# On the employment of Inexact Restoration for the minimization of functions whose evaluation is subject to errors * 

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#### Abstract

Inexact Restoration is a well established technique for continuous minimization problems with constraints. Recently, it has been used by Krejić and Martínez for optimization of functions whose evaluation is necessarily inexact and comes from an iterative process. This technique will be generalized in the present paper and it will be applied to stochastic optimization and related problems. New convergence results will be given and numerical results will be presented.


Key words: Inexact Restoration, stochastic programming, global convergence, numerical experiments.

## 1 Introduction

We will consider the problem

$$
\begin{equation*}
\text { Minimize } f(x) \text { subject to } x \in \Omega \text {, } \tag{1}
\end{equation*}
$$

where $\Omega \subseteq \mathbb{R}^{n}$ is nonempty, closed, and convex. We assume that, for different reasons, the exact evaluation of $f(x)$ is not available. In Stochastic Optimization the objective function uses to be (inexactly) evaluated employing a sample $\left\{\xi^{1}, \ldots, \xi^{N}\right\}$ of the random variable under consideration. As a consequence, the evaluation of $f(x)$ is always inexact, and precision is increased when the size $N$ of the sample grows. In other cases, the objective function $f(x)$ is very hard to evaluate, if not impossible. As a consequence, approximations need to be employed.

[^0]For example, suppose that $\xi_{1}, \ldots, \xi_{n}$ denote the future prices of $n$ assets and that $x_{j}$ is the fraction of the budget invested on asset $j$. (In this case we need the constraints $\sum_{j=1}^{n} x_{j}=$ 1 and $x \geq 0$.) Then the return is a random variable given by $g(x, \xi)=\sum_{j=1}^{n} \xi_{j} x_{j}$ and it is natural to maximize its expectation under additional suitable constraints.

The Sample Path Optimization or Sample Average Approximation (SAA) method (see, for example, [23]) applies to optimization problems for which the objective function cannot be computed exactly (or is unaffordable) but can be estimated by simulation. Given $N$ and an independent and identically distributed (i.i.d.) sample $\xi^{1}, \xi^{2}, \ldots, \xi^{N}$ of $N$ replications of the random vector $\xi$, one considers the function

$$
f_{N}(x) \equiv \frac{1}{N} \sum_{i=1}^{N} g\left(x, \xi^{i}\right)
$$

and the problem given by

$$
\begin{equation*}
\operatorname{Minimize}_{x \in \Omega} f_{N}(x) \tag{2}
\end{equation*}
$$

is approximately solved. Problem (2) is deterministic and, in principle, any standard nonlinear programming solver may be employed for its solution.

Assume that we are able to solve (2) globally up to a tolerance $\varepsilon_{\text {glo }}$, i.e. we can compute $x_{N}^{*}\left(\varepsilon_{\text {glo }}\right)$ such that $f_{N}\left(x_{N}^{*}\left(\varepsilon_{\text {glo }}\right)\right) \leq f_{N}(x)+\varepsilon_{\text {glo }}$, for all $x \in \Omega$. Let $\varepsilon_{\text {prob }} \in(0,1)$ and consider the fulfilment of

$$
\begin{equation*}
\left|f_{N}\left(x_{N}^{*}\left(\varepsilon_{\text {glo }}\right)\right)-f^{*}\right| \leq \varepsilon_{\text {glo }} \text { with probability at least } 1-\varepsilon_{\text {prob }}, \tag{3}
\end{equation*}
$$

where $f^{*}$ is the optimal value of problem (11). Then, the decision to be made in the SAA method is to determine the value of $N$ for which solving a single problem of the form (2) guarantees that (3) will hold. A conservative answer to this question, thus yielding to a potentially large value of $N$ and, in consequence, a costly SAA problem (2), is given in [24]. Usually, an infinite sequence of SAA problems with sample size $N_{k}$ and optimality tolerance $\left(\varepsilon_{\text {glo }}\right)_{k}$, such that $N_{k} \rightarrow \infty$ and $\left(\varepsilon_{\text {glo }}\right)_{k} \rightarrow 0$, needs to be solved (see, for example, [25]). The sequence $\left\{N_{k}\right\}_{k=1}^{\infty}$ is known as schedule. Many papers in the literature deal with the theoretical aspects and practical issues (such as convergence, convergence rates, and efficiency) associated with the determination of a schedule and the properties of the optimization method used to solve the SAA optimization problems. See, for example, [3, 4, 14, 16, 17, 18, 21, 22] and the references therein.

Inexact Restoration (IR) [1, 2, [5, 6, 9, 11, 12, 15, 19, 20] is a paradigm for nonlinear programming that have been proved to be effective in solving constrained optimization problems in which some structure of the feasible set induces a natural way of recovering feasibility. In this work, we will show in which way the Inexact Restoration approach can be used to tackle problem (1).

The rest of this work is organized as follows. In Section 2 the proposed IR algorithm is presented. Convergence theory is established in Section 3. Section 4 approaches crucial implementation details and presents numerical experiments. Final remarks are given in Section 5.

Notation. $\|\cdot\|$ denotes the Euclidean norm. $\mathbb{R}_{+}$denotes the set of non-negative real numbers. $\mathbb{N}_{+}$denotes the positive integer numbers.

## 2 Inexact Restoration algorithm

The Inexact Restoration (IR) algorithm described in this section aims to approximate the minimum of $f(x)$ by means of successive evaluations of $f_{N}(x)$ employing samples of $N$ elements. More generally, we may think that we wish to minimize $f(x)$ and that $\lim _{N \rightarrow \infty} f_{N}=f$ in some sense. The number of sample elements $N$ varies (perhaps decreasing) after each evaluation of $f_{N}$ as a result of the IR tests and decisions. A restoration step in the IR algorithm corresponds to increase the number of elements in the sample whereas the optimization step corresponds to improving the value of $f_{N}(x)$. Strictly speaking the value of $f_{N}(x)$ also depends on the chosen sample elements (say, $f_{N}(x)=f_{N}\left(\xi^{1}, \ldots, \xi^{N}, x\right)$ but we will maintain the notation that avoids the mention of the sample elements in order to simplify the notation. The IR approach presented in this paper is inspired by classical Inexact Restoration algorithms for constrained optimization. This idea has been used in [18] in the context of the minimization of functions whose evaluation depends on an iterative process. The main idea consists of considering that the exact evaluation of the objective function corresponds to feasibility of an implicit constrained optimization problem. In such a way, restoration steps correspond to steps where we increase the precision employed to evaluate the objective function. At optimization steps one essentially maintains the sample and one tries to reduce the corresponding sample-based objective function. The IR machinery fits well with this approach and suggests practical ideas for implementation.

For our analysis, we consider a continuous function $h: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$such that $h(\delta)=0$ if and only if $\delta=0$ and an integer-valued function $\bar{N}: \mathbb{R}_{+} \rightarrow \mathbb{N}_{+}$. The practical meaning of these functions will be given after proving the main IR results in order to stress the fact that the IR machinery is independent of this meaning. Moreover, define, for all $x \in \mathbb{R}^{n}, \delta \in \mathbb{R}_{+}, \theta \in(0,1)$, and $N \geq \bar{N}(\delta)$, the function

$$
\Phi(x, \delta, N, \theta)=\theta f_{N}(x)+(1-\theta) h(\delta) .
$$

Inspired by classical IR algorithms, the function $\Phi$ aims to take into account optimality defined by the functional values on a sample and feasibility defined by $h$.

Algorithm 2.1. Let $x_{0} \in \Omega, \delta_{0} \in \mathbb{R}_{+}, N_{0} \geq \bar{N}\left(\delta_{0}\right), \theta_{0} \in(0,1), r \in(0,1), \alpha>0$, and $\beta>0$ be given. Set $k \leftarrow 0$.

Step 1. Restoration Phase
Define $\delta_{k}^{\mathrm{re}} \in \mathbb{R}_{+}$in such a way that

$$
\begin{equation*}
h\left(\delta_{k}^{\mathrm{re}}\right) \leq r h\left(\delta_{k}\right) \tag{4}
\end{equation*}
$$

and a sample with $N_{k+\frac{1}{2}}$ elements such that $N_{k+\frac{1}{2}} \geq \bar{N}\left(\delta_{k}^{\mathrm{re}}\right)$ and

$$
\begin{equation*}
f_{N_{k+\frac{1}{2}}}\left(x_{k}\right) \leq f_{N_{k}}\left(x_{k}\right)+\beta h\left(\delta_{k}\right) . \tag{5}
\end{equation*}
$$

Step 2. Updating the penalty parameter

If

$$
\begin{equation*}
\Phi\left(x_{k}, \delta_{k}^{\mathrm{re}}, N_{k+\frac{1}{2}}, \theta_{k}\right) \leq \Phi\left(x_{k}, \delta_{k}, N_{k}, \theta_{k}\right)+\frac{1-r}{2}\left(h\left(\delta_{k}^{\mathrm{re}}\right)-h\left(\delta_{k}\right)\right) \tag{6}
\end{equation*}
$$

set $\theta_{k+1}=\theta_{k}$. Otherwise, set

$$
\begin{equation*}
\theta_{k+1}=\frac{(1+r)\left(h\left(\delta_{k}\right)-h\left(\delta_{k}^{\mathrm{re}}\right)\right)}{2\left(f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)-f_{N_{k}}\left(x_{k}\right)+h\left(\delta_{k}\right)-h\left(\delta_{k}^{\mathrm{re}}\right)\right)} . \tag{7}
\end{equation*}
$$

Step 3. Optimization Phase
Step 3.1. Choose $d_{k} \in \mathbb{R}^{n}$ such that $x_{k}+d_{k} \in \Omega$.
Step 3.2. Choose $\delta_{k+1} \in \mathbb{R}_{+}$(perhaps bigger than $\left.\delta_{k}^{\mathrm{re}}\right)$ and $N_{k+1} \geq \bar{N}\left(\delta_{k+1}\right)$. If

$$
\begin{equation*}
f_{N_{k+1}}\left(x_{k}+d_{k}\right) \leq f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)-\alpha\left\|d_{k}\right\|^{2} \tag{8}
\end{equation*}
$$

and

$$
\begin{equation*}
\Phi\left(x_{k}+d_{k}, \delta_{k+1}, N_{k+1}, \theta_{k+1}\right) \leq \Phi\left(x_{k}, \delta_{k}, N_{k}, \theta_{k+1}\right)+\frac{1-r}{2}\left(h\left(\delta_{k}^{\mathrm{re}}\right)-h\left(\delta_{k}\right)\right) \tag{9}
\end{equation*}
$$

set $t_{k}=1, x_{k+1}=x_{k}+d_{k}, k \leftarrow k+1$, and go to Step 1 .
Step 3.3. Re-define $\delta_{k+1} \leftarrow \delta_{k}^{\text {re }}$ and $N_{k+1} \leftarrow N_{k+\frac{1}{2}}$. Find $t_{k} \in\left\{1,10^{-1}, 10^{-2}, \ldots\right\}$ as large as possible such that

$$
\begin{equation*}
f_{N_{k+1}}\left(x_{k}+t_{k} d_{k}\right) \leq f_{N_{k+1}}\left(x_{k}\right)-\alpha t_{k}\left\|d_{k}\right\|^{2} . \tag{10}
\end{equation*}
$$

Set $x_{k+1}=x_{k}+t_{k} d_{k}, k \leftarrow k+1$, and go to Step 1 .
Remark. At Step 3.2 the algorithm may decide that the number of elements in the sample must decrease. The idea is to implement the algorithm in such a way that a decrease of $N$ is tried when the iterate is far from the solution. The rationale behind this is that very expensive evaluations are not worthwhile when the functional value is still large. The IR framework provides a test to decide if the reduction of $N$ is acceptable. Note that, at first sight, the occurrence of infinitely many steps 3.2 could lead to the generation of a bounded sequence $\left\{N_{k}\right\}$ and, so, to the non-convergence of $f_{N}$ to $f$. This possibility will be discarded by the convergence theory. In [18] we proved that, for a similar approach, a finite precision defined by a bound for $N$ is reached. The proof that the satisfactory optimality result holds with $N$ going to infinity is given in the present paper. The motivations for the tests that guarantee arbitrary precision and optimality are provided by the Inexact Restoration techniques developed for smooth constrained optimization. Step 3.2 corresponds to large trial steps in the tangent approximation to the approximate feasible set and Step 3.3 corresponds to backtracking in the tangent space.

## 3 Convergence theory

In this section we will prove that the sequence generated by Algorithm 2. 1 converges to a suitable approximation of a solution to (11). Roughly speaking, limit points of the sequence
will be stationary points with a probability as high as desired and the approximation of the sample-computed function $f_{N}$ will be satisfactory also with probability arbitrarily close to 1 . In the theoretical proofs we will employ the following assumptions.

Assumption A1 The Restoration Phase defined by Step 1 of Algorithm 2.1 can be computed in finite time for all $k$.

Assumption A2 There exists $c_{L}>0$ such that, for all $k$,

$$
\left\|\nabla f_{N_{k+\frac{1}{2}}}(x)-\nabla f_{N_{k+\frac{1}{2}}}(y)\right\| \leq c_{L}\|x-y\| \text { for all } x, y \in \Omega,
$$

i.e. the gradient of $f_{N_{k+\frac{1}{2}}}(\cdot)$ is Lipschitz-continuous with the same Lipschitz constant for all $k$. In addition, there exists $c>0$ such that, for all $k$, we have that $f_{N_{k}}(x) \leq c$ for all $x \in \Omega$.

Assumption A3 The iterates $x_{k}$ generated by Algorithm 2.1 lie in a bounded set. (This assumption is obviously satisfied if $\Omega$ is bounded.)

Assumption A4 There exists $\tau>0$ such that, for all $k$, the chosen direction $d_{k}$ satisfies

$$
\begin{equation*}
f_{N_{k+\frac{1}{2}}}\left(x_{k}+t d_{k}\right) \leq f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)-\alpha t\left\|d_{k}\right\|^{2} \text { for all } t \in[0, \tau] . \tag{11}
\end{equation*}
$$

The plausibility of Assumption A4 will be shown in the following two lemmas. In Lemma 3.1 we prove that Assumption A4 is satisfied in the unconstrained case under a very general gradientrelated choice for $d_{k}$. In Lemma 3.2 we prove that Assumption A4 also holds for a general closed and convex $\Omega$ if we take $d_{k}$ as a gradient projection. Note that the value of $\tau>0$ in Assumption A4 does not depend on $k$. This is the reason for which, in Lemma 3.1 below, the upper bound on the norm of $d_{k}$ is required.

Lemma 3.1 Suppose that Assumption A2 holds. If there exist $c_{\text {angle }} \in(0,1)$ and $c_{\text {big }}>c_{\text {small }}>$ 0 such that for all $k \in \mathbb{N}$ we choose $d_{k}$ satisfying

$$
\begin{equation*}
d_{k}^{T} \nabla f_{N_{k+\frac{1}{2}}}\left(x_{k}\right) \leq-c_{\text {angle }}\left\|d_{k}\right\|\left\|\nabla f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)\right\| \tag{12}
\end{equation*}
$$

and

$$
\begin{equation*}
c_{\text {big }}\left\|\nabla f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)\right\| \geq\left\|d_{k}\right\| \geq c_{\text {small }}\left\|\nabla f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)\right\|, \tag{13}
\end{equation*}
$$

then Assumption A4 holds for any $\alpha \in\left(0, \frac{1}{2} c_{\text {angle }} / c_{\text {big }}\right]$ taking $\tau=2 \alpha / c_{\mathrm{L}}$.
Proof: By Assumption A2 and elementary Calculus, we have that, for all $t>0$,

$$
\begin{equation*}
f_{N_{k+\frac{1}{2}}}\left(x_{k}+t d_{k}\right) \leq f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)+t d_{k}^{T} \nabla f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)+c_{\mathrm{L}} t^{2}\left\|d_{k}\right\|^{2} / 2 . \tag{14}
\end{equation*}
$$

Consequently, by (12),

$$
\begin{equation*}
f_{N_{k+\frac{1}{2}}}\left(x_{k}+t d_{k}\right) \leq f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)-t c_{\text {angle }}\left\|d_{k}\right\|\left\|\nabla f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)\right\|+c_{\mathrm{L}} t^{2}\left\|d_{k}\right\|^{2} / 2 . \tag{15}
\end{equation*}
$$

Therefore, by (13),

$$
\begin{equation*}
f_{N_{k+\frac{1}{2}}}\left(x_{k}+t d_{k}\right) \leq f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)+\left(c_{\mathrm{L}} t^{2} / 2-t c_{\text {angle }} / c_{\text {big }}\right)\left\|d_{k}\right\|^{2} \tag{16}
\end{equation*}
$$

and, since $\alpha \leq \frac{1}{2} c_{\text {angle }} / c_{\text {big }}$,

$$
\begin{equation*}
f_{N_{k+\frac{1}{2}}}\left(x_{k}+t d_{k}\right) \leq f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)+\left(c_{\mathrm{L}} t^{2} / 2-2 \alpha t\right)\left\|d_{k}\right\|^{2} . \tag{17}
\end{equation*}
$$

Therefore, (11) holds taking $\tau=2 \alpha / c_{\mathrm{L}}$.
Lemma 3.2 Suppose that Assumption A2 holds. Assume that $\sigma_{\max } \geq \sigma_{\min }>0$ are given parameters and that, for all $k \in \mathbb{N}$, we choose $\sigma_{k} \in\left[\sigma_{\min }, \sigma_{\max }\right]$ and compute $d_{k}=P_{\Omega}\left(x_{k}-\right.$ $\left.\sigma_{k} \nabla f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)\right)-x_{k}$, where $P_{\Omega}$ denotes the projection operator onto $\Omega$. Then, Assumption $A_{4}$ holds taking $\alpha=1 /\left(4 \sigma_{\max }\right)$ and $\tau=\min \left\{1,1 /\left(2 \sigma_{\max } c_{\mathrm{L}}\right)\right\}$.

Proof: By the definition of $d_{k}$ we have that $x_{k}+d_{k}$ is the projection of $x_{k}-\sigma_{k} \nabla f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)$ onto the convex set $\Omega$. Therefore the points $x_{k}, x_{k}-\sigma_{k} \nabla f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)$, and $x_{k}+d_{k}$ are the vertices of a triangle, the angle corresponding to $x_{k}+d_{k}$ is obtuse, and the opposite side to this vertex is the biggest one. Therefore,

$$
\left\|d_{k}+\sigma_{k} \nabla f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)\right\|^{2} \leq\left\|\sigma_{k} \nabla f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)\right\|^{2} .
$$

So,

$$
\left\|d_{k}\right\|^{2}+\sigma_{k}^{2}\left\|\nabla f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)\right\|^{2}+2 \sigma_{k} d_{k}^{T} \nabla f_{N_{k+\frac{1}{2}}}\left(x_{k}\right) \leq\left\|\sigma_{k} \nabla f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)\right\|^{2} .
$$

Thus,

$$
\begin{equation*}
\nabla f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)^{T} d_{k} \leq-\frac{1}{2 \sigma_{k}}\left\|d_{k}\right\|^{2} \tag{18}
\end{equation*}
$$

Since $x_{k}$ and $x_{k}+d_{k} \in \Omega$ and $\Omega$ is convex, we obtain, by Assumption A2 and elementary Calculus, that

$$
f_{N_{k+\frac{1}{2}}}\left(x_{k}+t d_{k}\right) \leq f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)+t \nabla f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)^{T} d_{k}+\frac{1}{2} t^{2} c_{\mathrm{L}}\left\|d_{k}\right\|^{2}
$$

for all $t \in[0,1]$. Therefore, by (18),

$$
f_{N_{k+\frac{1}{2}}}\left(x_{k}+t d_{k}\right) \leq f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)-\frac{1}{2 \sigma_{k}} t\left\|d_{k}\right\|^{2}+\frac{1}{2} t^{2} c_{\mathrm{L}}\left\|d_{k}\right\|^{2}
$$

for all $t \in[0,1]$. Then,

$$
f_{N_{k+\frac{1}{2}}}\left(x_{k}+t d_{k}\right) \leq f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)-\left[\frac{1}{2 \sigma_{\max }}-\frac{1}{2} t c_{\mathrm{L}}\right] t\left\|d_{k}\right\|^{2}
$$

for all $t \in[0,1]$. Thus, since $\tau<1$, for all $t \in[0, \tau]$, we have:

$$
f_{N_{k+\frac{1}{2}}}\left(x_{k}+t d_{k}\right) \leq f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)-\left[\frac{1}{2 \sigma_{\max }}-\frac{1}{2} \tau c_{\mathrm{L}}\right] t\left\|d_{k}\right\|^{2}
$$

and, by the hypothesis on $\tau$,

$$
f_{N_{k+\frac{1}{2}}}\left(x_{k}+t d_{k}\right) \leq f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)-\left[\frac{1}{2 \sigma_{\max }}-\frac{1}{4 \sigma_{\max }}\right] t\left\|d_{k}\right\|^{2}
$$

for all $t \in[0, \tau]$. Then, by the choice of $\alpha$ in the hypothesis, Assumption A4 holds.

In the following results we will assume, without mentioning explicitly, that Assumptions A1A 4 hold. In the next lemma we prove that the penalty parameter $\theta_{k}$ is non-increasing and bounded away from zero. This result is crucial in the IR framework to guarantee that, ultimately, the necessity of minimizing the objective function is preserved and that this goal is not overwhelmed by the search of feasibility.

Lemma 3.3 The sequence $\left\{\theta_{k}\right\}$ generated by Algorithm 2. 1 remains in $(0,1)$, is non-increasing, and is bounded below by a positive quantity that only depends on $\beta$, $r$, and $\theta_{0}$.

Proof: At each iteration, we have that $\theta_{k+1}=\theta_{k}$ when (6) takes place and that $\theta_{k+1}$ is computed by (7) otherwise. If (6) does not hold, it turns out that

$$
\begin{equation*}
\Phi\left(x_{k}, \delta_{k}^{\mathrm{re}}, N_{k+\frac{1}{2}}, \theta\right)>\Phi\left(x_{k}, \delta_{k}, N_{k}, \theta\right)+\frac{1-r}{2}\left(h\left(\delta_{k}^{\mathrm{re}}\right)-h\left(\delta_{k}\right)\right) \tag{19}
\end{equation*}
$$

for $\theta=\theta_{k}$. On the other hand, since $h\left(\delta_{k}^{\mathrm{re}}\right)<h\left(\delta_{k}\right)$ and, thus,

$$
h\left(\delta_{k}^{\mathrm{re}}\right)-h\left(\delta_{k}\right) \leq \frac{1-r}{2}\left(h\left(\delta_{k}^{\mathrm{re}}\right)-h\left(\delta_{k}\right)\right)
$$

we have that the negation of (19) always holds strictly if $\theta=0$. Since $\Phi$ is a linear function on $\theta$, this means that there exists a unique $\theta \in\left(0, \theta_{k}\right)$ that verifies

$$
\begin{equation*}
\Phi\left(x_{k}, \delta_{k}^{\mathrm{re}}, N_{k+\frac{1}{2}}, \theta\right)=\Phi\left(x_{k}, \delta_{k}, N_{k}, \theta\right)+\frac{1-r}{2}\left(h\left(\delta_{k}^{\mathrm{re}}\right)-h\left(\delta_{k}\right)\right) \tag{20}
\end{equation*}
$$

This $\theta$ is the quantity $\theta_{k+1}$ computed at (7). Therefore, in order to prove the lemma, it only remains to be shown that $\theta_{k+1}$ defined by (7) is bounded away from zero. In fact, by (7) and (5) we have:

$$
\begin{gathered}
\frac{1}{\theta_{k+1}}=2\left[\frac{f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)-f_{N_{k}}\left(x_{k}\right)}{h\left(\delta_{k}\right)-h\left(\delta_{k}^{\mathrm{re}}\right)}+1\right] \frac{1}{1+r} \leq \frac{1}{1+r}\left[\frac{f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)-f_{N_{k}}\left(x_{k}\right)}{(1-r) h\left(\delta_{k}\right)}+1\right] \\
\leq \frac{1}{1+r}\left[\frac{\beta h\left(\delta_{k}\right)}{(1-r) h\left(\delta_{k}\right)}+1\right]=\frac{1}{1+r}\left[\frac{\beta}{(1-r)}+1\right]
\end{gathered}
$$

This completes the proof.

Lemma 3.4 shows that each iteration of the algorithm satisfies both the requirement of decreasing the approximate objective function and the function $\Phi$, that combines optimality and feasibility. Note that, together with Assumption A4, Lemma 3.4 implies that Algorithm 2.1 is well defined.

Lemma 3.4 For all $k=0,1,2, \ldots$, (8) and (9) hold and $t_{k}=1$ is set at Step 3.2, or, after a finite number of evaluations which only depend on $\tau$, 10) holds and $t_{k} \in\left\{1,10^{-1}, 10^{-2}, \ldots\right\}$ is set at Step 3.3. In addition, in this case, we have that $t_{k} \geq \tau / 10$ and

$$
\begin{equation*}
\Phi\left(x_{k+1}, \delta_{k+1}, N_{k+1}, \theta_{k+1}\right) \leq \Phi\left(x_{k}, \delta_{k}, N_{k}, \theta_{k+1}\right)+\frac{1-r}{2}\left(h\left(\delta_{k}^{\mathrm{re}}\right)-h\left(\delta_{k}\right)\right) . \tag{21}
\end{equation*}
$$

Proof: The fact that (10) takes place after a finite number of evaluations that only depends on $\tau$ is a trivial consequence of Assumption A4, Also by Assumption A4, the first value of $t$ such that (10) holds is not smaller than $\tau / 10$. Let us show now that, at Step 3.3, when 10) holds, we also obtain (21).

If (6) holds then we have that $\theta_{k+1}=\theta_{k}$ and, therefore, (6) also holds substituting $\theta_{k}$ with $\theta_{k+1}$. If (6) does not hold, the choice (7) of $\theta_{k+1}$ is such that (6) holds (by equality) substituting $\theta_{k}$ with $\theta_{k+1}$. Therefore, in any case we have:

$$
\begin{equation*}
\Phi\left(x_{k}, \delta_{k+1}, N_{k+1}, \theta_{k+1}\right) \leq \Phi\left(x_{k}, \delta_{k}, N_{k}, \theta_{k+1}\right)+\frac{1-r}{2}\left(h\left(\delta_{k}^{\mathrm{re}}\right)-h\left(\delta_{k}\right)\right) . \tag{22}
\end{equation*}
$$

On the other hand,

$$
\begin{aligned}
\Phi\left(x_{k+1}, \delta_{k+1}, N_{k+1}, \theta_{k+1}\right) & =\theta_{k+1} f_{N_{k+1}}\left(x_{k+1}\right)+\left(1-\theta_{k+1}\right) h\left(\delta_{k+1},\right) \\
& \leq \theta_{k+1}\left[f_{N_{k+1}}\left(x_{k}\right)-\alpha t_{k}\left\|d_{k}\right\|^{2}\right]+\left(1-\theta_{k+1}\right) h\left(\delta_{k+1}\right) \\
& =\theta_{k+1} f_{N_{k+1}}\left(x_{k}\right)+\left(1-\theta_{k+1}\right) h\left(\delta_{k+1}\right)-\theta_{k+1} \alpha t_{k}\left\|d_{k}\right\|^{2} \\
& =\Phi\left(x_{k}, \delta_{k+1}, N_{k+1}, \theta_{k+1}\right)-\theta_{k+1} \alpha t_{k}\left\|d_{k}\right\|^{2} \\
& \leq \Phi\left(x_{k}, \delta_{k+1}, N_{k+1}, \theta_{k+1}\right)
\end{aligned}
$$

and (21) holds by (22).
In the following theorem we prove that the algorithm is satisfactory with respect to feasibility, which, in this case, means accurate evaluation of the objective function. This property is represented by the statement $h\left(\delta_{k}^{\mathrm{re}}\right) \rightarrow 0$ when $k \rightarrow \infty$. Note that validity of the theorem is independent of the meaning of $h$ which, up to now, has been associated with accuracy (or feasibility) only in terms of motivation.

Theorem 3.1 The series $\sum_{k=0}^{\infty} h\left(\delta_{k}\right)$ and $\sum_{k=0}^{\infty} h\left(\delta_{k}^{\mathrm{re}}\right)$ are convergent.
Proof: By (4) we only need to prove that $\sum_{k=0}^{\infty} h\left(\delta_{k}\right)$ is convergent. By Lemma 3.3, there exists $\theta_{*}>0$, the limit of the non-increasing sequence $\left\{\theta_{k}\right\}$. Define $\rho_{k}=\left(1-\theta_{k}\right) / \theta_{k}$ for all $k \in \mathbb{N}$. Clearly $\left\{\rho_{k}\right\}$ is non-decreasing and bounded above by $1 / \theta_{*}-1$. Therefore,

$$
\begin{equation*}
\sum_{k=0}^{\infty}\left(\rho_{k+1}-\rho_{k}\right)=\lim _{k \rightarrow \infty} \rho_{k+1}-\rho_{0}<\infty \tag{23}
\end{equation*}
$$

Now, since $\left\{h\left(\delta_{k}\right)\right\}$ is bounded, (23) implies that

$$
\begin{equation*}
\sum_{k=0}^{\infty}\left(\rho_{k+1}-\rho_{k}\right) h\left(\delta_{k}\right)<\infty \tag{24}
\end{equation*}
$$

On the other hand, by (21) in Lemma 3.4, (11), and the definition of $x_{k+1}$, we have that

$$
\begin{equation*}
\Phi\left(x_{k+1}, \delta_{k+1}, N_{k+1}, \theta_{k+1}\right) \leq \Phi\left(x_{k}, \delta_{k}, N_{k}, \theta_{k+1}\right)-\frac{(1-r)^{2}}{2} h\left(\delta_{k}\right) \tag{25}
\end{equation*}
$$

Inequality (25) is equivalent to

$$
f_{N_{k+1}}\left(x_{k+1}\right)+\rho_{k+1} h\left(\delta_{k+1}\right) \leq f_{N_{k}}\left(x_{k}\right)+\rho_{k+1} h\left(\delta_{k}\right)-\frac{(1-r)^{2}}{2 \theta_{k+1}} h\left(\delta_{k}\right) .
$$

Adding and subtracting $\rho_{k} h\left(\delta_{k}\right)$ and rearranging, as $\theta_{k+1} \in(0,1)$, we get:

$$
\frac{(1-r)^{2}}{2} h\left(\delta_{k}\right) \leq\left(\rho_{k+1}-\rho_{k}\right) h\left(\delta_{k}\right)+\left[\left(f_{N_{k}}\left(x_{k}\right)+\rho_{k} h\left(\delta_{k}\right)\right)-\left(f_{N_{k+1}}\left(x_{k+1}\right)+\rho_{k+1} h\left(\delta_{k+1}\right)\right)\right]
$$

Summing over all $k$, we obtain:

$$
\sum_{k=0}^{\infty} \frac{(1-r)^{2}}{2} h\left(\delta_{k}\right) \leq \sum_{k=0}^{\infty}\left(\rho_{k+1}-\rho_{k}\right) h\left(\delta_{k}\right)+\left[\left(f_{N_{0}}\left(x_{0}\right)+\rho_{0} h\left(\delta_{0}\right)\right)-\lim _{k \rightarrow \infty}\left(f_{N_{k}}\left(x_{k}\right)+\rho_{k} h\left(\delta_{k}\right)\right)\right]
$$

that, by (24) and the boundedness of $\left\{x_{k}\right\}$ and $\left\{\rho_{k}\right\}$, implies that

$$
\sum_{k=0}^{\infty} \frac{(1-r)^{2}}{2} h\left(\delta_{k}\right)<\infty
$$

Thus,

$$
\sum_{k=0}^{\infty} h\left(\delta_{k}\right)<\infty
$$

as we wanted to prove.
Finally, in Theorem 3.2, we prove that we are able to find points at which the gradient of the approximate sample-based function is as small as desired.

Theorem 3.2 The series $\sum_{k=0}^{\infty}\left\|d_{k}\right\|^{2}$ is convergent.
Proof: The inequalities (8) and (10), and the lower bound for $t_{k}$ stated at Lemma 3.4 (say $t_{k} \geq t_{*}$ for all $k$ ) imply that

$$
\begin{aligned}
f_{N_{k+1}}\left(x_{k+1}\right)-f_{N_{k}}\left(x_{k}\right) & =f_{N_{k+1}}\left(x_{k+1}\right)-f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)+f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)-f_{N_{k}}\left(x_{k}\right) \\
& \leq-\alpha t_{k}\left\|d_{k}\right\|^{2}+f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)-f_{N_{k}}\left(x_{k}\right) \\
& \leq-\alpha t_{*}\left\|d_{k}\right\|^{2}+\left|f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)-f_{N_{k}}\left(x_{k}\right)\right| \\
& \leq-\alpha t_{*}\left\|d_{k}\right\|^{2}+\beta h\left(\delta_{k}\right) .
\end{aligned}
$$

Thus, for an arbitrary integer $k$, we have

$$
f_{N_{k}}\left(x_{k}\right)-f_{N_{0}}\left(x_{0}\right) \leq-\alpha t_{*} \sum_{j=0}^{k-1}\left\|d_{j}\right\|^{2}+\beta \sum_{j=k}^{k-1} h\left(\delta_{j}\right) .
$$

Therefore, by Assumptions A2 and A3, and Theorem 3.1,

$$
\alpha t_{*} \sum_{j=0}^{\infty}\left\|d_{j}\right\|^{2}<\infty
$$

as we wanted to prove.
If $h$ is associated with feasibility then Theorem 3.1 means that limit points of the sequence generated by Algorithm 21 are feasible points; while if $h$ is associated with precision of evaluation of the objective function of (1], Theorem 3.1 means that in the limit the objective function is evaluated exactly. On the other hand, Theorem 3.2 implies that limit points satisfy an optimality condition as well.

Corollary 3.1 If $\Omega=\mathbb{R}^{n}$ and we choose $\alpha$ and $d_{k}$ as in the hypotheses of Lemma 3.1, we have that

$$
\lim _{k \rightarrow \infty}\left\|\nabla f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)\right\|=0
$$

Proof: The desired result follows from (13) and Theorem 3.2.
Corollary 3.2 If $\Omega$ is an arbitrary closed and convex set and we choose $\alpha$ and $d_{k}$ as in the hypotheses of Lemma 3.2, we have that

$$
\lim _{k \rightarrow \infty}\left\|P_{\Omega}\left(x_{k}-\sigma_{k} \nabla f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)\right)-x_{k}\right\|=0
$$

for all $k \in \mathbb{N}$.
Proof: The desired result follows from Theorem 3.2.

## 4 Implementation features and numerical experiments

Algorithm 2.1, as well as the Projected Gradient method (that will be used for comparison purposes), were implemented in Fortran 90. All tests were conducted on a computer with 3.5 GHz Intel Core i7 processor and 16GB 1600 MHz DDR3 RAM memory, running OS X Yosemite (version 10.10.4). Codes were compiled by the GFortran Fortran compiler of GCC (version 4.9.2) with the -O3 optimization directive enabled. Regarding the parameters of Algorithm 2.1, arbitrarily but based on previous experimentation with Inexact Restoration methods, we set $\delta_{0}=0.01, N_{0}=\bar{N}\left(\delta_{0}\right), \theta_{0}=0.9$, and $\alpha=10^{-4}$. The remaining parameters will be mentioned below.

### 4.1 Implementation details

Defining $h(\delta)=\delta$, and $\bar{N}(\delta)=\lceil 1 / \delta\rceil$, the interpretation of the convergence results presented in the previous section is that, given $\varepsilon_{\text {opt }}>0$ and $\underline{N}>0$, Algorithm 2.1 obtains, in finite time, a sample with $N_{k+1} \geq \underline{N}$ elements and a point $x_{k}$ such that $x_{k}$ satisfies an optimality condition
for minimizing $f_{N_{k+1}}(x)$ subject to $x \in \Omega$ with precision $\varepsilon_{\mathrm{opt}}$. At Step 3.1 of Algorithm 2. 1 , we considered

$$
d_{k}=P_{\Omega}\left(x_{k}-\nabla f_{N_{k+1}}\left(x_{k}\right)\right)-x_{k} .
$$

So, by Lemma 3.2, if Assumption A2 holds, Assumption A4 holds as well. With this choice, a natural stopping criterion for Algorithm 2. 2 is given by

$$
\begin{equation*}
N_{k+1} \geq \underline{N} \text { and }\left\|P_{\Omega}\left(x_{k}-\nabla f_{N_{k+1}}\left(x_{k}\right)\right)-x_{k}\right\|_{\infty} \leq \varepsilon_{\mathrm{opt}} \tag{26}
\end{equation*}
$$

Moreover, due to the choice of the search direction, in order to evaluate the performance of Algorithm 2.1, it would be natural to compare it against the classical Projected Gradient (PG) method applied to the minimization of $f_{\underline{N}}(\cdot)$ and considering the stopping criterion

$$
\begin{equation*}
\left\|P_{\Omega}\left(x_{k}-\nabla f_{\underline{N}}\left(x_{k}\right)\right)-x_{k}\right\|_{\infty} \leq \varepsilon_{\mathrm{opt}} . \tag{27}
\end{equation*}
$$

At Step 3.2, we always try $\delta_{k+1}=0.01$ and $N_{k+1}=\bar{N}\left(\delta_{k+1}\right)$. The idea behind this radical choice is to drastically reduce sample sizes that may be unnecessarily large at the beginning or far for a solution.

The freedom in the condition (4) at Step 1 allows one to use an adaptive choice for the samples employed at each iteration of the method. The idea is that, if one is far from the solution, it is not worthwhile to use a very big sample, but large samples are justified if we think that we are close to a solution. Specifically, we define

$$
\delta_{k}^{\mathrm{re}}= \begin{cases}r_{1} \delta_{k}, & \text { if } N_{k} \geq \underline{N} \text { or }\left\|P_{\Omega}\left(x_{k-1}-\nabla f_{N_{k}}\left(x_{k-1}\right)\right)-x_{k-1}\right\|_{\infty}>\varepsilon_{\mathrm{opt}},  \tag{28}\\ r_{2} \delta_{k}, & \text { otherwise }\end{cases}
$$

where $r_{1}=1-10^{-6}$ and $r_{2}=0.1$ (the parameter $r$ being given by $\max \left\{r_{1}, r_{2}\right\}$ ). The effect of this should be that a suitable approximation to a solution should be obtained employing reasonable small samples being big samples reserved to the case in which we are close to a solution. There is no practical way to guarantee that, for a given value of $\beta>0$, the choice 28) of $\delta_{k}^{\text {re }}$ plus $N_{k+\frac{1}{2}}=\bar{N}\left(\delta_{k}^{\mathrm{re}}\right)$ will satisfy (5). Therefore, in practice, we compute

$$
\max _{k=0,1,2, \ldots}\left\{\left[\beta_{k}\right]_{+}\right\}
$$

where

$$
\beta_{k} \equiv \frac{f_{N_{k+\frac{1}{2}}}\left(x_{k}\right)-f_{N_{k}}\left(x_{k}\right)}{h\left(\delta_{k}\right)}
$$

and $[\cdot]_{+}=\max \{0, \cdot\}$, in order to verify whether this sequence appears to be bounded above (meaning that Assumption A1 holds) or not.

### 4.2 A classification problem

We now describe a classification scheme that can be written as an optimization problem. Assume that there exists an (oracle) unknown function $\omega: B \rightarrow\{1,-1\}$ that is used to classify any $\xi \in B \subseteq \mathbb{R}^{m}$ attributing to it 1 or -1 . To fix ideas we may think in $B=\left\{\xi \in \mathbb{R}^{m} \mid \xi_{\min } \leq \xi_{i} \leq\right.$
$\left.\xi_{\max }, i=1, \ldots, m\right\}$, where $\xi_{\max } \geq \xi_{\min }$ are given values. Usually, $N>0$, an independent and identically distributed sample or training set of labelled examples

$$
D=\left\{\left(\xi^{i}, \omega_{i}\right), i=1, \ldots, N, \text { where } \xi^{i} \in B \text { and } \omega_{i} \equiv \omega\left(\xi^{i}\right) \in\{1,-1\}\right\}
$$

are given, and the goal is to find a classifier $C_{x}(\xi)$ that depends on unknown parameters $x \in \mathbb{R}^{n}$ such that $C_{x}\left(\xi^{i}\right) \leq 0$ if $\omega\left(\xi^{i}\right)=-1$ and $C_{x}\left(\xi^{i}\right) \geq 0$, if $\omega\left(\xi^{i}\right)=1$. Since such classifier may not exist, one seeks to minimize the function given by

$$
f(x)=f_{N}\left(\xi^{1}, \ldots, \xi^{N}, x\right)=\frac{1}{N}\left[\sum_{i \in I} \max \left\{0, C_{x}\left(\xi^{i}\right)\right\}^{2}+\sum_{i \in O} \max \left\{0,-C_{x}\left(\xi^{i}\right)\right\}^{2}\right]
$$

where $I=\left\{i \in\{1, \ldots, N\} \mid \omega_{i}=-1\right\}$ and $O=\left\{i \in\{1, \ldots, N\} \mid \omega_{i}=1\right\}$, subject to appropriate constraints on the parameters $x$ (that may define the shape of the classifier). This was the problem considered in [8] (where it was arbitrarily considered the case $N=10000$ ). On the other hand, in this work, we consider the stochastic problem given by

$$
\text { Minimize } E\left(\sum_{i \in \mathcal{I}} \max \left\{0, C_{x}(\xi)\right\}^{2}+\sum_{i \in \mathcal{O}} \max \left\{0,-C_{x}(\xi)\right\}^{2}\right) \text { subject to } x \in \Omega
$$

where $\mathcal{I}=\{\xi \in B \mid \omega(\xi)=-1\}, \mathcal{O}=\{\xi \in B \mid \omega(\xi)=1\}$, and $\Omega$ represents the constraints on $x$. It is assumed that consulting the oracle $\omega(\cdot)$ is expensive and that the goal is to find $C_{x}(\cdot)$ in order to (a) mimic the oracle answer and (b) understanding the relevance (or not) of each component of $\xi$ in the classification.

Example 1. In this example, we consider $m=2, n=3, C_{x}(\xi)=C_{(c, r)}(\xi)=\|\xi-c\|^{2}-r^{2}$, and $\Omega \equiv \mathbb{R}^{n}$ meaning that the classifier is a circle of radius $|r|$ centered at $c$ in $\mathbb{R}^{2}$. Following [8], we consider four instances of this example corresponding to four different oracles as follows:

Instance 1.1. $\omega(\xi)=-1$ if $\xi$ belongs to a circle with radius 7 centered at the origin and $\omega(\xi)=1$, otherwise.

Instance 1.2. $\omega(\xi)=-1$ if $\xi$ belongs to a square of side 7 centered at the origin and $\omega(\xi)=1$, otherwise.

Instance 1.3. $\omega(\xi)=-1$ if $\xi$ belongs to a rectangle with height equal to 7 and width equal to 14 centered at the origin and $\omega(\xi)=1$, otherwise.

Instance 1.4. $\omega(\xi)=-1$ if $\xi$ belongs to a triangle with vertices $(-7,0),(0,-7)$, and $(7,7)$ and $\omega(\xi)=1$, otherwise.
In all cases we considered $B=\left\{\xi \in \mathbb{R}^{2} \mid-10 \leq \xi_{i} \leq 10, i=1,2\right\}$.
Example 2. In this example, we consider $m=2$, $n=q(q+1), C_{x}(\xi)=C_{(A, b)}(\xi)=\xi^{T} A \xi+$ $b^{T} \xi-1$, and $\Omega \subset \mathbb{R}^{q \times q}$ as the set of positive definite matrices $A \in \mathbb{R}^{q \times q}$ with all their eigenvalues $\lambda_{i}(A), i=1, \ldots, q$, satisfying

$$
0<\lambda_{\min } \leq \lambda_{i}(A) \leq \lambda_{\max }<+\infty
$$

where $0<\lambda_{\min } \leq \lambda_{\max }<+\infty$ are given parameters. This means that the classifier is an ellipse in $\mathbb{R}^{2}$. It is worth noting that, given a square matrix $A$, its projection onto the closed and convex set $\Omega$ can be computed in two steps [10, 13]. In the first step one symmetrizes $A$ and in the second step one computes the $Q D Q^{T}$ decomposition of the symmetrized matrix and replaces its eigenvalues by their projection onto the interval [ $\lambda_{\min }, \lambda_{\max }$ ]. In the numerical experiments, we consider $\lambda_{\min }=10^{-4}$ and $\lambda_{\max }=10^{4}$. Moreover, the same four instances corresponding to the oracles described for Example 1, that here we named Instances 2.1-2.4, are also considered.

### 4.3 Numerical results of the classification problem

As described in two previous subsections, in the numerical experiments we evaluated the performance of the Inexact Restoration method (Algorithm 2.1), called "IR method" from now on, and we compare it against the classical Projected Gradient method, called "PG method" from now on. In the stopping criteria (26) and (27) for the IR and the PG method, respectively, we considered $\varepsilon_{\text {opt }}=10^{-4}$ and $\underline{N} \in\left\{10^{4}, 10^{5}, \ldots, 10^{8}\right\}$.

The cost that dominates the computations in both methods is the evaluation of the objective function. In the case of the PG method, that is applied to the minimization of $f_{\underline{N}}(\cdot)$, the cost of evaluating the objective function is given by $c \underline{N}$, where $c>0$ is a constant that represents the cost of evaluating $C(p)$, i.e. verifying whether a point $p \in \mathbb{R}^{2}$ is "inside" or "outside" the classifier $C$. Therefore, the total cost of the PG method is given by the total number of functional evaluations, "\#fcnt" times $c \underline{N}$. In order to avoid big figures, we will scale this total effort dividing it by $c \underline{N}$ and, therefore, the reported "total effort" of the PG method will be given by the total number of functional evaluations $\#$ fcnt. In the case of the IR method, which evaluates the objective function with samples of varying sizes, the scaled total effort is given by

$$
\frac{c N_{0}}{\underline{N}}+\sum_{k=0}^{\# \mathrm{it}}\left[\frac{c\left(s_{k+\frac{1}{2}} N_{k+\frac{1}{2}}+s_{k+1} N_{k+1}\right)}{c \underline{N}}\right]=\frac{1}{\underline{N}}\left(N_{0}+\sum_{k=0}^{\# \mathrm{it}}\left[s_{k+\frac{1}{2}} N_{k+\frac{1}{2}}+s_{k+1} N_{k+1}\right]\right),
$$

where $\#$ it is the number of iterations performed until the stopping criterion is reached and $s_{k+\frac{1}{2}}$ and $s_{k+1}$ are the number of functional evaluations with samples of sizes $N_{k+\frac{1}{2}}$ and $N_{k+1}$, respectively, performed at iteration $k$.

Table 1 presents a comparison of the performances of both methods when applied to Instances 1.1-1.4 of Example 1 and Instances 2.1-2.4 of Example 2 varying $\underline{N} \in\left\{10^{4}, 10^{5}, \ldots, 10^{8}\right\}$. In all cases, except the case of the PG method applied to Instance 2.2, the methods achieved their respective stopping criteria and found qualitative equivalent solutions. In the case in which the PG method was unable to achieve its stopping criterion (marked with "-" in the table), it stopped due to a very small step in the line search (smaller than $10^{-16}$ ). Moreover, even in this case, the projected gradient was relatively small (of the order of $10^{-3}$ ) and the final iterate was a "reasonable approximation" to a solution. Figure 1 shows a pictorial representation of the solutions found.

A comparison of the figures in the table shows that the IR method is more efficient in 32 out of the 40 considered combinations of instances and $\underline{N}$ (i.e. $80 \%$ of the cases). The PG method is more efficient in seven other instances and there is a tie between the two methods in a single instance. The instances in which the PG method is competitive are related to small values of $\underline{N}$;

|  |  | Example 1: Circular classifier |  | Example 2: Ellipsoidal classifier |  |
| :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | Total Effort |  | Total Effort |  |
|  |  | PG method | IR method | PG method | IR method |
| $\underline{N}=10^{4}$ |  | 100 | 114 | 350 | 153 |
|  |  | 64 | 63 | - | 1,667 |
|  |  | 31 | 54 | 831 | 2,525 |
|  |  | 31 | 46 | 1,005 | 1,559 |
| $\underline{N}=10^{5}$ | - circle | 66 | 13 | 340 | 17 |
|  |  | 67 | 55 | - | 1,457 |
|  | $\bigcirc$ Ofectangle | 37 | 35 | 1,205 | 1,784 |
|  | triangle | 29 | 29 | 1,081 | 747 |
| $\underline{N}=10^{6}$ | - circle | 66 | 3 | 183 | 6 |
|  | - square | 67 | 46 | - | 682 |
|  | $\bigcirc$ - rectangle | 37 | 24 | 884 | 1,614 |
|  | triangle | 28 | 24 | 686 | 459 |
| $\underline{N}=10^{7}$ | - circle | 60 | 2 | 197 | 3 |
|  | \% square | 67 | 40 | - | 671 |
|  | $\bigcirc$ - rectangle | 37 | 13 | 1,453 | 969 |
|  | triangle | 28 | 20 | 541 | 256 |
| $\underline{N}=10^{8}$ | - circle | 59 | 1 | 197 | 7 |
|  | - ${ }_{\text {U }}$ square | 67 | 5 | - | 466 |
|  | $\bigcirc$ - rectangle | 37 | 2 | 1,488 | 155 |
|  | triangle | 28 | 3 | 655 | 296 |

Table 1: Comparison between the Projected Gradient (PG) and the Inexact Restoration (IR) methods applied to Instances 1.1 to 1.4 of Example 1 and Instances 2.1 to 2.4 of Example 2 varying $\underline{N} \in\left\{10^{4}, 10^{5}, \ldots, 10^{8}\right\}$.


Figure 1: Solutions to the four considered instances of Examples 1 (circular classifier) and 2 (ellipsoidal classifier). In the case of instances (a) and (b) both solutions coincide.


Figure 2: Samples size as the method evolves in the four instances of Example 2 (ellipsoidal classifier).
while the instances in which the IR method is more efficient are the instances with very large values of $\underline{N}$. If we consider the instances with $\underline{N}=10^{8}$ only, the IR method is at least twice faster than the PG method and, in some cases, it is ten times faster and up to more than fifty times faster than the PG method. The cases in which the IR method appears to be extremely fast are cases in which samples with size $\underline{N}$ are required by the method (i.e. used to evaluate the objective function) only once or twice. Figure 2 displays, for the four instances of Example 2 with $\underline{N}=10^{8}$, the evolution of the size of the samples as the method evolves. It is very clear that, as desired, big samples are used only when the method approaches a solution.

### 4.4 A portfolio optimization problem

Assume that, using simulation, we are able to produce as many scenarios of future prices of a set of assets as we desire. One of the strategies used to compose a portfolio consists of maximizing the expected return minus a multiple of the estimated standard deviation. This objective function
satisfies adequately the requirement of combining future profits with a cautious attention to risk. Moreover, variations on the weight given to the standard deviation produce different solutions and the corresponding pairs return/deviation represent the classical Pareto curve of non-dominated decisions, among which the investor should chose the one that best fits his/her desires, perhaps based on non-quantifiable criteria. Fixing the weight of the standard deviation we are in presence of a problem of the type analyzed in this paper. As in previous examples, we are interested in the optimization of the original problem, employing a moderate number of scenarios at the portfolio trials that are far from a solution.

Let us consider a portfolio consisting of $n$ assets and assume that $\xi \in \mathbb{R}^{n}$ is a continuous random variable that represents a certain scenario of future prices of the assets. Thus, if $x_{j}$ is the present value of $z$ units of the asset $j$, then $\xi_{j} x_{j}$ is the inflation-discounted prize of $z$ units of the asset $j$ one period later under the scenario represented by $\xi$. The expected return of a portfolio $x \in \mathbb{R}^{n}$ is $E\left(\sum_{i=1}^{n} \xi_{j} x_{j}\right)$; while its standard deviation is given by $\sigma\left(\sum_{i=1}^{n} \xi_{j} x_{j}\right)$. Including the budget constraint $\sum_{j=1}^{n} x_{j}=1$ and non-negativity constraints $x \geq 0$, the objective of an investor could be

$$
\text { Maximize } \gamma\left[E\left(\sum_{j=1}^{n} \xi_{j} x_{j}\right)\right]+(1-\gamma)\left[-\sigma\left(\sum_{j=1}^{n} \xi_{j} x_{j}\right)\right] \text { subject to } \sum_{j=1}^{n} x_{j}=1 \text { and } x \geq 0
$$

where $\gamma \in[0,1]$ is a given parameter. As in the previous subsection, we assume that a large integer $\underline{N}>0$ and small tolerance $\varepsilon_{\text {opt }}>0$ are given and that the objective is to find a point $x^{*} \in \Omega$ such that

$$
\left\|P_{\Omega}\left(x^{*}-\nabla f_{N}\left(x^{*}\right)\right)-x^{*}\right\|_{\infty} \leq \varepsilon_{\mathrm{opt}}
$$

for some $N \geq \underline{N}$, where

$$
\begin{gathered}
\Omega=\left\{x \in \mathbb{R}^{n} \mid \sum_{j=1}^{n} x_{j}=1 \text { and } x \geq 0\right\} \\
f_{N}(x)=\gamma\left[\sum_{j=1}^{n} \bar{\xi}_{j} x_{j}\right]+(1-\gamma)\left[-\sqrt{\frac{1}{N} \sum_{i=1}^{N}\left(\sum_{j=1}^{n}\left[\xi_{i}\right]_{j} x_{j}-\sum_{j=1}^{n} \bar{\xi}_{j} x_{j}\right)^{2}}\right]
\end{gathered}
$$

and

$$
\bar{\xi}=\frac{1}{N} \sum_{i=1}^{N} \xi_{i} \in \mathbb{R}^{n}
$$

For the scenarios' generation, following [7], we considered the date range from January 14, 2013, to January 12, 2015, that, for the NYSE exchange, implies in 503 observable days. Companies were taken from the listing of the NYSE exchange available at http://www.nasdaq.com/ screening/company-list.aspx that contains 3300 companies (accessed on January 13, 2015). We considered the first 100 listed companies excluding those with less than 503 observations and those that contain the character ${ }^{\text {^ }}$ in its symbol. We took historical data (adjusted close value) from Yahoo! Finance using the web tool available at http://finance.jasonstrimpel.com,
that is able to download multiple stock data series on one spreadsheet. We named this real data $\hat{D} \in \mathbb{R}^{503 \times 100}$. To obtain a matrix $D \in \mathbb{R}^{503 \times n}$ with $n=101$, we added a last column given by $(1,1, \ldots, 1)^{T} \in \mathbb{R}^{503}$ corresponding to a risk-free asset. Using matrix $D$ as a source, we computed the matrix of growing factors $G \in \mathbb{R}^{502 \times n}$ whose elements $g_{i j}$ are given by $g_{i j}=d_{i+1, j} / d_{i j}$ for $i=1, \ldots, 502$ and $j=1, \ldots, n$. Each scenario $\xi \in \mathbb{R}^{n}$ will represent a scenario that corresponds to applying 250 random growing factors to the unitary vector $e=(1, \ldots, 1)^{T} \in \mathbb{R}^{n}$ (note that 250 observable days corresponds to a calendar year).

For comparison purposes, it is worth noting that the asset with largest expected profit is asset $j=21$ with $\bar{\xi}_{21} \approx 1.8601$ (when we consider $10^{6}$ as well as $10^{7}$ scenarios). Therefore, when $\gamma=1$, the solution to the problem at hand (that in this case reduces to a linear programming problem) is given by investing all the budget on asset $j=21$. On the other hand, when $\gamma=0$, the considered problem reduces to minimizing the standard deviation of the portfolio. In this case, the solution is given by investing all the budget in the risk-free asset $j=101$, that behaves exactly in the same way for all possible scenario. It is worth noting that the objective function is not differentiable at points where the standard deviation vanishes. For that reason, in the numerical experiments we considered $\gamma \in\{1,0.9, \ldots, 0.4\}$. We also considered that $\varepsilon_{\text {opt }}=10^{-4}$ and $\underline{N} \in\left\{10^{6}, 10^{7}\right\}$. Table 2 shows the results. In all cases both methods (PG and IR) satisfied the stopping criterion and found equivalent solutions. In the case $\underline{N}=10^{6}$, the Inexact Restoration method used approximately half of the computational effort required by the Projected Gradient method; while in the larger case ( $\underline{N}=10^{7}$ ), the IR method was almost one order of magnitude faster than the PG method. Figure 3 shows the value of the expected return versus the standard deviation in the solutions found for the different values of the parameter $\gamma$.

## 5 Concluding remarks

Perhaps the main theoretical meta-optimization consequence of this work is that the key word in Inexact Restoration is "inexactness" instead of "feasibility" or even "restoration". In [18] this consequence was timidly suggested as we still relied on a reformulation of the minimization of $f(x)$ using an auxiliary constraint $z=f(x)$. Here we discarded such interpretation at all and, as a matter of fact, our present preference is the reverse one: Instead of thinking inexactness through the model of infeasibility, we believe that it is more useful to think infeasibility through the model of inexactness. For example, in nonlinear programming, the evaluation of $f(x)$ at an infeasible point may be considered an inexact evaluation of $f(x)$ and we conjecture that many useful ideas may come from the encounter between the IR framework and the treatment of inexactness in mathematical problems.

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Table 2: Comparison between the Projected Gradient (PG) and the Inexact Restoration (IR) methods applied to the portfolio optimization problem with $\underline{N} \in\left\{10^{6}, 10^{7}\right\}$ and $\gamma \in$ $\{1,0.9, \ldots, 0.4\}$.


Figure 3: Expected return versus standard deviation for $\gamma$ varying from 1 to 0.4 .
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