$. We have $\left(z_{k}, z_{k}\right)=1=\int_{m}^{M} \nu_{k}(d \lambda)$ and $\nu_{k}$ is a probability measure on the Borel sets of $(0, \infty)$, satisfying
$\nu_{k}([m, M])=1$ for all $k$ and with moments still defined by (6). One iteration of a gradient algorithm with control variable $\beta_{k}$ thus gives in terms of $\nu_{k}$

$$
\nu_{k+1}(\mathcal{A})=\frac{\int_{\mathcal{A}}\left(\lambda-\beta_{k}\right)^{2} \nu_{k}(d \lambda)}{\beta_{k}^{2}-2 \beta_{k} \mu_{1}^{(k)}+\mu_{2}^{(k)}},
$$

for $\mathcal{A}$ any measurable subset of $[m, M]$, see (8). The properties obtained for the finite dimensional case remain valid and only a few adaptations are required.

Theorem 3 Assume that the sequence $\left\{\beta_{k}\right\}$ has asymptotic distribution function $F(\beta)$ which is supported on an interval $\left[m^{\prime}, M^{\prime}\right]=[m+\varepsilon, M-\varepsilon]$ with $0<\varepsilon<(M-m) / 2$. Suppose, moreover, that $I(\lambda)=$ $\int \log (\beta-\lambda)^{2} d F(\beta)$ is a continuous function of $\lambda$ on $\left(m^{\prime}, M^{\prime}\right)$ and that

$$
\begin{equation*}
I(\lambda)<\max \left\{\int \log (M-\beta)^{2} d F(\beta), \int \log (\beta-m)^{2} d F(\beta)\right\}, \quad \forall \lambda \in\left(m^{\prime}, M^{\prime}\right) \tag{39}
\end{equation*}
$$

and that $\nu_{0}\{[m, m+\gamma)\}>0$ and $\nu_{0}\{(M-\gamma, M]\}>0$ for all $\gamma>0$. Then, the measure $\nu_{k}$ converges to a two-point measure supported at $m$ and $M$, in the sense that there exists $k_{0}$ such that, for any function $g(\lambda)$ continuous on $[m, M]$ and any $\delta>0$, there exists $\gamma>0$ such that
$\max \left\{\left|\int_{m}^{C} g(\lambda) \nu_{k}(d \lambda)-g(m) \int_{m}^{C} \nu_{k}(d \lambda)\right|,\left|\int_{C}^{M} g(\lambda) \nu_{k}(d \lambda)-g(M) \int_{C}^{M} \nu_{k}(d \lambda)\right|\right\}<\delta+C_{\gamma} \alpha_{\gamma}^{k}, k>k_{0}$, where $C=(m+M) / 2$ and $C_{\gamma}>0, \alpha_{\gamma} \in(0,1)$ are constants depending on $\gamma$. If, moreover, the control variables $\beta_{k}$ are generated by symmetric pairs for large $k$, that is, $\beta_{2 j+1}=M+m-\beta_{2 j}$ for all $j \geq j_{0}$, then the asymptotic rate $R$ satisfies (21).

Proof. The proof of convergence of $\nu_{k}$ to a two-point measure follows the same arguments as for Theorem 1. Suppose that $F(\cdot)$ satisfies (14). We still have for the first term of the sum $S_{k}(\lambda, m)$ defined by (15)

$$
I_{k}(m)=\frac{1}{k} \sum_{j=0}^{k-1} \log \left(\beta_{j}-m\right)^{2} \rightarrow I(m)=\int \log (\beta-m)^{2} d F(\beta), k \rightarrow \infty
$$

Concerning the second term $I_{k}(\lambda)=(1 / k) \sum_{j=0}^{k-1} \log \left(\beta_{j}-\lambda\right)^{2}$ we need now a bound uniform in $\lambda$, that is, we need to show that

$$
\begin{equation*}
\forall \epsilon>0, \exists K_{0} \text { such that: } \sup _{\lambda \in\left(m^{\prime}, M^{\prime}\right)} I_{k}(\lambda)-I(\lambda)<\epsilon, \forall k>K_{0} \tag{40}
\end{equation*}
$$

Take a ball $\mathcal{B}\left(\lambda_{1}, \delta\right)=\left\{\lambda:\left|\lambda-\lambda_{1}\right| \leq \delta\right\}$ and consider $\bar{a}_{\delta}(\beta)=\sup _{\lambda \in \mathcal{B}\left(\lambda_{1}, \delta\right)} \log (\beta-\lambda)^{2}$, which is an increasing function of $\delta, \bar{a}_{\delta}(\beta)=2 \log \left(\left|\beta-\lambda_{1}\right|+\delta\right)$. We have

$$
\lim _{\delta \rightarrow 0} \int \bar{a}_{\delta}(\beta) d F(\beta)=\int\left[\lim _{\delta \rightarrow 0} \bar{a}_{\delta}(\beta)\right] d F(\beta)=I\left(\lambda_{1}\right)
$$

and therefore, there exists $\delta_{1}=\delta_{1}\left(\lambda_{1}\right)$ such that $\int \bar{a}_{\delta}(\beta) d F(\beta)<I\left(\lambda_{1}\right)+\epsilon / 3$ for $\delta<\delta_{1}$. Now,

$$
\sup _{\lambda \in \mathcal{B}\left(\lambda_{1}, \delta\right)} I_{k}(\lambda) \leq(1 / k) \sum_{j=0}^{k-1} 2 \log \left(\left|\beta_{j}-\lambda\right|+\delta\right)<\int \bar{a}_{\delta}(\beta) d F(\beta)+\epsilon / 3
$$

for all $k$ larger than some $K_{1}=K_{1}\left(\lambda_{1}, \delta\right)$. Also, from the continuity of $I(\lambda)$, there exists $\delta_{2}=\delta_{2}\left(\lambda_{1}\right)$ such that $\inf _{\lambda \in \mathcal{B}\left(\lambda_{1}, \delta\right)} I(\lambda)>I\left(\lambda_{1}\right)-\epsilon / 3$ for $\delta<\delta_{2}$. Altogether it gives $\sup _{\lambda \in \mathcal{B}\left(\lambda_{1}, \delta\right)} I_{k}(\lambda)-I(\lambda)<\epsilon$ for
$\delta<\delta_{0}\left(\lambda_{1}\right)=\min \left(\delta_{1}, \delta_{2}\right)$ and $k>K_{1}$. It only remains to cover [ $\left.m^{\prime}, M^{\prime}\right]$ with a finite number of such balls $\mathcal{B}\left(\lambda_{i}, \delta\right)$, with $\delta<\min _{i} \delta_{0}\left(\lambda_{i}\right)$ to obtain the result (40). Since $\log (\beta-\lambda)^{2}$ is a decreasing (resp. increasing) function of $\lambda$ in $\left[m, m^{\prime}\right]$ (resp. in $\left[M^{\prime}, M\right]$ ), together with the condition (39) it implies that for any set $\mathcal{S} \subset(m, M), \limsup _{k \rightarrow \infty} \sup _{\lambda \in \mathcal{S}} S_{k}(\lambda, m) \leq-\delta$ for some $\delta=\delta(\mathcal{S})>0$. Therefore, there exists $k_{0}$ such that, $\forall k>k_{0}$, $\sup _{\lambda \in\left(m^{\prime}, M^{\prime}\right)} H_{k}(\lambda) / H_{k}(m) \leq \theta_{\varepsilon}^{k}$ where $\theta_{\varepsilon}=\exp \left(-\delta_{\varepsilon}\right)<1$.

Consider now a function $g(\lambda)$ continuous on $[m, M$ ] and define

$$
\Delta_{k}=\left|\int_{m}^{C} g(\lambda) \nu_{k}(d \lambda)-g(m) \int_{m}^{C} \nu_{k}(d \lambda)\right|
$$

where $C=(m+M) / 2$. We show below that

$$
\begin{equation*}
\forall \delta>0, \exists \gamma>0 \text { such that } \Delta_{k}<\delta+2 \frac{D_{g}}{\int_{m}^{m+\gamma} \nu_{0}(d \lambda)} \alpha_{\gamma}^{k} \text { for all } k>k_{0} \tag{41}
\end{equation*}
$$

for some $\alpha_{\gamma}<1$, where $D_{g}=\max _{\lambda \in[m, C]}|g(\lambda)-g(m)|$. We have $\Delta_{k}<\int_{m}^{C}|g(\lambda)-g(m)| \nu_{k}(d \lambda)=$ $\Delta_{k, 1}+\Delta_{k, 2}+\Delta_{k, 3}$, with
$\Delta_{k, 1}=\int_{m}^{m+2 \gamma}|g(\lambda)-g(m)| \nu_{k}(d \lambda), \Delta_{k, 2}=\int_{m+2 \gamma}^{m^{\prime}}|g(\lambda)-g(m)| \nu_{k}(d \lambda), \Delta_{k, 3}=\int_{m^{\prime}}^{C}|g(\lambda)-g(m)| \nu_{k}(d \lambda)$, $\gamma<\varepsilon / 2$. From the continuity of $g(\lambda)$, we can take $\gamma$ small enough to have $\Delta_{k, 1}<\delta \int_{m}^{m+2 \gamma} \nu_{k}(d \lambda) \leq \delta$. Next, $\Delta_{k, 2}<D_{g} \int_{m+2 \gamma}^{m^{\prime}} \nu_{k}(d \lambda)=D_{g} \int_{m+2 \gamma}^{m^{\prime}} H_{k}(\lambda) \nu_{0}(d \lambda)$ with $H_{k}(\lambda)$ defined by (13). Since $\beta_{k} \in\left[m^{\prime}, M^{\prime}\right]$ for all $k, H_{k}(\lambda)$ is a decreasing function of $\lambda$ for $\lambda \in\left[m, m^{\prime}\right]$, and for $m+2 \gamma<\lambda<m^{\prime}$ it satisfies

$$
H_{k}(\lambda)<H_{k}(m+2 \gamma)<H_{k}(m+\gamma)\left(\frac{M-m-\varepsilon-2 \gamma}{M-m-\varepsilon-\gamma}\right)^{2 k}
$$

Since $\int_{m}^{m+\gamma} \nu_{k}(d \lambda)=\int_{m}^{m+\gamma} H_{k}(\lambda) \nu_{0}(d \lambda) \geq H_{k}(m+\gamma) \int_{m}^{m+\gamma} \nu_{0}(d \lambda)$, we obtain

$$
\Delta_{k, 2}<\frac{D_{g}}{\int_{m}^{m+\gamma} \nu_{0}(d \lambda)}\left[\frac{M-m-\varepsilon-2 \gamma}{M-m-\varepsilon-\gamma}\right]^{2 k}
$$

We also obtain for the last term,

$$
\Delta_{k, 3}<D_{g} \int_{m^{\prime}}^{C} H_{k}(\lambda) \nu_{0}(d \lambda)<D_{g} \theta_{\varepsilon}^{k} H_{k}(m) \int_{m^{\prime}}^{C} \nu_{0}(d \lambda)<D_{g} \theta_{\varepsilon}^{k} H_{k}(m) \text { for } k>k_{0}
$$

For $\lambda \in\left[m, m^{\prime}\right]$ we have $H_{k}(\lambda) / H_{k}(m) \geq\left(m^{\prime}-\lambda\right)^{2 k} / \varepsilon^{2 k}$ so that

$$
1 \geq \int_{m}^{m+\gamma} \nu_{k}(d \lambda) \geq H_{k}(m) \int_{m}^{m+\gamma}\left[\left(m^{\prime}-\lambda\right) / \varepsilon\right]^{2 k} \nu_{0}(d \lambda)>H_{k}(m)[(\varepsilon-\gamma) / \varepsilon]^{2 k} \int_{m}^{m+\gamma} \nu_{0}(d \lambda)
$$

Therefore, for $k>k_{0}$,

$$
\Delta_{k, 3}<\frac{D_{g}}{\int_{m}^{m+\gamma} \nu_{0}(d \lambda)}\left[\frac{\theta_{\varepsilon} \varepsilon^{2}}{(\varepsilon-\gamma)^{2}}\right]^{k}
$$

We have $\theta_{\varepsilon} \varepsilon^{2} /(\varepsilon-\gamma)^{2}<1$ for $\gamma<\varepsilon\left(1-\sqrt{\theta_{\varepsilon}}\right)$ so that (41) is satisfied for $\alpha_{\gamma}=\max \left\{\theta_{\varepsilon} \varepsilon^{2} /(\varepsilon-\gamma)^{2},(M-\right.$ $\left.m-\varepsilon-2 \gamma)^{2} /(M-m-\varepsilon-\gamma)^{2}\right\}$ and $\alpha_{\gamma}<1$ for $\gamma$ small enough. One can show a similar property for $\Delta_{k}^{\prime}=\left|\int_{C}^{M} g(\lambda) \nu_{k}(d \lambda)-g(M) \int_{C}^{M} \nu_{k}(d \lambda)\right|$.

Finally, we apply the property above to $g(\lambda)=\lambda$ and $g(\lambda)=\lambda^{2}$ and, following the same line as in the proof of Theorem 2, we then obtain for the product of rates at two successive iterations with control variables $\beta_{2 j}$ and $\beta_{2 j+1}=m+M-\beta_{2 j}$ :

$$
R_{2}^{2}\left(\beta_{2 j}\right)\left[1-\frac{A_{\gamma} \alpha_{\gamma}^{2 j}+B \delta}{R_{2}^{2}\left(\beta_{2 j}\right)}\right]<r_{2 j} r_{2 j+1}<R_{2}^{2}\left(\beta_{2 j}\right)\left[1+\frac{A_{\gamma} \alpha_{\gamma}^{2 j}+B \delta}{R_{2}^{2}\left(\beta_{2 j}\right)}\right]
$$

for some $A_{\gamma}>0, B>0$ and $j>k_{0} / 2$. Therefore,

$$
\log R_{2}\left(\beta_{2 j}\right)-A_{\gamma}^{\prime} \alpha_{\gamma}^{2 j}-B^{\prime} \delta<\log \sqrt{r_{2 j} r_{2 j+1}}<\log R_{2}\left(\beta_{2 j}\right)+A_{\gamma}^{\prime} \alpha_{\gamma}^{2 j}+B^{\prime} \delta
$$

with $A_{\gamma}^{\prime}=A_{\gamma} / R_{2}^{2}(m+\varepsilon)$ and $B^{\prime}=B / R_{2}^{2}(m+\varepsilon)$, for $j$ large enough. Since $\sum_{j=0}^{\infty} \alpha_{\gamma}^{2 j}=1 /\left(1-\alpha_{\gamma}^{2}\right)<\infty$, we obtain from (12),

$$
\left|\log R-\int \log R_{2}(\beta) d F(\beta)\right|=\left|\lim _{k \rightarrow \infty} \frac{1}{k} \sum_{j=0}^{k-1} \log \sqrt{r_{2 j} r_{2 j+1}}-\int \log R_{2}(\beta) d F(\beta)\right|<B^{\prime} \delta
$$

Since $\delta$ is arbitrary, the asymptotic rate of convergence is thus the same as in the finite dimensional case.

Remark 9 Note that the condition $I(\lambda)$ being a continuous function of $\lambda$ is satisfied for all examples considered in Section 3. It is also satisfied when the distribution function $F(\cdot)$ has density $\phi(\cdot)$ with derivative $\phi^{\prime}(\cdot)$ uniformly bounded on $\left(m^{\prime}, M^{\prime}\right)$. Indeed, one can write $I(\lambda)=\int_{\lambda-M^{\prime}}^{\lambda-m^{\prime}} \phi(\lambda-t) \log t^{2} d t$ which has derivative $I^{\prime}(\lambda)=\phi\left(m^{\prime}\right) \log \left(\lambda-m^{\prime}\right)^{2}-\phi\left(M^{\prime}\right) \log \left(\lambda-M^{\prime}\right)^{2}+\int_{m^{\prime}}^{M^{\prime}} \phi^{\prime}(t) \log (\lambda-t)^{2} d t$; this derivative is bounded, which implies the continuity of $I(\lambda)$.

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