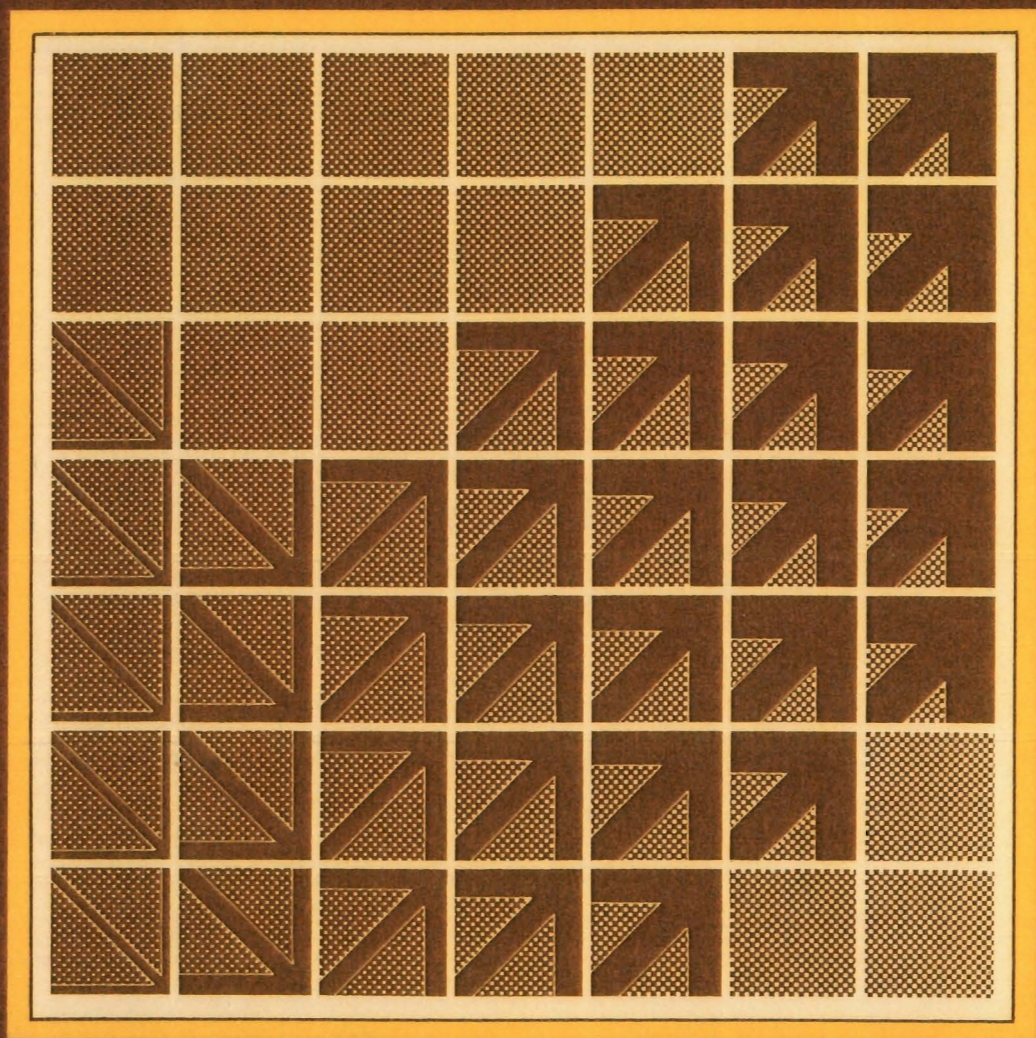


IIASA PROCEEDINGS SERIES

Nonsmooth Optimization

Proceedings of a IIASA Workshop
March 28-April 8, 1977

Claude Lemarechal and Robert Mifflin, Eds.



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Volume 3

**Nonsmooth
Optimization**

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NONSMOOTH OPTIMIZATION

Proceedings of a IIASA Workshop,
March 28 - April 8, 1977

CLAUDE LEMARECHAL
ROBERT MIFFLIN
Editors



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PREFACE

Task 2 of the System and Decision Sciences Area, Optimization, is a central methodological tool of systems analysis. It is used and needed by many Tasks at IIASA, including those of the Energy Systems and the Food and Agriculture Programs. In order to deal with large-scale applications by means of decomposition techniques, it is necessary to be able to optimize functions that are not differentiable everywhere. This is the concern of the subtask Nonsmooth Optimization. Methods of nonsmooth optimization have been applied to a model for determining equilibrium prices for agricultural commodities in world trade. They are also readily applicable to some other IIASA models on allocating resources in health care systems.

This volume is the result of a workshop on Nonsmooth Optimization that met at IIASA in the Spring of 1977. It consists of papers on the techniques and theory of nonsmooth optimization, a set of numerical test problems for future experimentation, and a comprehensive research bibliography.

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INTRODUCTION

This volume is the result of a workshop on nonsmooth optimization held at IIASA from March 28 to April 8, 1977. The designation "First World Conference on Nonsmooth Optimization", proposed in jest by one of the participants after noting that there were only nine others in the room with him, is, however, appropriate because of the various countries represented, and because the number of scientists doing research in this field at that time was rather small.

The small number of participants, and the workshop's unusual length, made it possible to achieve a substantial exchange of information. Each morning (three working hours) was devoted to the talk of one participant, who therefore could present his work quite thoroughly. During the afternoons, discussions took place on related topics, such as: systems of inequalities, constrained problems, test problems and numerical experiments, smooth approximation of nonsmooth functions, optimization with noise, direction-finding procedures and quadratic programming, line searches, general decomposition, However, this workshop format would have been a failure were it not for the fact that everyone was alert and active even when not "in the spotlight". We are very grateful to all the participants, who contributed to the success of the workshop by their numerous questions and interruptions during both the formal and informal presentations.

This workshop was held under the name Nondifferentiable Optimization, but it has been recognized that this is misleading, because it suggests "optimization without derivatives". As we view it, nonsmooth optimization (NSO) is concerned with problems having functions for which gradients exist almost everywhere, but are not continuous, so that the usual gradient-based methods and results fail. The contents of these Proceedings should convince the reader of the importance of being able to compute (generalized) gradients in NSO.

We have adopted the following topical classification for the papers: subgradient optimization (three papers), descent methods (four papers), and field of applicability (one paper).

The first paper, by B.T. Poljak, exhaustively surveys the Soviet work on subgradient optimization done since 1962. For this method he gives the

most important results obtained and the various extensions that have been developed.

J.L. Goffin studies rates of convergence in subgradient optimization. He shows under which conditions linear convergence can be obtained and provides bounds on the best possible rate of convergence. These bounds are given in terms of condition numbers that do not depend on derivative continuity.

The paper by R. Chaney and A.A. Goldstein addresses the question: What is the most general framework for the method of subgradient optimization to be applicable and convergent? Hence, they present the method in an abstract setting and study the minimal hypotheses required to ensure convergence.

One of the important conclusions of this workshop has been that nonsmooth optimization and nonlinear programming (NLP) are, in fact, equivalent fields. It was known that NLP is contained in NSO via exact penalty function methods, but B.N. Pshenichnyi's paper demonstrates the reverse containment via feasible direction methods.

In his paper, C. Lemarechal describes, in a unified setting, descent methods developed recently in Western countries. He also provides ideas for improvement of these methods.

Many methods for solving constrained optimization problems require the repeated solution of constrained least squares problems for search direction determination. An efficient and reliable algorithm for solving such subproblems is given in the paper by R. Mifflin.

The paper by P. Wolfe is concerned with line searches. He gives an APL algorithm that effectively deals with the issues of when to stop a line search with a satisfactory step size and how to determine the next trial step size when the stopping criterion is not met.

The last paper, by J. Gauvin, studies the differential properties of extremal value functions. This is important for the application of various decomposition schemes for solving large-scale optimization problems, because these approaches require the solution of nonsmooth problems involving extremal-value functions, and in order to guarantee convergence we need to know whether certain "semismoothness" conditions (such as Lipschitz continuity) are satisfied.

We then give four nonsmooth optimization test problems. They were selected because they are easy to work with and because they are representative both of the field of applicability and of the range of difficulty of

NSO. Problems 1 and 3 are examples of minimax problems and are not very difficult. Problem 2 is a nonconvex problem coming from a well-known NLP test problem, and problem 4 involves a piecewise-linear function. The last two are sufficiently difficult to slow down considerably the speed of convergence of any of the NSO methods we know of.

We conclude this volume with a large NSO bibliography. It was compiled by the participants and is an update of the bibliography given in *Mathematical Programming Study 3*. We wish to thank D.P. Bertsekas, V.F. Demjanov, M.L. Fisher, and E.A. Nurminskii for the items they communicated to us.

On behalf of all the participants we would like to acknowledge IIASA's generous support and to thank I. Beckey, L. Berg, A. Fildes, and G. Lindelof for their optimal organizational contributions, which led to a smooth-running workshop.

We are especially indebted to M.L. Balinski who was instrumental in establishing a Nonsmooth Optimization group at IIASA and who spent much of his time and energy to secure a truly international participation at this workshop.

C. Lemarechal
R. Mifflin

SUBGRADIENT METHODS: A SURVEY OF SOVIET RESEARCH

B. T. Poljak

This paper reviews research efforts by Soviet authors concerning the subgradient technique of nondifferentiable minimization and its extensions. It does not cover the works based on the concept of steepest descent (by V.F. Demjanov, B.N. Pshenichnyi, E.S. Levitin, and others) or on the use of a specific form of the minimizing function (for example minimax techniques). The paper essentially uses the review by N.Z. Shor [1]. The theorems given below are mostly simplified versions of results shown in the original papers.

1. THE SUBGRADIENT METHOD

Let $f(x)$ be a convex continuous function in the space R^n . A vector $\partial f(x) \in R^n$ is called its subgradient at the point x , if it satisfies the condition

$$f(x+y) \geq f(x) + (\partial f(x), y) \quad , \quad \forall y \in R^n \quad . \quad (1)$$

A subgradient exists (although, generally speaking, it may be not unique) for all $x \in R^n$. If $f(x)$ is differentiable, the subgradient is unique and coincides with the gradient $\partial f(x) = \nabla f(x)$. The rules of subgradient calculation for various types of functions are well known [2,3]. In particular, with $f(x) = \max_{1 \leq i \leq m} f_i(x)$ where $f_i(x)$ are convex differentiable functions, it is true that

$$\partial f(x) = \sum_{i \in I(x)} \alpha_i \nabla f_i(x) \quad , \quad \alpha_i \geq 0 \quad ,$$

$$\sum_{i \in I(x)} \alpha_i = 1 \quad , \quad I(x) = \{i : f_i(x) = f(x)\}$$

(for instance one may take $\partial f(x) = \nabla f_i(x)$ where $i \in I(x)$ is arbitrary).

The subgradient minimization method for $f(x)$ on R^n is an iterative process of the form

$$x_{k+1} = x_k - \gamma_k \partial f(x_k) / \|\partial f(x_k)\| \quad (2)$$

where $\gamma_k \geq 0$ is a step size. For differentiable functions this method coincides with the gradient one. The major difference between the gradient and the subgradient methods is that, generally speaking, the direction $-\partial f(x_k)$ is not a descent direction at the point x_k ; i.e., the values of $f(x_k)$ for nondifferentiable functions do not decrease monotonically in the method (2).

The subgradient method was developed in 1962 by N.Z. Shor and used by him for solving large-scale transportation problems of linear programming [4]. Although published in a low-circulation publication, this pioneering work became widely known to experts in the optimization area in the USSR. Also of great importance for the propagation of nondifferentiable optimization concepts were the reports by the same author presented in a number of conferences in 1962-1966.

Publication of papers [5,6,7] giving a precise statement of the method and its convergence theorems may be regarded as the culmination of the first stage in developing subgradient techniques.

Let us get down to describing the basic results concerning the subgradient method. As is known, the gradient method for minimization of smooth functions employs the following ways to regulate the step-size:

$$\gamma_k = \alpha \|\partial f(x_k)\| \quad ,$$

i.e.

$$x_{k+1} = x_k - \alpha \nabla f(x_k) \quad , \quad 0 < \alpha < \bar{\alpha}$$

(the ordinary gradient method);

$$\gamma_k = \arg \min_{\gamma} f(x_k - \gamma \partial f(x_k) / \|\partial f(x_k)\|)^\dagger$$

(the steepest descent method).

Simple examples may be constructed to show that neither of these methods converges in nondifferentiable minimization; this necessitates the construction of new principles of selecting the step size. Consider the major ones. Hereinafter we shall assume $f(x)$ to be convex and continuous and denote $f^* = \inf f(x)$ and $X^* = \text{Arg min } f(x)$.

(a) $\gamma_k = \gamma > 0$. This constant-step method was suggested in [4]. The simplest example, $f(x) = |x|$, $x \in \mathbb{R}^1$, explicitly proves that this method does not converge. One may show, however, that it gives a solution "with an accuracy of γ ".

Theorem 1 [4]

Let X^* be nonempty. Then for any $\delta > 0$ there exists $\bar{\gamma} > 0$ such that in the method (2) with $\gamma_k = \gamma$, $0 < \gamma < \bar{\gamma}$ we have $\liminf f(x_k) < f^* + \delta$.

Reference [4] has described the following way of step-size regulation resting upon this result, although it has not been entirely formalized. A certain $\gamma > 0$ is chosen and the computation is made with $\gamma_k = \gamma$ until the values of $f(x_k)$ start to oscillate about a certain limit. After this γ is halved and the process is repeated.

(b) The sequence γ_k is chosen a priori regardless of the computation results and satisfies the condition

$$\sum_{k=0}^{\infty} \gamma_k = \infty, \quad \gamma_k \rightarrow 0. \quad (3)$$

This way of choosing the step-size has been suggested in [5] and [6] independently.

[†]Hereafter $\arg \min_{\gamma} \rho(\gamma)$ will mean an arbitrary minimum point of the function $\rho(\gamma)$, $\text{Arg min}_{\gamma} \rho(\gamma)$ is the set of all minimum points.

Theorem 2 [5,6]

In the method (2), (3) $\liminf f(x_k) = f^*$. If X^* is nonempty and bounded then $\rho(x_k, X^*) \rightarrow 0$, where

$$\rho(x, X^*) = \min_{x^* \in X^*} \|x - x^*\| \quad .$$

It is clear that in the general case the method (2), (3) cannot converge faster than γ_k tends to zero. In particular, this method never converges at the rate of geometrical progression or at the rate

$$O(k^{-s}), \quad s > 1 \quad .$$

(c) In certain cases the value of f^* is known. For instance, if

$$f(x) = \sum_{i=1}^m f_i(x)_+ \quad ,$$

where $f_i(x)$ are convex functions,

$$f_i(x)_+ = \max \{0, f_i(x)\} \quad ,$$

and the system of inequalities $f_i(x) \leq 0$ $i = 1, \dots, m$ is solvable, then X^* is the set of solutions of this system and $f^* = 0$. Then one may take

$$\gamma_k = \lambda \frac{(f(x_k) - f^*)}{\|\partial f(x_k)\|} \quad , \quad 0 < \lambda < 2 \quad . \quad (4)$$

In solving systems of inequalities the method (3), (4) coincides with the known relaxation method of Kaczmarz, Agmon, Motzkin, Schoenberg, and Eremin [8]. The method for general problems of nonsmooth function minimization has in essence been suggested by I.I. Eremin [9] and systematically developed in [10].

Theorem 3 [9,10]

Let x^* be the unique minimum point for $f(x)$. Then in the method (2), (4) $x_k \rightarrow x^*$. If the condition

$$f(x) - f^* \geq \ell \|x - x^*\| \quad , \quad \ell > 0 \quad (5)$$

holds, the method converges with the rate of a geometrical progression.

The advantages of the method (2), (4) are the simplicity of selecting the step size (since no auxiliary problems should be solved and no characteristics of $f(x)$ other than f^* should be known) and its applicability, since for a smooth strongly convex $f(x)$ the method also converges with the rate of a geometrical progression [10]. Reference [10] has shown a way to modify the technique when f^* is unknown.

(d) N.Z. Shor [11] has suggested an essentially different method for choosing γ_k :

$$\gamma_k = \gamma_0 q^k \quad , \quad 0 < q < 1 \quad . \quad (6)$$

Note that the condition (3) is not satisfied for (6).

Theorem 4 [11,12,13]

Let the condition

$$(\partial f(x), x - x^*) \geq \ell \|\partial f(x)\| \|x - x^*\| \quad , \quad \ell > 0 \quad (7)$$

hold. Then there exists a pair \bar{q} (which depends on ℓ) and $\bar{\gamma}$ (which depends on $\|x_0 - x^*\|, \ell$) such that with $1 > q \geq \bar{q}$, $\gamma_0 \geq \bar{\gamma}$ in the method (2), (6) we have

$$\|x_k - x^*\| \leq C(q, \gamma_0) q^k \quad .$$

The relationship of $\bar{q}(\ell)$ and $\bar{\gamma}(\|x_0 - x^*\|, \ell)$ may be expressed explicitly. However, practical implementation of the method (2), (6)

faces difficulties because generally the values of ℓ and $\|x_0 - x^*\|$ are unknown.

The above results prove that the convergence rate for any of the step-size regulating rules is linear at best. The denominator of the geometrical progression for the ill-conditioned problems (i.e. for functions with greatly extended level sets) is near unity. Thus the convergence rate of all the versions of the subgradient method may be rather poor.

2. ACCELERATING CONVERGENCE OF THE SUBGRADIENT METHOD

One of the reasons why the subgradient method converges so slowly lies in its Markov nature. The subsequent iteration makes no use of the information obtained at the previous steps. The major concept of all techniques for accelerating convergence is the use of this information (i.e. the values $f(x_i)$, $\partial f(x_i)$, $i=0, \dots, k-1$).

The first methods of this type were those developed by Kelley and by Cheney and Goldstein [14,15], based on piecewise-linear approximation of the function. An original technique suggested in [16] and [17] independently made use only of the values $\partial f(x_i)$. Let M_k be a polyhedron in R^n in which the minimum point is localized after k iterations. Then for an arbitrary $x_{k+1} \in M_k$ one may take

$$M_{k+1} = M_k \cap \{x: (\partial f(x_{k+1}), x - x_{k+1}) \leq 0\} .$$

If the center of gravity M_k is taken as the point x_{k+1} one may show [17] that for the volume V_k of the polyhedron M_k the following expression holds:

$$V_{k+1} \leq [1 - (1 - \frac{1}{N-1})^N] V_k ,$$

where N is the dimension of the space. Thus for problems of any dimension

$$V_k \leq V_0 q^k , \quad q = 1 - e^{-1} .$$

In other words, the method converges with the rate of geometrical progression, the denominator being independent of both the properties of the function and the dimension of the space. This result is mostly of theoretical interest since the auxiliary problem of finding the center of gravity for the polyhedron is very difficult to solve. References [17,18,19] give a number of modified techniques in which a simpler auxiliary problem is solved at each step. References [18,19] have also shown that the center-of-gravity technique is optimal in a certain sense. Roughly speaking, there is no algorithm that uses the same information and provides better convergence. A similar result for algorithms that uses only the values $f(x_i)$ rather than $\partial f(x_i)$ is given in Reference [20].

These methods [14,15,16,17,18,19,20] have been developed independently of the subgradient method. Let us turn now to the algorithms obtained as a direct extension of (2).

Let the value of f^* be known. Reference [10] suggests the following technique. At the k^{th} step the quadratic programming problem is solved:

$$\begin{aligned} \min \quad & \|x - x_k\|^2 \\ & f(x_k) + (\partial f(x_k), x - x_k) \leq f^* \\ & f(x_{k-1}) + (\partial f(x_{k-1}), x - x_{k-1}) \leq f^* \\ & \text{-----} \\ & f(x_{k-m}) + (\partial f(x_{k-m}), x - x_{k-m}) \leq f^* \end{aligned} \tag{8}$$

and its solution is taken as x_{k+1} . The value $m \geq 0$ is arbitrary and may depend on k . In particular, with $m = 0$ the method is identical to the subgradient algorithm (2), (4) with $\lambda = 1$. If m is smaller, then instead of (8) it is more convenient to solve the dual problem

$$\min \left\{ \frac{1}{2} \left\| \sum_{i=k-m}^k \lambda_i \partial f(x_i) \right\|^2 - \sum_{i=k-m}^k \lambda_i (f(x_i) - f^* - (\partial f(x_i), x_i - x_k)) : \lambda_i \geq 0 \right\} \tag{9}$$

and, denoting its solution as λ_i^k , to obtain

$$x_{k+1} = x_k - \sum_{i=k-m}^k \lambda_i^k \partial f(x_i) \quad . \quad (10)$$

This form of stating the method shows that it is quite closely related to the steepest descent, the conjugate subgradient, and other methods. Reference [10] shows that the method (8) converges at least no more slowly than the subgradient method (2), (4) with $\lambda = 1$. Moreover, if $f(x)$ is a piecewise-linear function with a non-singular minimum and $m \geq N$, the method is finite. The latter property is a great advantage of the method, although, firstly, one needs to know f^* and, secondly, for large m the auxiliary problem 8) or (9,10) is rather laborious.

In a number of his works [21,22,23,24,25,26] N.Z. Shor has suggested that space metric be transformed at each iteration to accelerate convergence of the subgradient method. A general algorithm for such an approach is in the following: let $s \in \mathbb{R}^n$, $\|s\| = 1$, $\alpha \geq 0$. Then a linear operator $R_\alpha(s)$ such that

$$R_\alpha(s)x = x + (\alpha - 1)ss^t x \quad (11)$$

is referred to as the space-dilation operator acting in the direction s with a coefficient α . It is clear that $R_\alpha(s)s = \alpha s$, and for x orthogonal to s , $R_\alpha(s)x = x$. By making the transformations of the space $R_{\alpha_k}(s_k)$ at each step, computing subgradients for the new variables, and then transforming back to the original variables, we shall have the method

$$x_{k+1} = x_k - \gamma_k P_k P_k^t \partial f(x_k) \quad , \quad (12)$$

$$P_k = P_{k-1} R_{\alpha_k}^{-1}(s_k) \quad , \quad P_0 = I \quad .$$

N.Z. Shor has considered two feasible versions for selecting directions s_k . The first [21,22,24,25] provides dilation of the space in the direction of a subsequent subgradient, i.e.

$s_k = \partial f(x_k)$. Let us put down this method in a symmetrical form suggested by V.A. Skokov [27]:

$$\begin{aligned} x_{k+1} &= x_k + \gamma_k p_k, \quad p_k = -H_k \partial f(x_k), \\ H_k &= H_{k-1} + \left(1 - \frac{1}{\alpha_k^2}\right) \frac{p_{k-1}^t p_{k-1}}{(\partial f(x_{k-1}), p_{k-1})}, \quad H_0 = I. \end{aligned} \quad (13)$$

Obviously, the matrices H_k are symmetric. In the above papers a number of ways to select the parameters γ_k and α_k are considered. In particular if f^* and the constant $M > 1$ are known in the inequality

$$(\partial f(x), x - x^*) \leq M(f(x) - f^*), \quad (14)$$

one may choose

$$\begin{aligned} \gamma_k &= \left(\frac{2M}{M+1}\right) \left(\frac{f(x_k) - f^*}{\|\partial f(x_k)\|^2}\right), \\ \alpha_k &= \alpha = \frac{M+1}{M-1} \end{aligned} \quad (15)$$

Theorem 5 [22]

Let $f(x)$ be convex and continuous and let the condition (14) hold. Then the algorithm (13), (15) converges geometrically with the rate $\alpha^{-1/N}$:

$$\liminf \alpha^{k/N} (f(x_k) - f^*) < \infty.$$

A result for the extreme case $\alpha_k = \infty$ and

$$\gamma_k = \frac{2(f(x_k) - f^*)}{\|p_k\|^2} \quad (16)$$

is also known.

Theorem 6 [22,27]

If $f(x)$ is quadratic, then in the method (13), (16) $x_N = x^*$, and $H_N = 0$. In other words, for a quadratic case the method (16) is finite and coincides with one of the known orthogonalization methods for solving systems of linear equations.

The second version of the method [23,24,25,26] provides space dilation in the direction of the difference of two sequential subgradients, i.e. $s_k = \partial f(x_k) - \partial f(x_{k-1})$. In the statement of [27] the method takes the form

$$x_{k+1} = x_k + \gamma_k p_k, \quad p_k = -H_k \partial f(x_k),$$

$$e_k = \partial f(x_k) - \partial f(x_{k-1}), \quad (17)$$

$$H_k = H_{k-1} - \left(1 - \frac{1}{\alpha_k}\right) \frac{H_{k-1} e_k e_k^t H_{k-1}}{(H_{k-1} e_k, e_k)}, \quad H_0 = I.$$

Unlike all other versions of the subgradient method, the step size is chosen from the condition of the steepest descent:

$$\gamma_k = \arg \min_{\gamma} f(x_k + \gamma p_k). \quad (18)$$

In practical computations the value α_k was taken equal to 2 or 3 and further increase of α_k did not affect the convergence rate. References [23,24,25,26] describe certain conditions that guarantee the convergence of the algorithm. Its rate of convergence has not been given enough attention.

Consider an extreme version of the method, $\alpha_k = \infty$.

Theorem 7 [26,27]

Let $\alpha_k = \infty$ and $f(x)$ be quadratic. Then in the method (17), (18)

$$x_N = x^*, \quad H_N = 0.$$

The method (17), (18) in this case is reduced to one of the versions of the conjugate directions method [28].

Thus, the subgradient methods with space dilation are conceptually close to variable-metric methods used for minimization of smooth functions, and their limiting versions have the property that they are finite for quadratic functions. The important question of the convergence rate for space-dilation methods for essentially nonsmooth functions (e.g. piecewise-linear) remains unanswered.

3. EXTENSIONS

We have thus far discussed unconditional minimization of a convex continuous function on R^N . Now let us concentrate on the potential of the subgradient method in solving more general problems.

(a) *Infinite-Dimensional Space.* Let it be required to find the minimum of the convex continuous functional $f(x)$ in the Hilbert space H . The subgradient $\partial f(x)$ is in this case defined exactly as in (1) and the subgradient method has the same form as (2). It has been proved [6] that the theorem on convergence of the method (2), (3) remains valid even in Hilbert spaces, and the same has been proved [10] for the methods (2), (4), and (8). Nevertheless some of the methods of Section 2 (e.g. the center-of-gravity method) are specific to finite-dimensional spaces.

(b) *Problems with Constraints.* Let us consider the minimization problem with constraints

$$\begin{aligned} \min f(x) \quad , \quad x \in H \\ g(x) \leq 0 \\ x \in Q \quad , \end{aligned} \tag{19}$$

where f and g are continuous convex functionals and Q is a convex closed set. An arbitrary problem of convex programming can be reduced to the form (19). Thus if there are several constraints

$g_i(x) \leq 0, i = 1, \dots, m$ then it can be assumed that

$$g(x) = \max_i g_i(x)$$

or

$$g(x) = \sum_{i=1}^m g_i(x)_+ .$$

The set Q is assumed to be of simple structure; in particular, the problem of finding a projection onto Q has a simple solution. The cases of

$$Q = \{x \in \mathbb{R}^N : a \leq x \leq b\} , \quad Q = \{x \in H : \|x - a\| \leq r\}$$

are typical. Reference [6] has proposed an extension of the sub-gradient method for solution of (19):

$$x_{k+1} = P_Q(x_k - \gamma_k s_k)$$

$$s_k = \begin{cases} \partial f(x_k) / \|\partial f(x_k)\| & \text{if } g(x_k) \leq 0 \\ \partial g(x_k) / \|\partial g(x_k)\| & \text{if } g(x_k) > 0 \end{cases} , \quad (20)$$

where P_Q is the projection operator on Q and convergence of this method has been proved under the condition (3).

(c) *Nonconvex Functions.* The function $f(x)$ to be minimized has thus far been assumed convex. Now let us consider the case of quasiconvex functions (such that the sets $\{x : f(x) \leq C\}$ are convex). In this case the subgradient at the point x_k can be replaced by a vector s_k which is support for the set $\{x : f(x) \leq f(x_k)\}$, i.e. $(s_k, x - x_k) \leq 0$ for all x such that $f(x) \leq f(x_k)$. It can be proved [6] that with this replacement the method retains convergence under the assumptions made in the convex case.

Another generalization of greater importance is using nearly-differentiable functions [25]. A generalized gradient of f at x

is a vector $s = \lim_{k \rightarrow \infty} \nabla f(x_k)$, where $x_k \rightarrow x$ is an arbitrary sequence of points where the gradient exists. For convex continuous $f(x)$ the set of generalized gradients coincides with the set of subgradients. References [25,29] have proved the convergence of the method (2) where the subgradient is replaced by the generalized gradient. A somewhat different important class of weakly convex functions to which the subgradient method can be extended has been studied by E.A. Nurminskii [30,31].

(d) *Non-Unique Minimum.* Let the set X^* of the minimum points of minima of the convex continuous function $f(x)$ on the convex closed set Q consist of more than one point. Then the subgradient minimization method is nonstable, in the sense that for different initial points it can converge to different solutions (and for some variants of the method even that cannot be guaranteed). In a similar way, in the infinite-dimensional case the subgradient method may not converge even if there is a unique minimum solution. This kind of problem can be solved by using a regularization method. A regularizing function $f_1(x)$ is chosen that is strictly convex (in the infinite-dimensional case uniformly convex), e.g. $f_1(x) = \|x\|^2$. There is a unique point x_k minimizing the regularized function $f(x) + \alpha_k f_1(x)$, where $\alpha_k > 0$ is the regularization parameter. Then it can be shown [32] that $x_k \rightarrow x^*$ as $\alpha_k \rightarrow 0$ where x^* is the point in the set of minimum points of $f(x)$ on Q for which $f_1(x)$ is minimal. The subgradient method can then be made stable by using this idea. Consider the method

$$x_{k+1} = P_Q[x_k - \gamma_k (\partial f(x_k) + \alpha_k \partial f_1(x_k))] \quad . \quad (21)$$

In other words, one step of the subgradient method for minimization of the regularized function is made; following this, the regularization parameter is changed.

Theorem 8 [33]

Let $f(x)$ and $f_1(x)$ be convex continuous functions on \mathbb{R}^N , let $f_1(x)$ be strongly convex, let Q be a convex closed and bounded set,

and let $\alpha_k \rightarrow 0, \gamma_k/\alpha_k \rightarrow 0, \sum \gamma_k \alpha_k = \infty, \alpha_k/\alpha_{k+1} = 1 + o(\gamma_k \alpha_k)$ (for instance $\gamma_k = k^{-1/2}, \alpha_k = k^{-\rho}, 0 < \rho < 1/2$). Then in the method (21) $x_k \rightarrow x^*$, where $x^* = \arg \min_{x \in X^*} f_1(x), X^* = \text{Arg} \min_{x \in Q} f(x)$.

(e) *A Continuous Analog of the Subgradient Method.* To the discrete gradient method, $x_{k+1} = x_k - \gamma_k \nabla f(x_k)$, there corresponds the continuous analog $\dot{x} = -\nabla f(x)$. Similarly there exists a continuous version of the subgradient method

$$\dot{x} \in -\partial f(x) \quad . \quad (22)$$

Methods of this type have been in use (without much justification) since the mid-fifties for solution of linear programming problems on analog computers. To prove that (22) converges is a non-trivial matter, however, since (22) is a differential equation with a multivalued discontinuous right-hand side; therefore its solution needs an appropriate definition. The existence of a solution to such an equation also needs special study. References [34,35,36] deal with this subject.

(f) *The Subgradient Method for Finding Saddle Points.* The point

$$x^* \in Q \subset \mathbb{R}^n, y^* \in S \subset \mathbb{R}^m$$

is called a saddle point of the function $\phi(x,y)$ on $Q \times S$ if

$$\phi(x,y^*) \geq \phi(x^*,y^*) \geq \phi(x^*,y)$$

for all

$$x \in Q, y \in S$$

or

$$\min_{x \in Q} \max_{y \in S} \phi(x, y) = \phi(x^*, y^*) = \max_{y \in S} \min_{x \in Q} \phi(x, y) .$$

If the function $\phi(x, y)$ is smooth, then to minimize it the gradient method can be applied:

$$\begin{aligned} x_{k+1} &= P_Q(x_k - \gamma_k \nabla_x \phi(x_k, y_k)) , \\ y_{k+1} &= P_S(y_k + \gamma_k \nabla_y \phi(x_k, y_k)) . \end{aligned}$$

The similar subgradient algorithm for a nonsmooth case has been proposed in [5,37]:

$$\begin{aligned} x_{k+1} &= P_Q(x_k - \gamma_k \partial_x \phi(x_k, y_k)) \\ y_{k+1} &= P_S(y_k + \gamma_k \partial_y \phi(x_k, y_k)) . \end{aligned} \tag{23}$$

The method (23) has been validated in [38]. The function $\phi(x, y)$ is called stable [38] if for all $x^* \in X^*$, $y^* \in Y^*$ ($X^* \times Y^*$ is the set of saddle points of $\phi(x, y)$ on $Q \times S$) one has

$$\text{Arg min}_{x \in Q} \phi(x, y^*) = X^* , \quad \text{Arg max}_{y \in S} \phi(x^*, y) = Y^* .$$

In particular, if $\phi(x, y)$ is strictly convex with respect to x and strictly concave with respect to y , then it is stable.

Theorem 9 [38]

Let $\phi(x, y)$ be continuous on $R^n \times R^m$, convex with respect to x for all $y \in S$, concave with respect to y for all $x \in Q$, and stable, let the sets Q and S be convex, closed and bounded, and let $\gamma_k \rightarrow 0$, $\sum \gamma_k = \infty$. Then in the method (23),

$$\rho(x_k, X^*) \rightarrow 0 , \quad \rho(y_k, Y^*) \rightarrow 0 .$$

Results on the convergence of (23) in infinite-dimensional spaces have been given in [33].

The convergence of the method (23) without stability (e.g. for the Lagrange function in a convex programming problem) remains an open question. For the smooth case and with stability only with respect to x , this question has been answered in [39]. In the general case the method can be modified by regularization [33].

$$\begin{aligned} x_{k+1} &= P_Q [x_k - \gamma_k (\partial_x \phi(x_k, y_k) + \alpha_k \partial f_1(x_k))] \\ y_{k+1} &= P_S [y_k + \gamma_k (\partial_y \phi(x_k, y_k) - \alpha_k \partial f_2(y_k))] \end{aligned} \quad (24)$$

where $f_1(x), f_2(y)$ are strongly convex functions.

In [33] convergence of this method has been proved under the same assumptions on γ_k, α_k as in Theorem 8. No strict convexity-concavity or stability of $\phi(x, y)$ is needed.

(g) *Difference Approximation of the Subgradient.* In a number of problems the subgradients $\partial f(x)$ are inaccessible and only values $f(x)$ at arbitrary points are known. In this case the subgradient $\partial f(x)$ can probably be replaced by its finite-difference approximation, e.g. by the vector

$$s(x, \alpha) = \sum_{i=1}^N \frac{f(x + \alpha e_i) - f(x - \alpha e_i)}{2\alpha} e_i \quad (25)$$

where e_i are coordinate orths; α is the size of the test step. This procedure works, however, only in the smooth case; it may stop at non-minimum points when minimizing a nondifferentiable convex function. Convergence may be obtained by introducing additional smoothing through randomization. One of the simplest methods of this type was proposed by A.M. Gupal [40]:

$$\begin{aligned} x_{k+1} &= x_k - \gamma_k s(y_k, \alpha_k) \\ y_k &= x_k + \alpha_k z_k \end{aligned} \quad (26)$$

where Z_k is a random vector uniformly distributed on a unit sphere and $s(y, \alpha)$ is computed by formula (25). It has been proved [4] that with a certain ratio of the step sizes α_k, γ_k (namely, with $\sum \gamma_k = \infty, \sum \gamma_k^2 < \infty, \alpha_k \rightarrow 0, \gamma_k/\alpha_k \rightarrow 0, \frac{\alpha_k - \alpha_{k+1}}{\alpha_k \gamma_k} \rightarrow 0$) and under some natural assumptions on $f(x)$, this method does converge. In [41] a similar method was used for minimization of discontinuous as well as nondifferentiable functions.

4. THE SUBGRADIENT METHOD IN THE PRESENCE OF NOISE

In many real-world problems the gradient or the subgradient cannot be precisely computed (for instance in system parameter estimation, identification, learning, and pattern recognition [42,43]) because of incomplete data on the function to be minimized which is the expected value of a certain quantity whose distribution law is not exactly known. In other cases the errors are caused by computation errors, experimentation in a real process, etc. In any case, usually we know only an approximate value of the vector $\partial f(x)$, denoted as $\partial F(x)$. The error

$$\xi(x) = \partial F(x) - \partial f(x) \quad (27)$$

may contain both the random and the deterministic components $\eta(x)$ and $a(x)$:

$$M\xi(x) = a(x) \quad , \quad \eta(x) = \xi(x) - a(x) \quad , \quad M\eta(x) = 0 \quad . \quad (28)$$

Then the subgradient method for minimizing $f(x)$ on R^N is of the form

$$x_{k+1} = x_k - \gamma_k \partial F(x_k) \quad . \quad (29)$$

The pioneering effort in the study of stochastic methods of the form (29) was made by Yu.M. Ermol'ev [44,45,46]. His investigations and those of his followers have been summarized in a monograph [47] and a survey [48]. Among other works on optimization methods in the

presence of random noise, the book by V.Ya. Katkovnik [49] is worth mentioning. The methods of type (29) may be regarded as Robbins-Monro stochastic approximation procedures, and the results obtained in the theory of stochastic approximation (e.g. [43,47,50,51]) can be used in their analysis.

Let us cite a simple result on convergence of the algorithm (29). Its modifications and extensions (deterministic errors, nonunique or nonexistent extremum, noise with infinite variance, mean square convergence, etc.) are to be found in [46,47,52,53, 54].

Theorem 10

Let $f(x)$ be convex, and continuous and have a unique minimum point $x^* \in R^N$; suppose the noise $\xi(x)$ is purely random, is independent at different points, and has a mean $M\xi(x) = 0$, and a variance $\sigma^2(x) = M\|\xi(x)\|^2$, and assume the following constraints on the growth of $\partial f(x)$ and $\sigma^2(x)$:

$$\begin{aligned} \|\partial f(x)\|^2 &\leq C_1(1 + \|x - x^*\|^2) \quad , \\ \sigma^2(x) &\leq C_2(1 + \|x - x^*\|^2) \quad , \end{aligned} \tag{30}$$

Let γ_k satisfy the condition

$$\sum_{k=0}^{\infty} \gamma_k = \infty \quad , \quad \sum_{k=0}^{\infty} \gamma_k^2 < \infty \quad . \tag{31}$$

Then in the method (29), $x_k \rightarrow x^*$ with a probability 1.

As for convergence rate, it can be shown [47,55] that if the condition $f(x) \geq f^* + \ell\|x - x^*\|^2$, $\ell > 0$ or $f(x) \geq f^* + \ell\|x - x^*\|$, $\ell > 0$ is valid, and given the choice $\gamma_k = \gamma/k$ ($\gamma > 0$ being large enough), then the rate of decrease of order $O(k^{-1})$ can be guaranteed for $\|x_k - x^*\|^2$ in some probabilistic sense. This is, generally speaking, the highest possible rate of convergence. Thus for

$$f(x) = x^2, x \in \mathbb{R}^1, M\xi^2(x) = \sigma^2 > 0, \quad ,$$

the iterative process (29) with any method of selecting the step size cannot decrease the value of Mx_k^2 faster than $O(k^{-1})$ [56].

On the contrary, for the functions satisfying the condition (5), and with constrained noise

$$\|\xi(x)\| \leq c < \ell \quad (32)$$

and a step-size rule of type (6), there is geometric convergence.

Theorem 11 [57]

Let $f(x)$ be convex and continuous and the conditions (5) and (32) hold. Then for any x_0 there are γ_0 and $q < 1$ such that with $\gamma_k = \gamma_0 q^k$ the estimate $\|x_k - x^*\| \leq \|x_0 - x^*\| q^k$ is valid for the method (29).

The extensions of the subgradient method in Section 3 are, as a rule, applicable to problems in the presence of noise. Thus a stochastic analog of the method (20) has been proposed [58]. E.A. Nurminskii has applied the stochastic subgradient method to a class of non convex functions [59]. The iterative regularization method (21) in the presence of random noise has been discussed in Reference [60]. Finally, the behavior of the subgradient algorithm for finding saddle points (23) in the stochastic case has been studied in Reference [44,45,46,47] and in combination with the regularization method (an analog of (24)), in Reference [50].

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Abbreviations of Russian Journal Titles

- EMM - Ekonomika i Matematicheskie Metody (Economics and Mathematical Methods)
- *Kibern. - Kibernetika (Kiev)
- *DAN - Doklady Akademii Nauk SSSR (Soviet Mathe. Doklady)
- *ZVMMF - Zurnal Vycislitel'noi Matematiki i Matematicheskoi Fiziki (Journal of Computational Mathematics and Mathematical Physics)
- *AT - Avtomatika i Telemekhanika (Automation and Remote Control)

*These journals are translated into English.

NONDIFFERENTIABLE OPTIMIZATION AND THE RELAXATION METHOD*

J. L. Goffin

1. INTRODUCTION

The relaxation method for solving systems of linear inequalities, as defined by Agmon [1] and Motzkin and Schoenberg [7], is closely connected to the relaxation method for solving systems of linear equalities. The relationship was made precise in Agmon's paper.

Subgradient optimization is a technique that attempts to solve the problem of maximizing a general, maybe nondifferentiable, concave function (or minimizing a convex function). One of the main expectations from subgradient optimization is that it could be used to solve some large-scale problems; computational results reported in [5] and [6] gave some credence to that idea. Subgradient optimization is closely related to the relaxation method for solving systems of linear inequalities, which has been used with success on some very large-scale problems with special structure: this seems to justify a closer look at subgradient optimization.

In this paper we will make explicit the relationship between the relaxation method for linear inequalities and subgradient optimization. The speed of convergence of both methods depends on condition numbers which have been defined in [3] and [4]. It will be shown that the two theorems on convergence are almost identical.

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2. THE RELAXATION METHOD AND SUBGRADIENT OPTIMIZATION

Let $\langle \alpha^i, x \rangle + \beta^i \geq 0$, $i \in I$ be a finite system of linear inequalities where $\alpha^i \in \mathbb{R}^n$, $\beta^i \in \mathbb{R}$, $x \in \mathbb{R}^n$. (2.1)

Let $P = \{x \in \mathbb{R}^n: \langle \alpha^i, x \rangle + \beta^i \geq 0, \forall i \in I\}$, the solution set.

Define the function f_1 by:

$$f_1(x) = \text{Min} \{ \langle \alpha^i, x \rangle + \beta^i, i \in I \}$$

and let f_1^* be the maximum of $f(x)$. It is clear that

$$f_1^* > 0 \Rightarrow \dim P = n ;$$

$$f_1^* < 0 \Rightarrow P \text{ is empty .}$$

If we define $f(x) = \text{Min} \{0, f_1(x)\}$, it follows that P is not empty if and only if $\text{Max}_{x \in \mathbb{R}^n} f(x) = 0$, and that P is the set of points on

which f assumes its maximum value. Furthermore, P is empty if and only if $\text{Max}_{x \in \mathbb{R}^n} f(x) = \text{Max}_{x \in \mathbb{R}^n} f_1(x) = f_1^* < 0$; the set of points

on which f assumes its maximum value has been defined as the Chebyshev solution to the infeasible system of linear inequalities.

If we let $a^i = \frac{\alpha^i}{\|\alpha^i\|}$, $b^i = \frac{\beta^i}{\|\alpha^i\|}$ (we assume that $\alpha_i \neq 0 \forall i \in I$), where $\|\cdot\|$ means the Euclidean norm, then the system

$$\langle a^i, x \rangle + b^i \geq 0, \quad i \in I \quad (2.2)$$

is equivalent to (2.1).

Let $w_1(x) = \text{Min} \{ \langle a^i, x \rangle + b^i, i \in I \}$ and $w(x) = \text{Min} \{ w_1(x), 0 \}$; also let $w^* = \text{Max}_x w(x)$.

Clearly $\{x \in \mathbb{R}^n: w(x) = w^*\} = P$ provided that P is not empty (and $w^* = 0$). If P is empty then $w^* < 0$, and the set of x which maximizes $w(x)$ is not necessarily the same as the set of x which

maximizes $f(x)$ (it is clear though that $f^* < 0$ iff $w^* < 0$). The functions f, f_1, w, w_1 are all concave and piecewise-linear.

A "subgradient" set can be defined for each x and each concave function, say w_1 :

$$\partial w_1(x) = \{v \in \mathbb{R}^n: w_1(y) \leq w_1(x) + \langle v, y - x \rangle, y \in \mathbb{R}^n\} .$$

Letting $Iw_1(x) = \{i \in I: w_1(x) = \langle a^i, x \rangle + b^i\}$, then $\partial w_1(x) = \text{Conv}\{a^i: i \in Iw_1(x)\}$, where Conv means convex hull. Let $I_w(x) = \{i \in I: w(x) = \langle a^i, x \rangle + b^i\}$. It can be seen that

$$\begin{aligned} \partial w(x) &= \partial w_1(x) && \text{if } x \notin P \\ \partial w(x) &= \text{Conv}(\partial w_1(x) \cup \{0\}) && \text{if } x \in \text{bd } P \\ \partial w(x) &= \{0\} && \text{if } x \in \text{int } P \end{aligned}$$

(where $\text{bd } P$ means boundary of P and $\text{int } P$ means the interior of P).

The same definitions and ideas are of course valid for f and f_1 .

Three different implementations of the relaxation method have been given by Agmon--the maximal distance, the maximal residual, and the cyclical methods.

Let

$$H^i = \{x: \langle \alpha^i, x \rangle + \beta^i \geq 0\} = \{x: \langle a^i, x \rangle + b^i \geq 0\}$$

and

$$E^i = \{x: \langle \alpha^i, x \rangle + \beta^i = 0\} = \{x: \langle a^i, x \rangle + b^i = 0\} .$$

The notation $d(x, S)$ will indicate the distance between a point x and a set S . Clearly,

$$d(x, E^i) = \frac{|\langle \alpha^i, x \rangle + \beta^i|}{\|\alpha^i\|} = |\langle a^i, x \rangle + b^i|$$

$$d(x, H^i) = \frac{\text{Max} (-(\langle \alpha^i, x \rangle + \beta^i), 0)}{\|\alpha^i\|} = \text{Max} (-(\langle a^i, x \rangle + b^i), 0) .$$

Note that $w(x) = - \text{Max}_{i \in I} d(x, H^i)$.

The relaxation method applied to (2.1) constructs a sequence of points in the following manner.

2.2.1 Choose $x^0 \in R^n$ arbitrary.

2.2.2 If $x^q \in P$, the sequence terminates.

If $x^q \notin P$, then determine i^q by one of the three methods below.

2.2.2.1 The maximal distance method:

let i^q be the index of a halfspace H^i which is the furthest distance away from x^q ; i.e.,

$$d(x^q, H^{i^q}) \geq d(x^q, H^i) \quad \forall i \in I.$$

2.2.2.2 The maximal residual method:

let i^q be the index of a most violated constraint;

$$\text{i.e., } \langle \alpha^{i^q}, x^q \rangle + \beta^{i^q} \leq \langle \alpha^i, x^q \rangle + \beta^i \quad \forall i \in I.$$

2.2.2.3 The cyclical method:

assume that $I = \{1, 2, 3, \dots, m\}$, then take

$$i^q = q + 1 \pmod{m}.$$

2.2.3 Then set

$$x^{q+1} = x^q + \sigma_q d(x^q, H^{i^q}) \frac{\alpha^{i^q}}{\|\alpha^{i^q}\|} ,$$

where usually $\sigma_q \in (0, 2]$, and go back to 2.2.2 with $q + q + 1$.

In the cyclical method it is quite possible that $d(x^q, H^{i^q}) = 0$, so that no move is taken at Step 3. We will not discuss it, as it does not seem easy to compare it to subgradient optimization.

If $\sigma_q = 1$ then x^{q+1} is the orthogonal projection of x^q on H^{i^q} ; if $\sigma_q = 2$, then x^{q+1} is the reflexion of x^q on E^{i^q} ; if $\sigma_q \in (1, 2)$ one talks of overprojection, and if $\sigma_q \in (0, 1)$ of underprojection.

Now let us describe the subgradient optimization algorithm, say for the function f :

2.3.1 Choose $x^0 \in R^n$.

2.3.2 Compute a subgradient of f at x^q : $u^q \in \partial f(x^q)$.
If $u^q = 0$, an optimal point has been found.

2.3.3 The next point x^{q+1} of the sequence will be obtained by moving from x^q in the direction of u^q by a certain step size. Go back to 2.3.2 with $q+1$ replacing q . Various proposals have been made for the step size:

2.3.3.1 Shor [10]: $x^{q+1} = x^q + \lambda_q \frac{u^q}{\|u^q\|}$ where $\lambda_q > 0$.

2.3.3.2 Held and Karp [5]: $x^{q+1} = x^q + \rho_q u^q$ where $\rho_q > 0$.

2.3.3.3 Held, Wolfe and Crowder [6], Eremin [2], Poljak [9], Oettli [8]:

$$x^{q+1} = x^q + \sigma_q \frac{\hat{f} - f(x^q)}{\|u^q\|^2} u^q ,$$

where $\sigma_q > 0$ and $\hat{f} > f(x^q)$ is a guess of the optimal value f^* of $f(x)$ which can either be:

2.3.3.3.1 an overestimate: $\hat{f} > f^*$,

2.3.3.3.2 an underestimate: $\hat{f} < f^*$,

2.3.3.3.3 the exact estimate: $\hat{f} = f^*$.

It can now be seen that:

- (i) The maximal distance relaxation method for solving (2.1) (which is the same as the maximal residual relaxation method for solving (2.2)) is equivalent to the subgradient algorithm applied to the function $w(x)$, using the step size given in 2.3.3.3.3 if P is not empty, and the step size of 2.3.3.3.1 with $\hat{f} = 0$ if P is empty.
- (ii) The maximal residual relaxation method for solving (2.1) is equivalent to the subgradient algorithm applied to the function $f(x)$, with the same use of step size as above.

For the maximal distance relaxation method (2.2.2.1) it is clear that the index i^q selected maximizes $d(x^q, H^{i^q})$, and thus satisfies

$$\begin{aligned} -\langle a^{i^q}, x^q \rangle - b^{i^q} &= \frac{-\langle a^{i^q}, x^q \rangle - \beta^{i^q}}{\|a^{i^q}\|} = d(x^q, H^{i^q}) \\ &= \text{Max}_{i \in I} d(x^q, H^i) = -w(x^q) \quad ; \end{aligned}$$

thus

$$i^q \in Iw(x^q) \quad \text{and} \quad a^{i^q} \in \partial w(x^q) \quad .$$

The only minor difference is that in subgradient optimization we are explicitly permitted the use of any $v^q \in \partial w(x^q)$. It would require the following extension of the relaxation method: if $Iw(x^q)$, the set of indices of the most distant halfspaces, contains more than one element, then any direction $v^q = \sum_{i \in Iw(x^q)} \dots$

$\eta_i \alpha^i$ can be used, where

$$\sum_{i \in Iw(x^q)} \eta_i = 1, \quad \eta_i \geq 0 \quad \forall i \in Iw(x^q),$$

and the next iterate would be:

$$x^{q+1} = x^q + \sigma_q \frac{-\langle v^q, x^q \rangle - c^q}{\|v^q\|} \frac{v^q}{\|v^q\|},$$

with $c^q = \sum_{i \in Iw(x^q)} \eta_i b^i$. Note that

$$\langle v^q, x^q \rangle + c^q = \sum_{i \in Iw(x^q)} \eta_i (\langle a^i, x^q \rangle + b^i) = w(x^q),$$

and that we are simply projecting in the direction of the half-space $\{x: \langle v^q, x \rangle + c^q \geq 0\}$. Also $\|v^q\| \leq 1$.

In the next section, on condition numbers and convergence rates, this distinction will require some attention.

One key point to notice is that if the system of inequalities (2.1) is feasible then $w^* = 0$, and thus the maximum value of $w(x)$ is known; that is exactly the case 2.3.3.3.3 of the subgradient optimization technique. That the step sizes are identical is trivial to check.

The identity between the maximal residual relaxation method and subgradient optimization as applied to f can be shown in the same fashion. In this case, also, if (2.1) is feasible then the maximum value of $f(x)$ is also known to be equal to zero.

Now if the system (2.1) has no solution, it is clear that $f^* < 0$ is not known, then the search for a Chebyshev solution to (2.1) becomes exactly equivalent to the problem of maximizing $f(x)$, which can be solved through subgradient optimization (in the case where the objective value is not known). One of the first examples of subgradient optimization can be found in some work by Eremin [2] on Chebyshev solutions of incompatible systems

of linear inequalities (as can the first result on convergence that I have seen). This also means that the results on rates of convergence of subgradient optimization gathered by Shor and the author are valid for the Chebyshev problem.

The main difference between the relaxation method for solving feasible systems of inequalities and the subgradient method for maximizing a concave function is that in the first case the optimal value of the objective is known, while in the second it is not. This has some implications for the theories of convergence for the two methods: in the first case a rule for the step size ($\sigma_q \in (0,2]$) can be given so that for each q , $d(x^{q+1}, P) \leq d(x^q, P)$; in the second case no step size rule can guarantee the monotonic decrease of anything. This also implies that a cyclical rule is not easy to use in subgradient optimization as it would not be clear if one should move or not in a direction a^i .

3. CONDITION NUMBERS AND RATES OF CONVERGENCE

Condition numbers have been defined and rates of convergence studied, for the relaxation method by Agmon [1] and the author [3], and for subgradient optimization by Shor [11] and the author [4]. As the two methods are quite similar, it is not surprising that the similarity will extend to the study of rates of convergence. In what follows, we will refer extensively to [3] and [4].

The study will be made for the comparison of the maximal distance method applied to (2.1) (or equivalently to (2.2)) and subgradient optimization applied to $w(x)$. The function $w(x)$ has the property that every extreme point of $\partial w(x)$ has norm one (except possibly for $x \in P$, where 0 is a subgradient).

We want to emphasize that not all of the results that follow can be extended to the comparison of f and (2.1).

3.1 The Relaxation Method

We will assume that P is not empty. For every $x^* \in \text{bd } P$, let

$$N_P(x^*) = \{v \in \mathbb{R}^n: \langle v, y - x^* \rangle \leq 0, \forall y \in P\} \quad ,$$

the normal cone to P at x^* . The normal cone can be written as:

$$N_P(x^*) = \{v \in \mathbb{R}^n: v = - \sum_{i \in I_W(x^*)} \eta_i a_i, \eta_i \geq 0\} \quad .$$

For any $x \in \mathbb{R}^n/P = \{x \in \mathbb{R}^n: x \notin P\}$, let $x^*(x)$ be the closest point to x in P . Then $x^*(x)$ can be characterized by the fact that

$$x - x^*(x) \in N_P(x^*(x)) \quad .$$

Define, for every $x^* \in \text{bd } P$ [3],

$$\mu^*(x^*) = \inf_{v \in N_P(x^*)} \max_{i \in I_W(x^*)} \{-\langle a^i, v \rangle\} \quad .$$

$$\|v\| = 1$$

It can be shown that $\mu^*(x^*) > 0$, and that

$$\mu^* = \inf_{x^* \in \text{bd } P} \mu^*(x^*) > 0 \quad .$$

The definition of the condition number μ^* is a sharpened version of a similar quantity defined by Agmon [1].

It can be shown [3] that (remember that $w^* = 0$):

$$\mu^* d(x, P) \leq -w(x) \leq d(x, P) \quad ,$$

where $d(x, P) = \|x - x^*(x)\|$, the distance between x and P . Both bounds can be shown to be tight, and thus μ^* could be defined by

$$\mu^* = \operatorname{Inf}_{x \in \mathbb{R}^n/P} \frac{-w(x)}{d(x,P)} .$$

Having defined μ^* , it is very easy to prove the geometric convergence of the relaxation method, as one can show that

$$\begin{aligned} d^2(x^{q+1}, P) &\leq d^2(x^q, P) - \sigma_q(2 - \sigma_q) w^2(x^q) \\ &\leq d^2(x^q, P) [1 - \sigma_q(2 - \sigma_q)\mu^{*2}] . \end{aligned}$$

It follows that, provided that $\sigma_q \in [0, 2]$, then $d(x^{q+1}, P) \leq d(x^q, P)$, so monotonic decrease is guaranteed; and also that if $\sigma_q \in [\varepsilon, 2 - \varepsilon]$, then

$$d(x^{q+1}, P) \leq d(x^q, P) \sqrt{1 - \varepsilon\mu^{*2}} ,$$

so that geometric convergence occurs.

Two other concepts, which are also trying to measure the well conditioning of a problem, were introduced in [3].

For $x^* \in \operatorname{bd} P$, let

$$C_P(x^*) = \{u \in \mathbb{R}^n: \langle a^i, u \rangle \geq 0, \forall i \in I_w(x^*)\} ,$$

the tangent cone to P at x^* . Clearly $N_P(x^*)$ is the polar cone to $C_P(x^*)$. A polytope P will be called obtuse if every tangent cone is obtuse:

$$-N_P(x^*) \subseteq C_P(x^*) \quad \forall x^* \in \operatorname{bd} P .$$

For a cone C , let $\nu(C)$ be the sine of the half aperture angle of the largest spherical cone contained in C , and let

$$\nu = \operatorname{Inf}_{x^* \in \operatorname{bd} P} \nu(C_P(x^*)) .$$

It is clear that $\nu > 0$ iff $\dim P = n$, and that $\nu \geq \frac{\sqrt{2}}{2}$ implies that P is obtuse. It can also be shown that $\mu^* \geq \nu$ [3]. It can then be shown that the maximal distance relaxation method applied to (2.1) or (2.2), assuming that $\dim P = n$, converges finitely if:

(i) P is obtuse and $\sigma_q \in [1, 2]$ for all q ;

or

(ii) $\sigma_q \in \left[\frac{2}{1 + 2\nu\sqrt{1 - \nu^2}} + \epsilon, 2 \right]$, where $\epsilon > 0$, for all q .

Thus, it has been proved that, if $\dim P = n$, there is a range of values of σ_q which leads to finite convergence, and to geometric convergence (until termination) [3].

3.2 Subgradient Optimization

For the concave function $w(x)$, a condition number can be defined:

$$\mu(x) = \inf_{u \in \partial w(x)} \frac{\langle u, x^*(x) - x \rangle}{\|u\| \cdot \|x^*(x) - x\|},$$

and

$$\mu = \inf_{x \in \mathbb{R}^n/P} \mu(x).$$

Rates of convergence for subgradient optimization depend upon this quantity.

Theorem 3.1 [4]

Let $\{x^q\}$ be a sequence which satisfies $x^{q+1} = x^q + \lambda_0 \rho^q \frac{u^q}{\|u^q\|}$,

where $u^q \in \partial w(x^q)$, $\lambda_0 > 0$, $\rho \in (0, 1)$.

Let

$$C = \text{Max} \left\{ \frac{1}{\rho}, \frac{\mu - \sqrt{\mu^2 - (1 - \rho^2)}}{1 - \rho^2} \right\}, \quad D = \frac{\mu + \sqrt{\mu^2 - (1 - \rho^2)}}{1 - \rho^2}$$

and

$$z(\mu) = \begin{cases} \sqrt{1 - \mu^2} & \text{if } \mu \leq \frac{\sqrt{2}}{2} \\ \frac{1}{2\mu} & \text{if } \mu > \frac{\sqrt{2}}{2} \end{cases} .$$

Then:

- (i) $\rho \geq z(\mu)$ and $d(x^0) \in [\lambda_O C, \lambda_O D]$ implies that for all q : $d(x^q) \leq d(x^0) \rho^q$;
- (ii) $\rho \geq z(\mu)$ and $d(x^0) < \lambda_O C$ implies that for all q : $d(x^q) \leq \lambda_O C \rho^q$;
- (iii) $\rho < z(\mu)$ or $d(x^0) > \lambda_O D$ may lead to convergence of $\{x^q\}$ to a non optimal point.

Theorem 3.2

The condition numbers μ and μ^* are equal.

Proof

This theorem can be proved in many different ways. It will be proved here in a way that gives new definitions of μ and μ^* , and also does not assume any other result.

By definition

$$\mu(x) = \text{Inf}_{u \in \partial w(x)} \frac{\langle u, x^*(x) - x \rangle}{\|u\| \cdot \|x^*(x) - x\|} .$$

But $\langle u, x^*(x) - x \rangle$ is positive for any $u \in \partial w(x)$, and

$$\frac{\langle u, x^*(x) - x \rangle}{\|u\| \cdot \|x^*(x) - x\|}$$

is a quasiconcave function of u on the domain $\{u: \langle u, x^*(x) - x \rangle \geq 0\}$; it thus follows that the infimum is attained at some of the extreme points of $\partial w(x)$; i.e.

$$\mu(x) = \text{Min}_{i \in Iw(x)} \frac{\langle a^i, x^*(x) - x \rangle}{\|x^*(x) - x\|} \quad \text{as } \|a^i\| = 1 .$$

If one uses the definition of concavity for w , one gets, as $w(x^*(x)) = 0$:

$$\langle a^i, x^*(x) - x \rangle \geq -w(x) \geq \langle a^j, x^*(x) - x \rangle$$

for any $i \in Iw(x)$, $j \in Iw(x^*(x))$.

Furthermore, if $w'(x; \delta)$ denotes the directional derivative at x in the direction δ , then

$$w'(x; x^*(x) - x) = \text{Inf}_{i \in Iw(x)} \langle a^i, x^*(x) - x \rangle$$

$$w'(x^*; x - x^*(x)) = \text{Inf}_{j \in Iw(x^*)} \langle a^j, x - x^*(x) \rangle ,$$

and thus:

$$w'(x; x^*(x) - x) \geq -w(x) \geq -w'(x^*; x - x^*(x)) .$$

Define for every x :

$$\delta(x) = \frac{x^*(x) - x}{\|x^*(x) - x\|} ;$$

then one has

$$w'(x; \delta(x)) \geq \frac{-w(x)}{d(x, P)} \geq -w'(x^*; -\delta(x)) .$$

It is also clear that

$$\mu(x) = w'(x; \delta(x)) \quad ,$$

so

$$\mu = \inf_{x \in \mathbb{R}^n/P} w'(x; \delta(x)) \quad .$$

Also,

$$-w'(x^*; -\delta(x)) = \sup_{j \in Iw(x^*)} \langle a^j, \delta(x) \rangle \geq \mu^*(x^*)$$

as

$$-\delta(x) \in N_P(x^*) \quad , \quad \text{and} \quad \|\delta(x)\| = 1 \quad .$$

And

$$\inf_{x \in x^* + N_P(x^*)} \{-w'(x^*; -\delta(x))\} = \mu^*(x^*) \quad ,$$

so that

$$\mu^* = \inf_{x^* \in \text{bd } P} \inf_{\substack{-\delta \in N_P(x^*) \\ \|\delta\| = 1}} (-w'(x^*; -\delta)) \quad .$$

From this, it is clear that $\mu \geq \mu^*$ (this could also be concluded from Theorem 3.1 in [4]). Now to show that $\mu = \mu^*$, the simplest way is to use v and x^* such that

$$\mu^* = \max_{i \in Iw(x^*)} \langle a^i, v \rangle$$

where

$$v \in -N_P(x^*) \quad \text{and} \quad \|v\| = 1 .$$

Now for any $\alpha \geq 0$, one has that

$$x^*(x^* - \alpha v) = x^* ,$$

i.e., the projection of $x^* - \alpha v$ on P is x^* . So

$$\mu^* = \text{Max}_{i \in Iw(x^*)} \langle a^i, \frac{x^* - x}{\alpha} \rangle$$

where $x = x^* - \alpha v$ and $\alpha = \|x^* - x\| = d(x, P)$. But,

$$\begin{aligned} \langle a^i, x^* - x \rangle &= \langle a^i, x^* \rangle + b^i - \langle a^i, x \rangle - b^i \\ &= -\langle a^i, x \rangle - b^i \end{aligned}$$

$$\text{as } \langle a^i, x^* \rangle + b^i = 0 \quad \forall i \in Iw(x^*) ;$$

so

$$\mu^* = \text{Max}_{i \in Iw(x^*)} \frac{-\langle a^i, x \rangle - b^i}{d(x, P)} .$$

It is easy to check that there exists $\epsilon > 0$ such that $Iw(x^* - \eta v) \subseteq Iw(x^*)$ for any $0 \leq \eta \leq \epsilon$. Thus, if $z = x^* - \epsilon v$,

$$\begin{aligned} \frac{-w(z)}{d(z, P)} &= \text{Max}_{i \in Iw(z)} \frac{-\langle a^i, z \rangle - b^i}{d(z, P)} \leq \mu^* \leq \text{Max}_{i \in I} \frac{-\langle a^i, z \rangle - b^i}{d(z, P)} \\ &= \frac{-w(z)}{d(z, P)} . \end{aligned}$$

This, by the way, proves that

$$\mu^* = \text{Inf}_{x \notin P} \frac{-w(x)}{d(x, P)} .$$

Also:

$$\begin{aligned} \mu^* &= -\frac{w(z)}{d(z,P)} = \operatorname{Max}_{i \in Iw(z)} \frac{-\langle a^i, z \rangle - b^i}{d(z,P)} \\ &= \frac{-\langle a^i, z \rangle - b^i}{d(z,P)} \quad \forall i \in Iw(z) \\ &= \operatorname{Min}_{i \in Iw(z)} \frac{\langle a^i, x^* - z \rangle}{d(z,P)} = \mu(z) \quad . \end{aligned}$$

From this it follows that

$$\mu^* = \mu(z) \geq \mu \quad ,$$

and as $\mu \geq \mu^*$, the theorem follows. ||

The proof given here is somewhat messy. One reason for its use is that it gives two new definitions of μ or μ^* in terms of directional derivatives. We will state a few related results, whose proofs are similar to, or included in, the one given above.

(i) If x is close enough to P , then

$$\langle a^i, x^*(x) - x \rangle = \langle a^j, x^*(x) - x \rangle \quad \forall i, j \in Iw(x) \quad .$$

This means that every extreme point of the subgradient set $\partial w(x)$ makes the same angle with the direction to the closest point to the optimal set.

(ii) Let

$$x^* \in \operatorname{bd} P \quad \text{and} \quad i \in Iw(x^*) \quad ;$$

then

$$\partial w(x^* - \alpha a^i) = \{a^i\} \quad \forall \alpha > 0 \quad .$$

Proof

$$\begin{aligned}
-\alpha &= a^i(x^* - \alpha a^i) + b^i \geq w(x^* - \alpha a^i) \\
&= \underset{j \in I}{\text{Min}} (\langle a^j, x^* \rangle + b^j - \alpha \langle a^j, a^i \rangle) \\
&\geq \underset{j \in I}{\text{Min}} (\langle a^j, x^* \rangle + b^j) + \underset{j \in I}{\text{Min}} (-\alpha \langle a^j, a^i \rangle) \\
&= 0 - \alpha \underset{j \in I}{\text{Max}} \langle a^j, a^i \rangle \\
&= -\alpha \quad \text{as} \quad \langle a^j, a^i \rangle \leq \langle a^i, a^i \rangle = 1 \quad \forall j \in I .
\end{aligned}$$

If $a^k \neq a^i$, then

$$a^k(x^* - \alpha a^i) + b^k \geq -\alpha \langle a^k, a^i \rangle > -\alpha ,$$

and thus $\partial w(x^* - \alpha a^i)$ contains only one element, i.e., $\{a^i\}$.

This last result depends crucially upon the assumption $\|a^i\| = 1 \quad i \in I$.

There is one result that is very important for subgradient optimization (for a general concave function), but which depends upon the assumption that all a^i have norm one. If we assume that an overestimate of μ is known, the only thing required to *guarantee* convergence is that an overestimate or $\mu d(x^0, P)$ be available [4].

Here it is clear that if $\tilde{w} \geq w^*$, then

$$\tilde{w} - w(x^0) \geq w^* - w(x^0) = -w(x^0) \geq \mu^* d(x^0, P) = \mu d(x^0, P) ,$$

and as any feasible solution to a dual problem to $\text{Max}_{x \in R^n} w(x)$ will provide a value $\tilde{w} \geq w^*$, it follows that $\tilde{w} - w(x^0)$ is an overestimate to $\mu d(x^0, P)$.

This is exactly the result that one would wish to extend to a general function $w(x)$. Unfortunately, given a proper definition of μ^* , this result does not extend without significant change.

Another minor result might be worth mentioning, about the quantity v ; it follows somewhat directly that if $\dim P = n$ and $x^* \in \text{bd } P$

$$v(C_P(x^*)) = \text{Max}_{\|e\|=1} \text{Min}_{i \in Iw(x^*)} \langle a^i, e \rangle ,$$

and thus $v(C_P(x^*)) = \text{Max}_{\|e\|=1} w_1'(x^*; e)$ which could be called the intensity of the steepest ascent at x^* of the function $w_1(x)$. It also follows that

$$v(C_P(x^*)) = \text{Min}_{u \in \partial w_1(x^*)} \|u\|$$

and thus

$$v = \text{Min}_{x^* \in \text{bd } P} \text{Min}_{u \in \partial w_1(x^*)} \|u\| .$$

4. CONCLUSION

The relationship between the maximal distance relaxation method and subgradient optimization as applied to a certain function $w(x)$, whose optimal value is known, has been made explicit.

Condition numbers defined in both theories are compared; in particular $\mu = \mu^*$, so that if $\mu \leq \frac{\sqrt{2}}{2}$, then the sustainable rate of convergence for subgradient optimization is exactly equal to the rate given by the relaxation method with $\lambda = 1$.

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AN EXTENSION OF THE METHOD OF SUBGRADIENTS*

Robin Chaney and Allen Goldstein

Until 1964, the only general method for the minimization of nondifferentiable convex functions was the so-called "cutting plane method" discovered by Cheney-Goldstein [1] and independently by Kelley [2].

In 1964, Shor [3] conceived the subgradient algorithm. Since that time, the method has been extensively developed in the Soviet Union. Some of these references are given in [3]-[12].

Lemarechal has reported at this conference that his computational experience showed that a modification of the subgradient method, due to Shor [8], was by far the most effective method for difficult nondifferentiable problems.

Generalizations of the subgradient method beyond convex functions have had partial success. The strongest result is due to Nurminskii [12]; this was corroborated by Poljak. An algorithm is formulated for a class of "weakly" convex functions. The rate of convergence obtained is slower than geometric.

The present extension of the subgradient method is to max families [13] and quasidifferentiable functions of B.N. Pshenichnyi. They are useful in many problems of applied optimization.

1. ALGORITHM

1.0 Hypotheses

Let E be a real Hilbert space and f a real-valued function defined on E . Let $x \in E$ and $\epsilon > 0$ be given. Let $B(x, \epsilon) = \{y: \|y - x\| \leq \epsilon\}$.

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Definition: f is called ε -non-stationary on $B(x, \varepsilon)$, if for any η , $0 < \eta \leq \varepsilon$ and for any minimizing sequence $\{x_k\}$ for f on $B(x, \eta)$, $\{\|x_k - x\|\} \rightarrow \eta$.

Let $m(x, \varepsilon) = \inf \{f(y) : y \in B(x, \varepsilon)\}$, $T(x, \varepsilon) = \{y \in E : f(y) \leq m(x, \varepsilon)\}$, and $D(x, \varepsilon) = \{y : \text{dist}(y, T(x, \varepsilon)) \leq \varepsilon\}$.

Assume there exist multifunctions $P : D(x, \varepsilon) \rightarrow$ subsets of $T(x, \varepsilon)$ and $\phi : D(x, \varepsilon) \rightarrow$ subsets of $B(0, 1)$. These functions, together with f , are assumed to satisfy [A], [B], and [C] below.

Assume there exist positive constants θ , μ , and L such that whenever $y \in D(x, \varepsilon)$ the following inequalities hold:

$$[A] \quad [\phi(y), y - P(y)] \geq \theta \|y - P(y)\|$$

$$[B] \quad \mu \|y - P(y)\| \leq f(y) - f(P(y)) \leq L \|y - P(y)\|$$

$$[C] \quad \|y - T(x, \varepsilon)\| \leq \|y - P(y)\| \leq$$

$$\|y - T(x, \varepsilon)\| \left(1 + \left(\frac{\mu}{\alpha}\right)^2\right)^{1/2}, \quad \alpha \geq L/\theta \quad .$$

1.1 Lemma

Assume [A], [B], and [C] and the number $m(x, \varepsilon)$ are given. Set $x_0 = x$, $\gamma_k = [f(x_k) - f(P(x_k))]/\alpha$ and $x_{k+1} = x_k - \gamma_k \phi(x_k)$. Then $\|x_k - P(x_k)\| \leq \beta^k \|x_0 - P(x_0)\|$ $0 < \beta = [1 - (\frac{\mu}{\alpha})^4]^{1/2} < 1$.

Proof: Assume $x_k \in D$. Compute $\|x_{k+1} - P(x_k)\|^2 = \|x_k - P(x_k) - \gamma_k \phi(x_k)\|^2 = \|x_k - P(x_k)\|^2 - 2\gamma_k [\phi(x_k), x_k - P(x_k)] + \gamma_k^2$. Now $[\phi(x_k), x_k - P(x_k)] \geq \theta \|x_k - P(x_k)\| \geq \theta [f(x_k) - f(P(x_k))]/L \geq [f(x_k) - f(P(x_k))]/\alpha_k = \gamma_k$. Hence $-2\gamma_k [\phi(x_k), x_k - P(x_k)] \leq -2\gamma_k^2$. Thus $\|x_{k+1} - P(x_{k+1})\|^2 \leq \|x_k - P(x_k)\|^2 - \gamma_k^2 = \|x_k - P(x_k)\|^2 - \left(\frac{f(x_k) - f(P(x_k))}{\alpha}\right)^2 \leq \|x_k - P(x_k)\|^2 \times$

$(1 - (\frac{\mu}{\alpha})^2) \leq \|x_k - T(x, \epsilon)\|^2 (1 - (\frac{\mu}{\alpha})^4)$. Since $\|x_{k+1} - T(x, \epsilon)\| \leq \|x_{k+1} - P(x_k)\|$ it follows that $x_{k+1} \in D(x, \epsilon)$.

Also, $\|x_{k+1} - P(x_{k+1})\| \leq \|x_{k+1} - T(x, \epsilon)\| (1 + (\frac{\mu}{\alpha})^2)^{1/2} \leq \|x_{k+1} - P(x_k)\| (1 + (\frac{\mu}{\alpha})^2)^{1/2}$. Thus $\|x_{k+1} - P(x_{k+1})\| \leq \|x_k - P(x_k)\| (1 - (\frac{\mu}{\alpha})^4)^{1/2}$.

1.2 Definition

By E_ϵ we denote the set $\{y \in E: f \text{ is } \epsilon\text{-non-stationary on } B(y, \epsilon)\}$.

1.3 Theorem

Assume f is bounded below on E , and assume [A], [B], and [C] hold for each $y \in E_\epsilon$. A finite sequence $\{x_k\}$ can be constructed such that for some integer M , $x_M \notin E_\epsilon$.

Proof: Observe first that by [C] $\inf \{m(x, \epsilon): x \in E_\epsilon\} \geq \mu\epsilon$. Take $x_0 \in E_\epsilon$, and let $\{x_j\}$ be the sequence constructed by the above theorem. Since $\|x_k - P(x_k)\| \rightarrow 0$, we have by [B] that $\{f(x_k)\} \rightarrow m(x_0, \epsilon)$. Again by [B], $f(x_0) - m(x_0, \epsilon) \geq \mu\|x_0 - P(x_0)\| = \mu\epsilon$. Hence $m(x_0, \epsilon) \leq f(x_0) - \mu\epsilon$. Therefore for some number k_1 , $f(x_{k_1}) \leq f(x_0) - \mu\epsilon/2$. Set $x'_0 = x_{k_1}$, and invoke the lemma with an ϵ ball around x'_0 and a new mapping P' assigning points to the set $T' = \{x \in E: f(x) \leq \inf\{f(y): y \in \beta(x'_0, \epsilon)\}\}$. Repeating this process a sequence of numbers $f(x_{k_1}), f(x_{k_2}), \dots$ will be generated. It follows that for some $k_s, x_{k_s} \notin E_\epsilon$; for otherwise $\{f(x_{k_i})\} \rightarrow -\infty$, contradicting that f is bounded below.

1.4 Remarks

In application, the number $m(x, \epsilon)$ is in general unknown. The estimate $f(x_0) - \mu\epsilon \geq m(x_0, \epsilon)$ may be employed in lieu of

$m(x_0, \varepsilon)$, however. We have that $f(x_0) - \mu\varepsilon = m(x_0, \eta)$ for some η , $0 < \eta \leq \varepsilon$, and fortunately the number η is not required in the lemma or theorem. If a lower bound for the infimum of f is known, the total number of iterations to reach x_{k_s} can be estimated from 1.1.

2. APPLICATION TO SOME MAX FAMILIES OF B.N. PSHENICHNYI

2.0 General Hypotheses

Let U be a metric compactum. Assume f is defined and continuous on $E_\varepsilon \times U$, that $f(\cdot, u)$ is Gateaux differentiable on E_ε , and that

$$|f'(x, u, h) - f'(y, u, h)| \leq K \|x - y\| \|h\| \quad (1)$$

for all x and y in E_ε , $u \in U$ and $h \in E$. It follows that f , and F below, are locally Lipschitz on E_ε . Let

$$F(x) = \max \{f(x, u) : u \in U\} \quad , \quad x \in E_\varepsilon \quad (2)$$

and

$$U(x) = \{u \in U : F(x) = f(x, u)\} \quad . \quad (3)$$

The Clarke differential of F [14] is defined as

$$F^\circ(x, h) = \limsup_{\lambda \downarrow 0 \quad k \rightarrow 0} \left\{ \frac{F(x + k + \lambda h) - F(x + k)}{\lambda} \right\} \quad .$$

The convex functional $h \rightarrow F^\circ(x, h)$ is the support function for a weakly compact convex set called the generalized gradient of F at x . It is denoted by $\partial F(x)$. See also [14], Propositions 1 and 2.

2.1 Lemma

$$\partial F(x) = \overline{co} \cup \partial f(x, u)$$

$$u \in U(x)$$

and

$$F^{\circ}(x, h) = F'_{+}(x, h) = \max \{f'(x, u, h) : u \in U(x)\} ,$$

where

$$F'_{+}(x, h) = \lim_{\lambda \rightarrow 0+} \left(\frac{F(x + \lambda h) - F(x)}{\lambda} \right)$$

Proof: By Proposition 4 (Clarke, [14]), $f'(x, u, h) = f^{\circ}(x, u, h)$.

Let $M(x) = \overline{co} \cup \partial f(x, u)$ and take $x \in S$ and $u \in U(x)$. Then $F(x + \lambda h) - F(x) \geq f(x + \lambda h, u) - f(x, u)$ and so

$$\liminf_{\lambda \rightarrow 0} \frac{F(x + \lambda h) - F(x)}{\lambda} \geq f'(x, u, h) = f^{\circ}(x, u, h)$$

$$= \max \{[\phi, h] : \phi \in \partial f(x, u)\} .$$

Since u is arbitrary in $U(x)$,

$$F^{\circ}(x, h) \geq \liminf_{\lambda \rightarrow 0} \frac{F(x + \lambda h) - F(x)}{\lambda}$$

$$\geq \max_{u \in U(x)} \max_{\phi \in \partial f(x, u)} [\phi, h] = \max_{\phi \in M(x)} [\phi, h] .$$

(See Pshenichnyi [13], p. 166.) Hence $M(x) \subseteq \partial F(x)$. By the definition of $F^{\circ}(x, h)$ there exist sequences $\{v_k\} \rightarrow 0$ and $\{\lambda_k\} \downarrow 0$ such that:

$$F^{\circ}(x, h) = \lim_{n \rightarrow \infty} \frac{F(x + v_n + \lambda_n h) - F(x + v_n)}{\lambda_n} .$$

Choose $u_n \in U(x + v_n + \lambda_n h)$ so that

$$F(x + v_n + \lambda_n h) = f(x + v_n + \lambda_n, u_n) \quad .$$

Then

$$\frac{F(x + v_n + \lambda_n h) - F(x + v_n)}{\lambda_n} \leq \frac{f(x + v_n + \lambda_n h, u_n) - f(x + v_n, u_n)}{\lambda_n} \quad .$$

We may assume $\{u_n\} \rightarrow u^*$.

Observe that $F(x) = \lim_{n \rightarrow \infty} f(x + v_n + \lambda_n h, u_n) = f(x, u^*)$; hence $u^* \in U(x)$. Therefore

$$\begin{aligned} \limsup_{\lambda \rightarrow 0} \frac{F(x + \lambda h) - F(x)}{\lambda} &\leq F^O(x, h) \\ &\leq \liminf_{n \rightarrow \infty} \frac{f(x + v_n + \lambda_n h, u_n) - f(x + v_n, u_n)}{\lambda_n} \\ &= \liminf \{f'(x + v_n, u_n, h) + f'(\xi_n, u_n, h) \\ &\quad - f'(x + v_n, u_n, h)\} \leq \liminf \{f'(x + v_n, u_n, h) \\ &\quad + K\lambda_n \|h\|^2\} \\ &= f'(x, u^*, h) = f^O(x, u^*, h) \leq \max \{f'(x, u, h) : u \in U(x)\} \\ &= \max_{u \in U(x)} \max_{\phi \in \partial f(x, u)} [\phi, h] = \max_{\phi \in M(x)} [\phi, h] \quad . \end{aligned}$$

Hence

$$\partial f(x) = M(x)$$

and

$$F'_+(x, h) = F^O(x, h) = \max_{\phi \in M(x)} [\phi, h] \quad .$$

2.2 Lemma

Given $\epsilon > 0$ and $\|h\| = 1$, we have that $(F(x + \epsilon h) - F(x))/\epsilon \geq F'(x, h) - K\epsilon$.

Proof:

$$\begin{aligned} & \frac{1}{\epsilon} \max \{f(x + \epsilon h, u) : u \in U\} - \frac{1}{\epsilon} \max \{f(x, u) : u \in U\} \\ & \geq \frac{1}{\epsilon} \max \{f(x, u) + f'(x, u, h) : u \in U(x)\} \\ & - K \epsilon \|h\|^2 - F(x)/\epsilon = \max \{f'(x, u, h) : u \in U(x)\} \\ & - K\epsilon \quad . \end{aligned}$$

To proceed, another hypothesis on our max family is convenient. We shall assume if the index values u are close that jumps in $f'(y, \dots, h)$ are bounded.

2.3 Lemma

Let x and y be arbitrary in E_ϵ and let $h = (y - x)/\|y - x\|$. Assume that for $\epsilon > 0$, $\|x - y\| \geq \epsilon$, $u_1 \in U(x)$ and $u_2 \in U(y)$ and $|F'_+(y, h)| \geq \sigma/2 > 0$; then

$$\frac{|f'(y, u_1, h) - f'(y, u_2, h)|}{\max \{|f'(y, u, h)| : u \in U(y)\}} \leq r < 1 \quad .$$

Then there exists $\epsilon_0 \leq \epsilon$ such that:

$$\frac{|F'_+(x, h) - F'_+(y, h)|}{|F'_+(y, h)|} \leq (1 + r)/2$$

whenever $\|x - y\| \leq \epsilon_0$.

Proof: Take $\|x - y\| \leq \varepsilon$. Then

$$\begin{aligned} |F'_+(x, h) - F'_+(y, h)| &= |\max \{f'(y, u, h) + f'(x, u, h) \\ &\quad - f'(y, u, h) : u \in U(x)\} - \max \{f'(y, u, h) : u \in U(y)\}| \\ &\leq \max \{|f'(x, u, h) - f'(y, u, h)| : u \in U(x)\} \\ &\quad + |\max \{f'(y, u, h) : u \in U(x)\} - \max \{f'(y, u, h) : \\ &\quad u \in U(y)\}| \leq k\varepsilon_0 + r|F'_+(y, h)|. \end{aligned}$$

Choose $\varepsilon_0 \leq \frac{(1-r)}{4K} \sigma$; then

$$\frac{|F'_+(x, h) - F'_+(y, h)|}{F'_+(y, h)} \leq \frac{1+r}{2}$$

whenever $\|x - y\| \leq \varepsilon_0$.

2.4 Lemma

Assume the hypotheses of 2.3 with $x = P(y)$. Then $F(y) - F(P(y)) \geq \sigma/2 \|y - P(y)\| (1 - q)$. By [18], if $|F'_+(x, h) - F'_+(y, h)| \leq q|F'_+(x, h)|$ then $(1 - q) |F'_+(y, y - P(y))| \leq |F(y) - F(P(y))|$. Since

$$q = \frac{1+r}{2} \quad \text{and} \quad F'_+\left(y, \frac{y - P(y)}{\|y - P(y)\|}\right) \geq \sigma/2,$$

the formula follows.

2.5 Hypotheses

We now collect some hypotheses.

(a) Given $x \in E_\varepsilon$, let $\sigma(x, \varepsilon) = \frac{f(x) - f(P(x))}{\varepsilon}$.

Since $f(P(x)) < f(x)$, $\sigma(x, \varepsilon) > 0$. We shall assume that: $\inf \{\sigma(x, \varepsilon) : x \in E_\varepsilon\} = \sigma > 0$.

(b) The hypothesis of 2.2.

(c) The hypothesis of 2.0.

2.6 Remark

Hypothesis (a) is proven in the finite dimensional case that follows in Section 3.

2.7 Theorem

As a consequence of the hypotheses of 2.5 the conditions A, B, and C are satisfied.

Proof:

$$\begin{aligned} \text{[A]} \quad & \text{By 2.1, if } \varepsilon < \sigma/2K, \text{ then } F' \left(y, \frac{P(y) - y}{\|P(y) - y\|} \right) \\ & \leq -\sigma(x, \varepsilon) + \sigma/2 \leq -\sigma + \sigma/2 = -\sigma/2. \end{aligned}$$

Take $\phi(y) \in \partial f(y, u)$, then $[\phi(y), P(y) - y]$

$$\leq F'(y, P(y) - y). \quad \text{Hence } [\phi(y), y - P(y)]$$

$$\geq \sigma/2 \|y - P(y)\|.$$

[B] By 2.0 F is Lipschitz continuous. For the lower bound, see 2.4.

[C] The existence of P follows by taking a minimizing sequence converging to $\text{dist}(y, T(x, \varepsilon))$.

3. APPLICATION TO THE QUASIDIFFERENTIABLE FUNCTION OF PSHENICHNYI

3.0* Hypotheses

The following two hypotheses will be used in the sequel. Let S be an arbitrary subset of E_n .

*In this section the one-sided differential $f'_+(x, h)$ will be simply written as $f'(x, h)$.

- [D] f is locally Lipschitz on S and there exists $\hat{\epsilon} > 0$ such that if $0 < \epsilon \leq \hat{\epsilon}$ and x is in S then f achieves its minima on $B(x, \epsilon)$ only on the boundary.
- [E] f is quasidifferentiable on S . Thus, f is locally Lipschitz on S , $f'(x, \cdot)$ exists for every x in S , and $f'(x, h) = \max \{[\phi, h] : \phi \in \partial f(x)\}$ for every x in S and h in E_n .

3.1 Remarks

The concept of quasidifferentiability is used here as it is in [13]; see also [15] and [16]. Clarke [17] has shown that many "max" functions are quasidifferentiable.

For a quasidifferentiable function f it is true that each $f'(x, \cdot)$ is continuous. Moreover, if 0 is not in $\partial f(x)$, it follows that $f'(x, \cdot)$ attains a unique minimum on the set of all unit vectors. To see this, note that

$$\begin{aligned}
 \min_{\|h\|=1} f'(x, h) &= \min_{\|h\| \leq 1} \max \{[\phi, h] : \phi \in \partial f(x)\} \\
 &= \max_{\|h\| \leq 1} \{\min [\phi, h] : \phi \in \partial f(x)\} \\
 &= \max \{[\phi, -\phi/\|\phi\|] : \phi \in \partial f(x)\} \\
 &= -\min \{\|\phi\| : \phi \in \partial f(x)\} ;
 \end{aligned}$$

hence $f'(x, \cdot)$ has its unique minimizer among unit vectors at $-\phi_0/\|\phi_0\|$, where ϕ_0 is the closest point to the origin in $\partial f(x)$.

3.2 Definitions

Assume 3.0[D] and let $\hat{\epsilon} > 0$ be as in 3.0[D]. Given $x \in S$ and ϵ in $[0, \hat{\epsilon}]$, let $h(x, \epsilon)$ be the set of all unit vectors h such

that $x + \epsilon h$ minimizes f on $B(x, \epsilon)$. Each vector in $h(x, \epsilon)$ is called a direction of ϵ -steepest descent, while

$$\frac{f(x + \epsilon h(x, \epsilon)) - f(x)}{\epsilon}$$

is called the rate of ϵ -steepest descent.

3.3 Lemma

Assume 3.0[E] and $x \in S$. Let $\{t_i\}$ be any positive sequence converging to 0, and let $h_i \in h(x, t_i)$ for each i . Let $\{t_k\}$ be a subsequence of $\{t_i\}$ such that $\{h_k\}$ converges to $h(0)$. Then

$$\begin{aligned} \lim_{k \rightarrow \infty} \frac{f(x + t_k h_k) - f(x)}{t_k} &= \min \{f'(x, h) : \|h\| = 1\} \\ &= f'(x, h(0)) \quad . \end{aligned}$$

Proof: By 3.1, $f'(x, \cdot)$ has a unique minimizer h_0 on the unit vectors. Then, $f(x + t_k h_k) - f(x) \leq f(x + t_k h_0) - f(x)$ for each k , and so

$$\limsup_{k \rightarrow \infty} \frac{f(x + t_k h_k) - f(x)}{t_k} \leq f'(x, h_0) \quad .$$

Also,

$$\begin{aligned} f'(x, h_0) &\leq f'(x, h(0)) = \lim_{k \rightarrow \infty} \frac{f(x + t_k h(0)) - f(x)}{t_k} \\ &= \liminf_{k \rightarrow \infty} \frac{f(x + t_k h_k) - f(x)}{t_k} + \lim_{k \rightarrow \infty} \frac{f(x + t_k h(0)) - f(x + t_k h_k)}{t_k} \\ &= \liminf_{k \rightarrow \infty} \frac{f(x + t_k h_k) - f(x)}{t_k} \quad , \end{aligned}$$

because of the Lipschitz condition on f . Therefore,

$$\lim_{k \rightarrow \infty} \frac{f(x + t_k h_k) - f(x)}{t_k} = f'(x, h_0) \quad ,$$

and moreover, $f'(x, h_0) = f'(x, h(0))$.

3.4 Remark

If 3.0[E] holds and if we put $h(x, 0) = \{h_0\}$ (for each x in S), then the multifunction $t \rightarrow h(x, t)$ is, for each x in S , upper semicontinuous at 0.

3.5 Theorem

Assume 3.0[D] and 3.0[E]. For each x in S , the multifunction $h(x, \cdot)$ is upper semicontinuous on $[0, \hat{\epsilon})$.

Proof: By 3.4, $h(x, \cdot)$ is upper semicontinuous at $t = 0$. Suppose $0 < t < \hat{\epsilon}$. Let $\{\theta_i\}$ be a sequence of real numbers converging to 0 and let $\{h_i\}$ be a sequence of vectors such that each h_i is in $h(x, t + \theta_i)$ and $\{h_i\}$ converges to \bar{h} . Let $q_i = (t + \theta_i) h_i$ for each i . Then $\{q_i\}$ converges to $t\bar{h}$ and it must be shown that \bar{h} is in $h(x, t)$. Two cases can be distinguished, accordingly as $\theta_i > 0$ for each i or $\theta_i < 0$ for each i .

Case (a). Suppose $\theta_i > 0$ for each i . Choose $x + p_i$ on the line segment which joins $x + th_i$ and $x + q_i$ such that $f(x + p_i) = f(x + th(x, t))$; this is possible because $f(x + th_i) \geq f(x + th(x, t)) > f(x + q_i)$. Hence, $|f(x + q_i) - f(x + th(x, t))| = |f(x + q_i) - f(x + p_i)| \leq L \|q_i - p_i\| \leq L\theta_i$, where L is a Lipschitz constant for f . Thus, $\lim_{i \rightarrow \infty} f(x + q_i) = f(x + t\bar{h}) = f(x + th(x, t))$ and so \bar{h} is in $h(x, t)$.

Case (b). Suppose $\theta_i < 0$ for each i . Choose \tilde{h} in $h(x, t)$ and choose $x + s_i$ on the line segment which joins $x + (t + \theta_i) \tilde{h}$

and $x + th^{\sim}$ such that $f(x + s_i) = f(x + q_i)$. This is possible, because $f(x + (t + \theta_i) h^{\sim}) \geq f(x + q_i) > f(x + th^{\sim})$. Much as before, one obtains $|f(x + q_i) - f(x + th^{\sim})| \leq L|\theta_i|$. Thus, $\lim_{i \rightarrow \infty} f(x + q_i) = f(x + th) = f(x + th^{\sim})$ and so \bar{h} is again in $h(x, t)$.

3.6 Miscellany

The following information is from [17] and [19].

Let f be a locally Lipschitz function on E_n . In this setting, the Clarke generalized gradient $\partial f(x)$ is the convex hull of the limit points of all sequences $\{\nabla f(x + h_i)\}$ for which $\{h_i\} \rightarrow 0$. The ε -generalized gradient $\partial_\varepsilon f(x)$ is the convex hull of all limit points of $\{\nabla f(x_i)\}$ where $\{x_i\}$ converges to some point in $B(x, \varepsilon)$. The sets $\partial f(x)$ and $\partial_\varepsilon f(x)$ are convex and compact and the multifunctions $x \rightarrow \partial f(x)$ and $x \rightarrow \partial_\varepsilon f(x)$ are upper semicontinuous (in the Hausdorff metric).

Let B be a closed ball and let Z be the set of all stationary points x (i.e., $0 \in \partial f(x)$) in B . Given $\delta > 0$, put $B_\delta = \{x \in B: \|x - Z\| \geq \delta\}$. Then, by [19], there exist numbers $\varepsilon_1 > 0$ and $\sigma > 0$ such that, for each x in B_δ , $0 \notin \partial_{\varepsilon_1} f(x)$ and $\|\nabla_{\varepsilon_1} f(x)\| \geq \sigma$, where $\nabla_{\varepsilon_1} f(x)$ is the point in $\partial_{\varepsilon_1} f(x)$ closest to the origin. If one puts $h = \nabla_{\varepsilon_1} f(x) / \|\nabla_{\varepsilon_1} f(x)\|$ then $[h, \phi] \geq \sigma$ for all ϕ in $\partial_{\varepsilon_1} f(x)$ and $x \in B_\delta$; moreover, if $0 < \lambda \leq \varepsilon_1$ then $f(x + \lambda h) - f(x) \leq -\sigma\lambda$ (for each x in B). It follows that $\min \{f'(x, h): \|h\| = 1\} \leq -\sigma$ for all x in B_δ .

Finally, it should also be noted that if B_δ is chosen as above and if $S = B_\delta$ then 3.0[D] is automatically satisfied with $\hat{\varepsilon} = \delta$; this is an immediate consequence of the fact that any unconstrained local minimizer of f is a stationary point.

3.7 Lemma

Assume 3.0[E] with $S = B_\delta$ and let $M > 1$ be given. Then there exists $\varepsilon^{\sim} > 0$ such that $f'(x, h) \leq -\sigma/2$ whenever $x \in B_\delta$, $\varepsilon^{\sim}/M \leq \varepsilon \leq \varepsilon^{\sim}$, and $h \in h(x, \varepsilon)$.

Proof: Let $\varepsilon_2 = \min(\delta, \varepsilon_1)$. For each x in B_δ define $\varepsilon(x) = \min\{\frac{1}{2}\varepsilon_2 \sup\{\varepsilon: f'(x, h) \leq -\sigma/2 \text{ for } \varepsilon/M \leq \gamma \leq \varepsilon \text{ and } h \in h(x, \gamma)\}\}$. Given x in B_δ , we know that $\varepsilon(x) > 0$, in view of 3.0[D], 3.5, and the inequality $f'(x, h(x, 0)) \leq -\sigma$.

To complete the proof, we must show that

$$\varepsilon \sim = \inf\{\varepsilon(x): x \in B_\delta\} > 0 \quad . \quad (1)$$

Take a sequence $\{x_k\}$ in B_δ so that $\{\varepsilon(x_k)\}$ decreases monotonically to $\varepsilon \sim$. Since B_δ is compact, we may require that $\{x_k\}$ converges to \bar{x} in B_δ . Now define $\varepsilon^* = \min\{\frac{1}{2}\varepsilon_2, \sup\{\varepsilon: f'(\bar{x}, h) \leq -3\sigma/4 \text{ for } 0 \leq \gamma \leq \varepsilon \text{ and } h \in h(\bar{x}, \gamma)\}\}$. As above, $\varepsilon^* > 0$.

Let $\varepsilon \geq 0$ be given with $\varepsilon \leq \varepsilon^*$. Let $T_\varepsilon = \{y: f(y) \leq f(\bar{x} + \varepsilon h(\bar{x}, \varepsilon))\}$ and let P_ε denote the multifunction which assigns to each y in B_δ the set of closest points in T_ε . There exists k_0 such that $\|x_k - \bar{x}\| < \varepsilon_2 - \varepsilon^*$ whenever $k \geq k_0$. For $k \geq k_0$,

$$\begin{aligned} \|x_k - P_\varepsilon(x_k)\| &\leq \|x_k - P_\varepsilon(\bar{x})\| \\ &\leq \|x_k - \bar{x}\| + \|\bar{x} - P_\varepsilon(\bar{x})\| < \varepsilon_2 \quad . \end{aligned} \quad (2)$$

Define $\gamma(x_k, \varepsilon) = \|x_k - P_\varepsilon(x_k)\|$. Since $0 \leq \gamma(x_k, \varepsilon) < \delta$, we know that f is minimized on $\beta(x_k, \gamma(x_k, \varepsilon))$ only at boundary points and we have the equations

$$P_\varepsilon(x_k) = x_k + \gamma(x_k, \varepsilon) h(x_k, \gamma(x_k, \varepsilon)) \quad , \quad k \geq k_0 \quad ,$$

and (3)

$$f(\bar{x} + \varepsilon h(\bar{x}, \varepsilon)) = f(x_k + \gamma(x_k, \varepsilon) h(x_k, \gamma(x_k, \varepsilon))) \quad .$$

It is clear that, for $k \geq k_0$, each function $\gamma(x_k, \cdot)$ is nondecreasing on $[0, \epsilon^*]$. At this point, we show that each such $\gamma(x_k, \cdot)$ is continuous on $[0, \epsilon^*]$.

Thus, for $k \geq k_0$, suppose $0 < \epsilon_0 \leq \epsilon^*$ and put $\alpha_k = \lim_{\epsilon \rightarrow \epsilon_0} \gamma(x_k, \epsilon)$. It is clear that

$$\alpha_k \leq \gamma(x_k, \epsilon_0) \quad . \quad (4)$$

Choose an increasing sequence $\{\epsilon^j\}$ which converges to ϵ_0 , and for each j , choose $h_j \in h(x_k, \gamma(x_k, \epsilon^j))$; we may assume that $\{h_j\}$ converges to a unit vector h^* . From 3.7 and (3), we have $f(x_k + \alpha_k h^*) = \lim_{j \rightarrow \infty} f(x_k + \gamma(x_k, \epsilon^j) h_j) = f(\bar{x} + \epsilon_0 h(\bar{x}, \epsilon_0))$ and so $x_k + \alpha_k h^*$ is in T_{ϵ_0} . Hence $\gamma(x_k, \epsilon_0) = \|x_k - P_{\epsilon_0}(x_k)\| \leq \|\alpha_k h^*\| = \alpha_k$. In view of (4), it follows that $\gamma(x_k, \cdot)$ is left continuous at ϵ_0 .

Let k remain fixed, suppose $0 \leq \epsilon_0 < \epsilon^*$, and put $b_k = \lim_{\epsilon \rightarrow \epsilon_0} \gamma(x_k, \epsilon)$. Again, it is clear that

$$b_k \geq \gamma(x_k, \epsilon_0) \quad . \quad (5)$$

Fix $\eta > 0$. Since $h(x_k, \cdot)$ is upper semicontinuous at $\gamma(x_k, \epsilon_0)$ there exists $\gamma > \gamma(x_k, \epsilon_0)$, h_0^* in $h(x_k, \gamma(x_k, \epsilon_0))$ and h_* in $h(x_k, \gamma)$ so $\gamma - \gamma(x_k, \epsilon_0) < \eta$ and $\|h_0^* - h_*\| < \eta$. Set $x_k^* = x_k + \gamma h_*$. Then $f(x_k^*) < f(x_k + \gamma(x_k, \epsilon_0) h(x_k, \gamma(x_k, \epsilon_0))) = f(\bar{x} + \epsilon_0 h(\bar{x}, \epsilon_0))$, by (3); hence, there exists $\epsilon > \epsilon_0$ so that $f(x_k^*) < f(\bar{x} + \epsilon h(\bar{x}, \epsilon))$. Therefore, $\gamma(x_k, \epsilon) = \|x_k - P_\epsilon(x_k)\| \leq \|x_k - x_k^*\| \leq \|x_k - (x_k + \gamma(x_k, \epsilon_0) h_0^*)\| + \|x_k + \gamma(x_k, \epsilon_0) h_0^* - (x_k + \gamma h_0^*)\| + \|x_k + \gamma h_0^* - (x_k + \gamma h_*)\| < \gamma(x_k, \epsilon_0) + \eta + \gamma\eta$. Since $\eta > 0$ was otherwise arbitrary, it follows that $b_k \leq \gamma(x_k, \epsilon_0)$ and, hence, from (5), that $\gamma(x_k, \cdot)$ is right continuous at ϵ_0 .

Now, for each $k \geq k_0$ and ϵ in $[\epsilon^*/2M, \epsilon^*]$, we choose $v_k(\epsilon)$ in the closed set $h(x_k, \gamma(x_k, \epsilon))$ so as to maximize $f'(x_k, \cdot)$ over

$h(x_k, \gamma(x_k, \epsilon))$. For each $k \geq k_0$ and ϵ in $[\epsilon^*/2M, \epsilon^*]$, we can, by 3.0[E], select $\phi_k(\epsilon)$ in $\partial f(x_k)$ so that

$$f'(x_k, v_k(\epsilon)) = [\phi_k(\epsilon), v_k(\epsilon)] \quad . \quad (6)$$

Assert now that an integer $k_1 \geq k_0$ exists so that

$$[\phi_k(\epsilon), v_k(\epsilon)] \leq -\frac{1}{2} \sigma \quad \text{for } k \geq k_1 \quad \text{and} \quad (7)$$

$$\epsilon^*/2M \leq \epsilon \leq \epsilon^* \quad .$$

If (7) were false, we could select an infinite subsequence $\{x_j\}$ of $\{x_k\}$ and a sequence $\{\epsilon_j\}$ in the interval $[\epsilon^*/2M, \epsilon^*]$ so that $[\phi_j(\epsilon_j), v_j(\epsilon_j)] > -\sigma/2$ for each j , $\{\epsilon_j\}$ converges to $\epsilon^\#$, $\{v_j(\epsilon_j)\}$ converges to $h^\#$, and $\{\phi_j(\epsilon_j)\}$ converges to $\phi^\#$. Since ∂f is upper semicontinuous at \bar{x} , $\phi^\#$ is in $\partial f(\bar{x})$. By 3.0[E] and the definition of ϵ^* , we get

$$-3\sigma/4 \geq f'(\bar{x}, h) \geq [\phi^\#, h] \quad , \quad \text{for every } h \text{ in } h(\bar{x}, \epsilon^\#) \quad . \quad (8)$$

Next, we shall show that $h^\#$ is in $h(\bar{x}, \epsilon^\#)$. If this were false, then there would be a unit vector h_1 so that

$$f(\bar{x} + \epsilon^\# h_1) < f(\bar{x} + \epsilon^\# h^\#) \quad . \quad (9)$$

For large j , the triangle inequality gives

$$| \|x_j - P_{\epsilon_j}(x_j)\| - \|\bar{x} - P_{\epsilon_j}(\bar{x})\| | \leq \|x_j - \bar{x}\| \quad ,$$

which amounts to $|\gamma(x_j, \epsilon_j) - \epsilon_j| \leq \|x_j - \bar{x}\|$. It follows that

$$\lim_{j \rightarrow \infty} \gamma(x_j, \epsilon_j) = \epsilon^\# \quad . \quad (10)$$

From (9), (10), and (3), we have for large j ,

$$\begin{aligned}
 f(\bar{x} + \varepsilon_j h_1) &< f(x_j + \gamma(x_j, \varepsilon_j) h(x_j, \gamma(x_j, \varepsilon_j))) \\
 &= f(\bar{x} + \varepsilon_j h(\bar{x}, \varepsilon_j)) \quad ;
 \end{aligned}
 \tag{11}$$

But (11) contradicts the definition of $h(\bar{x}, \varepsilon_j)$ and so it must be true that $h^\#$ is in $h(\bar{x}, \varepsilon^\#)$. Since $-\sigma/2 \leq \lim_{j \rightarrow \infty} [\phi_j(\varepsilon_j), v_j(\varepsilon_j)] = [\phi^\#, h^\#]$, the fact that $h^\#$ is in $h(\bar{x}, \varepsilon^\#)$ leads to a contradiction of (8). Hence (7) is established.

From (6), (7), and the continuity of each $\gamma(x_k, \cdot)$ we have, for $k \geq k_1$,

$$\begin{aligned}
 f'(x_k, h) &\leq -\frac{1}{2} \sigma \quad \text{for } \gamma(x_k, \varepsilon^*/2M) \leq \gamma \leq \gamma(x_k, \varepsilon^*) \quad \text{and} \\
 &h \in h(x_k, \gamma) \quad .
 \end{aligned}
 \tag{12}$$

Arguing as we did for (10), we find $\lim_{k \rightarrow \infty} \gamma(x_k, \varepsilon^*/2M) = \varepsilon^*/2M$ and $\lim_{k \rightarrow \infty} \gamma(x_k, \varepsilon^*) = \varepsilon^*$. Hence there is an integer $k_2 \geq k_1$ such that, whenever $k \geq k_2$, we have $\gamma(x_k, \varepsilon^*/2M) \leq 2\varepsilon^*/3M$ and $\gamma(x_k, \varepsilon^*) \geq 2\varepsilon^*/3$. Hence (12) leads to $f'(x_k, h) \leq -\sigma/2$ whenever $2\varepsilon^*/3M \leq \gamma \leq 2\varepsilon^*/3$ and $h \in h(x_k, \gamma)$, provided $k \geq k_2$.

Therefore, it is true that $\varepsilon(x_k) \geq 2\varepsilon^*/3$ for $k \geq k_2$. It follows that (1) is true and that the proof is complete.

3.8 Remarks

Suppose that 1.0[A] is replaced by the weaker hypothesis:

$$1.0[A'] \quad [\phi(y), y - P(y)] \geq \theta \quad \|y - P(y)\|$$

$$\text{for } y \text{ in } D, \text{ provided } \|y - P(y)\| \geq \varepsilon/N \quad ,$$

$$\text{where } N \geq 2L/\mu \quad .$$

If this replacement is made, then the argument which yields Lemma 1.1 can be used to obtain a version of Lemma 1.1 with the

following conclusion: "...there exists a smallest positive integer k^* such that $\|x_{k^*} - P(x_{k^*})\| < \epsilon/N$; moreover, for $k = 0, \dots, k^*-1$, one has $\|x_k - P(x_k)\| \leq \beta^k \|x_0 - P(x_0)\|$, where $\beta = [1 - \mu^4/\alpha^4]^{1/2}$."

It is then easy to show that Theorem 1.3 will still be valid, with essentially the same proof, if 1.0[A] is replaced by the weaker 1.0[A'].

3.9 Definitions

Assume 3.0[E] with $S = B_\delta$ and let $\epsilon > 0$. Given x in S , set $T(\epsilon, x) = \{y: f(y) \leq f(x + \epsilon h(x, \epsilon))\}$ and $D(\epsilon, x) = \{y \in B_\delta: \|y - T(\epsilon, x)\| \leq \epsilon\}$, and let $P = P_{\epsilon, x}$ be the multifunction which associates with each y in $D(\epsilon, x)$ the closest points to y in $T(\epsilon, x)$. Define $\phi_{\epsilon, x} = \phi$ by $\phi(y) = \{h/\|h\| : h \in \partial f(y)\}$.

3.10 Theorem

Let $\delta > 0$ and form the set B_δ as before. Assume 3.0[E] holds with $S = B_\delta$. There exists $\epsilon^{\sim} > 0$ and there exist constants θ, μ, L, α , and N such that for any x in B_δ and ϵ^* in $[\epsilon^{\sim}/2, \epsilon^{\sim}]$, with $T = T(\epsilon^*, x)$, $D = D(\epsilon^*, x)$, $\phi = \phi_{\epsilon^*, x}$, and $P = P_{\epsilon^*, x}$ as in 3.9, it is true that 1.0[A'], 1.0[B], and 1.0[C] hold.

Proof: Let L be the Lipschitz constant for f on B . Then put $\mu = \sigma$, $\theta = \sigma/(2L)$, $\alpha = L/\theta$, and $N = 2L/\sigma$. Let $M = 2N$ and choose ϵ^{\sim} as given by Lemma 3.8 to correspond to M .

Now suppose $\epsilon^{\sim}/2 \leq \epsilon^* \leq \epsilon^{\sim}$ and x is in B_δ . Since $\|y - T\| = \|y - P(y)\|$ for all y , it is obvious that 1.0[C] holds. If one sets $h = -\nabla_{\epsilon^*} f(y) / \|\nabla_{\epsilon^*} f(y)\|$, then by 3.6, $f(y) - f(P(y)) \geq f(y) - f(h + \|h - T\|h) \geq \sigma \|y - T\|$, for all y in D . Hence 1.0[B] is verified.

To see 1.0[A'], let y be in D with $\|y - T\| \geq \epsilon^*/N$. Put $\epsilon = \|y - T\| = \|y - P(y)\|$. Then $\epsilon \geq \epsilon^*/N = \epsilon^{\sim}/M$ while $\epsilon \leq \epsilon^* \leq \epsilon^{\sim}$. It follows from Lemma 3.7 that $f'(y, h(y, \epsilon)) \leq -\sigma/2$.

From 3.0[E] and from the equation $P(y) = y + \epsilon h(y, \epsilon)$, it follows that, for every ϕ in $\partial f(y)$, $[\phi, y - P(y)] = [\phi, -\epsilon h(y, \epsilon)] \geq \epsilon \sigma / 2 = \sigma / 2 \|y - P(y)\| \geq (\sigma / 2L) \|y - P(y)\| \|\phi\|$.

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NONSMOOTH OPTIMIZATION AND NONLINEAR PROGRAMMING

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We first give an algorithm of the penalization type for solving ordinary nonlinear programming problems in a "nondifferentiable optimization" context as in [1]. We investigate its rate of convergence. Then we apply it for solving some more specific problems such as finding a feasible point to a set of equality and inequality constraints; and, finally, we give some ideas which might lead to superlinearly convergent algorithms.

1. THE ALGORITHM

Let us consider an ordinary nonlinear programming problem

$$\left. \begin{array}{ll} \min f_0(x) & x \in R^n \\ f_i(x) \leq 0 & i \in I = \{1, 2, \dots, m\} \end{array} \right\} \quad (1)$$

in which we suppose the gradients $f_i'(x)$ to be Lipschitz continuous for $i = 0, 1, \dots, m$. For simplicity of notation, we suppose also that there is a degenerate constraint $f_i(x) \equiv 0 \forall x$. We do not require any Slater condition; in particular, the following development will accept equality constraints $\ell(x) = 0$, simulated by $f_i(x) = \ell(x)$, $f_{i+1}(x) = -\ell(x)$.

It would be possible to solve (1) by successive linearizations, i.e. by Newton's method, but the linearized problems

*This paper was written by C. Lemarechal from the tape recording of the talk.

usually have no solution. This motivates the following algorithm, which has an NSO background.

Set

$$F(x) = \max\{f_i(x)/i \mid i \in I\} \quad (\text{note that } F(x) \geq 0) \quad (2)$$

$$\phi_N(x) = f_0(x) + N F(x) \quad \begin{array}{l} \text{for some positive} \\ \text{number } N \text{ (}\phi_N \text{ is a} \\ \text{so-called exact} \\ \text{penalization)} \end{array} \quad (3)$$

$$I_\delta(x) = \{i \in I \mid f_i(x) \geq F(x) - \delta\} \quad \text{for some } \delta > 0 \quad (4)$$

We will give a descent algorithm for minimizing $\phi_N(x)$, starting from some x_0 .

Linearizing (1) about x , consider the following direction-finding problem in p :

$$\left. \begin{array}{ll} \min (f'_0(x), p) + \frac{1}{2}|p|^2 & p \in \mathbb{R}^n \\ (f'_i(x), p) + f_i(x) \leq 0 & i \in I_\delta(x) \end{array} \right\} \quad (5)$$

We suppose N and δ are chosen such that for any x in the level set $\{x \mid \phi_N(x) \leq \phi_N(x_0)\}$:

$$\text{Problem (5) is feasible} \quad (6)$$

and

its Lagrange multipliers u_i satisfy

$$\sum_{i \in I_\delta(x)} u_i \leq N \quad (7)$$

Therefore, during the iterations of the algorithm, we can check (5) and (6): if (5) is infeasible, we decrease δ . If (7) does not hold, we increase N .

Let $p(x)$ denote the solution of (5).

Theorem 1

$p(x) = 0$ if and only if x is feasible in (1) and it satisfies first order necessary optimality conditions.

It seems that, for computing $p(x)$, it is more convenient to solve the dual of (5), namely

$$\left. \begin{aligned} \min \quad & \frac{1}{2} \left| f'_0(x) + \sum_{i \in I_\delta(x)} u_i f'_i(x) \right|^2 - \sum_{i \in I_\delta(x)} u_i f_i(x) \\ & u \geq 0 \end{aligned} \right\} \quad (8)$$

which, in addition, allows us to check (7).

Now the algorithm for minimizing $\phi_N(x)$ is iterative: knowing x_k , solve (5) (or (8)) to obtain $p(x_k) = p_k$. Then perform a step in the direction p_k :

$$x_{k+1} = x_k + \alpha_k p_k \quad \alpha_k > 0 \quad (9)$$

The step size α_k gives a convergent algorithm provided it is chosen in the following way. Choose $\epsilon \in]0, 1[$. Try $\alpha = 1$. If

$$\phi_N(x_k + \alpha p_k) \leq \phi_N(x_k) - \epsilon \alpha |p_k|^2, \quad (10)$$

then take this value. Otherwise replace α by a fixed fraction of itself until (10) is satisfied.

In other words, p_k is a direction from x_k , in which it is possible to decrease $\phi_N(x)$ by a significant amount, given by (10).

Theorem 2

$F(x_k) \rightarrow 0$. Every cluster point x^* of the sequence $\{x_k\}$ is feasible in (1) and satisfies the first order necessary optimality conditions. Furthermore, if (1) is a linear program, convergence is finite.

Because of the presence of δ , (5) might contain only a small number of constraints compared to those appearing in (1).

2. RATES OF CONVERGENCE

If no constraints are present in (1), the algorithm reduces to the ordinary gradient method. The rate of convergence is therefore usually linear. However, it becomes quadratic when the solution of (1) is a vertex (intersection of n constraints), i.e. when (1) looks like a linear program.

Let us investigate this point more precisely, because it will show how to modify the algorithm to obtain superlinear convergence.

Let x^* be a solution of (1). We suppose the standard second order sufficient optimality conditions are satisfied, namely

- (i) $f'_i(x^*) \quad i \in I_0(x^*)$ are linearly independent;
- (ii) the (unique) Lagrange multipliers u^* are such that $u_i^* > 0, \quad i \in I_0(x^*);$
- (iii) $(p, L''_{xx}(x^*, u^*)p) > 0$ for any $p \neq 0$ such that $(p, f'_i(x^*)) = 0 \quad i \in I_0(x^*)$.

Concerning (i), a modification can be made to accept simulation of equality constraints (i.e., when $f_i(x) \leq 0$ and $-f_i(x) \leq 0$ are both present, one takes into account only one such constraint).

We can figure out the linear rate of convergence: let P be the projection mapping onto the subspace spanned by the active constraint gradients.

Lemma

The mapping $p(x)$ from R^n to R^n given by (5) is differentiable around x^* and its Jacobian is given by

$$p'(x^*) = - [P + (I - P)L''_{xx}(x^*, u^*)] \quad . \quad (11)$$

Now observe that $x_{k+1} = x_k + \alpha p(x_k)$ is just the standard process for solving $p(x) = 0$. It converges if the mapping $I + \alpha p'(x^*)$ has all its eigenvalues inside the unit circle. Then its rate of convergence is the largest modulus of all these eigenvalues.

It turns out that one can actually compute all these eigenvalues. Let $m = |I_0(x^*)|$ = number of active constraints at x^* . then $p'(x^*)$ has the following eigenvalues:

m of them are equal to -1 ;

the $n - m$ remaining ones are equal to $-\gamma_i$
 where γ_i are the $n - m$ strictly positive
 eigenvalues of the symmetric matrix
 $(I - P)L''_{xx}(I - P)$.

Therefore the eigenvalues of the process (9) (with fixed α) are

$1 - \alpha$ m times

and

$1 - \alpha\gamma_i$ $n - m$ times.

These numbers are obviously smaller than 1. For ensuring convergence, they must be also greater than -1 , which means that α must satisfy

$$\alpha < 2 \quad \text{and} \quad \alpha < \frac{2}{\gamma_i} \quad i = 1, \dots, n - m \quad .$$

The rate of convergence is then given by the smallest of these numbers. Hence, the Lagrange multipliers play a very important role for the rate of convergence. Roughly speaking, the larger m , the better the convergence. Ideally, when $m = n$, the choice $\alpha = 1$ makes all the eigenvalues zero, so that convergence is superlinear (Equation (9) reduces to Newton's method for solving a system of n equalities).

3. APPLICATIONS

Equations (5), (9) and (10) can be applied to solve a system of equalities and inequalities. In this case, $f_0(x) \equiv 0$ and (5) reduces to

$$\left. \begin{array}{l} \min |p|^2 \\ (f'_i(x), p) + f_i(x) \leq 0 \quad I \in I_\delta(x) \end{array} \right\} \quad (12)$$

which is Newton's method.

Here, the assumptions (6) and (7) have to be modified, and we suppose now that there exists $\delta > 0$ such that

$$\left. \begin{array}{l} \text{Problem (12) is feasible and} \\ \text{satisfies } |p(x)| \leq C|F(x)| \end{array} \right\} \quad (13)$$

This condition is satisfied for example if one supposes that the null-vector is never obtained by any *positive* combination of the active gradients: $\sum \lambda_i f'_i(x) \neq 0 \quad \forall \lambda_i \geq 0$.

Also the rule for choosing α is slightly different: we require

$$F(x_k + \alpha_k p_k) \leq (1 - \alpha_k \epsilon) F(x_k) \quad \text{for some } \epsilon \in]0, 1[\quad . \quad (14)$$

Theorem 3

- (i) $x_k \rightarrow x^*$ and $F(x^*) \leq 0$
- (ii) After a finite number of steps, $\alpha_k = 1$ satisfies (14)
- (iii) After a finite number of steps, $F(x_{k+1}) \leq C F^2(x_k)$
- (iv) $|x_k - x^*| \leq C q^{2^k}$ with $q < 1$
- (v) If the problem is linear, convergence is finite.

The algorithm can also be applied to minimize a max function of the form

$$f(x) = \max\{f_i(x) / i \in I\} \quad (I \text{ finite}) .$$

Then expressions (1) reduce to

$$\left. \begin{array}{l} \min v \\ f_i(x) - v \leq 0 \quad i \in I . \end{array} \right\} \quad (15)$$

It is then convenient not to consider (5) with $f_0 = v$, but to modify it slightly and solve

$$\left. \begin{array}{l} \min \eta + \frac{1}{2} |p|^2 \\ (f'_i(x), p) + f_i(x) - \eta \leq 0 \quad i \in I_\delta(x) . \end{array} \right\} \quad (16)$$

In this case, (6) and (7) are readily satisfied: (η, p) is feasible in (16) if η is large enough, and one can verify that the Lagrange multipliers sum up to 1. Hence, the choice of δ will depend only on computational convenience.

Theorem 2 applies. In particular, if the solution of (15) is a corner (i.e. $n + 1$ functions f_i are maximum at the solution), then convergence is quadratic. It would be interesting to generalize (16) to the case where the functions $f_i(x)$ are not explicitly known (in particular, when I is infinite).

4. IMPROVEMENT OF CONVERGENCE

Let us go back to (5). From (11), if $L''_{xx}(x^*, u^*) = I$, then $p'(x^*) = -I$; all its eigenvalues are -1 and the choice of $\alpha_k = 1$ yields superlinear convergence. Therefore, any information about the optimal Lagrange multipliers might allow us to make a linear transformation on the variables which would generate a very small maximum eigenvalue for $(I + \alpha p')$.

For example, when x is near x^* , replacing (5) by

$$\left. \begin{aligned} \min (f'_0(x), p) + \frac{1}{2}(p, L''_{xx}(x^*, u^*)p) \\ (f'_i(x), p) + f_i(x) \leq 0 \end{aligned} \right\} \quad (17)$$

would yield $p(x) = -(x - x^*) + o(|x - x^*|)$

and $u(x) = u^* + o(|x - x^*|)$.

To approximate (17), one could for example replace $L''(x^*, u^*)$ by $L''(x_k, u_{k-1})$ (u_{k-1} being the Lagrange multipliers of (17) with $x = x_{k-1}$).

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BUNDLE METHODS IN NONSMOOTH OPTIMIZATION

C. Lemarechal

This paper tries to synthesize what are called *conjugate sub-gradient methods*, and to extend them to a wider class of *bundle methods*. Also, we will show a connection with other methods primarily designed for solving ordinary mathematical programming problems. Our approach will be intuitive rather than algebraic: we will give not theorems, but ideas. Nothing essentially new will be said with respect to papers that have been published elsewhere.

We show that methods of conjugate gradients are perfectly justified as far as a local aspect is concerned, but that this local study is not enough for constructing efficient algorithms. We then try to replace the concept *local* by *finite neighborhood* and define the class of bundle methods. Finally we show that this class has a common background with many well-known methods.

Throughout the paper, $f(x)$ is a function defined on \mathbb{R}^n , convex, and Lipschitz, the latter hypothesis being the most important. \mathbb{R}^n is considered as a Hilbert space, i.e. we consider only the Euclidean norm. Also, we note that the dimension need not be finite; therefore we denote \mathbb{R}^n by H .

1. LOCAL ASPECT

In this section, we have a point x , fixed in H , and we ask the question: how can we find a direction of descent, i.e. an element $d \in H$ such that the directional derivative

$$f'(x, d) = \lim_{t \rightarrow 0} \frac{f(x + td) - f(x)}{t} \quad (1)$$

is strictly negative.

We follow Clarke's analysis [1]; as a Lipschitz function, f has a gradient almost everywhere in H . It is therefore possible to construct sequences $\{x_i\}$ such that $\nabla f(x_i)$ exists and $x_i \rightarrow x$. The corresponding sequences $\{\Delta f(x_i)\}$ are bounded and have (weak) cluster points. Define the set of all such cluster points:

$$M(x) = \{g \mid g = \lim \nabla f(x_i), x_i \rightarrow x, \nabla f(x_i) \text{ exists}\} . \quad (2)$$

Then there exists between (1) and (2) the following basic relation:

$$f'(x, d) = \sup \{(d, g) \mid g \in M(x)\} . \quad (3)$$

From this, several observations follow:

(a) The property that d is a descent direction implies that $-d$ makes an acute angle with every g in $M(x)$. The set of descent directions is the (open) polar cone of the convex cone generated by $M(x)$.

(b) $M(x)$ represents exactly the behavior of f in the neighborhood of x . Knowing a descent direction implies a *complete study* of f around x . It is not a trivial problem in general, unless f is differentiable at x ($M(x)$ is the singleton $\nabla f(x)$), or when f is made up of a *finite* number of *known differentiable* functions f_i which meet at x (then $M(x)$ is made up of the gradients $\nabla f_i(x)$).

(c) Any descent direction defines a hyperplane separating the convex sets $\{0\}$ and $\{\text{conv } M(x)\}$. Therefore, x is optimal iff there is no such hyperplane, i.e. $0 \in \text{conv } \{M(x)\}$. Accordingly, it is interesting to consider not $M(x)$ alone, but its convex hull $\partial f(x) = \text{conv } M(x)$ (see also (a) above) which is the ordinary subdifferential [11].

(d) The *best* descent direction is defined as solving the problem

$$\min_d f'(x, d) \iff \min_d \max \{(g, d) \mid g \in \partial f(x)\} . \quad (4)$$

Of course, the directional derivative being in essence a positively homogeneous function of d , it is necessary to normalize d . Then, when the Euclidean norm is chosen, it is possible to show, through some tedious calculations, that the optimal d is opposite to

$$\text{Nr } \partial f(x) = \text{Proj } 0/\partial f(x) \quad , \quad (5)$$

i.e. the point of minimal Euclidean norm in $\partial f(x)$. (We note here that, in the differentiable case, the optimal d in the ℓ_1 sense would be a vector of the canonical basis; in this framework, the *steepest descent* method would be Gauss-Seidel!)

(e) Of course, the probability for $M(x)$ not to be a singleton is zero. However, x might be so close to a point of nondifferentiability that it is impossible to find a numerically non-zero step size in the direction $-\nabla f(x)$. We are in fact interested in constructing directions that are numerically usable. From this viewpoint, there is no difference between a nondifferentiable function and a stiff function: they differ only on a set of measure 0. Accordingly, our development will be valid also for minimizing ordinary but ill-conditioned functions. ||

Now we give the basic ideas for constructing a descent direction or, equivalently, for constructing $M(x)$. Suppose we know k points in $M(x)$:

$$\{g_1, \dots, g_k\} \subset M(x)$$

for some integer k . This is initialized by computing g_1 directly. In view of (e) above, one generally has $g_1 = \nabla f(x)$. Knowledge of one point in $M(x)$ is really the minimal requirement for using a method based on gradients.

The question is: knowing this partial information about $M(x)$, is it possible to deduce a descent direction easily? If not, how can we determine some $g_{k+1} \in M(x)$ so as to improve the current approximation?

Since $f'(x, d) \geq \max \{(d, g_i) \mid i = 1, \dots, k\}$, we choose some d_k satisfying

$$(d_k, g_i) < 0 \quad i = 1, \dots, k \quad . \quad (6)$$

We hope that d_k is a descent direction; so, to check it, we tabulate the function $f(x + td_k)$ for $t \downarrow 0$. Then,

- either we find $t > 0$ such that $f(x + td_k) < f(x)$ and we are done,
- or $f(x + td_k) \geq f(x)$ for any t generated by this line search.

This is the only interesting case here. By convexity, for any $t > 0$ and for any $g \in M(x + td_k)$

$$f(x + td_k) \geq f(x) \geq f(x + td_k) + (g, x - x - td_k) \quad .$$

Passing to the limit and denoting by g_{k+1} any cluster point of g , $g_{k+1} \in M(x)$ by definition. Furthermore one has

$$(g_{k+1}, d_k) \geq 0 \quad . \quad (7)$$

Comparing (6) and (7), we see that, increasing k by 1 and computing a new d satisfying (6), we will certainly get a different direction.

Note that property (7) comes directly from convexity of f , which therefore seems essential. In fact it can be weakened to the so-called *weak upper semismoothness* [8] which, roughly speaking, implies: if the differential quotient $[f(x + td) - f(x)]/t$ goes to a positive limit, then the corresponding slope (g, d) (where $g \in M(x + td)$) goes also to a positive limit. ||

Now recall that we are trying to construct $M(x)$. In view of (6) and (7), in order for g_{k+1} to be *as good as possible*, d_k

should make the numbers (d_k, g_i) as negative as possible. This justifies a min-max strategy, which consists in computing d_k as solving

$$\min_d \max \{(d, g_i) \mid i = 1, \dots, k\} \quad . \quad (8)$$

Again it is necessary to bound d ; again, when using the Euclidean norm, (8) turns out to have the solution $d_k = -N\{g_1, \dots, g_k\}$. Note that this gives $d_1 = -g_1$ for $k = 1$. It is then just technical to prove that, if x is not optimal, the above process is finite, thanks to (7), (8), and the boundedness of $\{g_k\}$. When x is optimal $d_k \rightarrow 0$ (strongly), and, since $-d_k \in \partial f(x)$, this provides a stopping test. Note that when $M(x)$ is finite, the process is finite anyway.

To conclude this section we state that, knowing a convex set G included in $\partial f(x)$, the best we can do is to compute its vector of minimal length.

- If $G = \partial f(x)$ we then get the steepest descent direction.
- If G is a sufficient approximation of $\partial f(x)$, we get some descent direction.
- If G is too poor an approximation we can generate a new point in $M(x)$ and improve G by an infinite line-search. Repeating the process, it is then possible to find a descent direction, if any.

2. NUMERICAL ASPECT: ENLARGEMENT OF THE SUBDIFFERENTIAL

In the previous development, several questions remain open.

(a) Keep in mind Section 1(e). Strictly speaking, the process for constructing $M(x)$ is useless since it is probably a singleton.

(b) Suppose we try to minimize f by the following algorithm:

- x_k being given, first compute a descent direction d_k by the process of Section 1.

- Then move along d_k , for example with an optimal step size.

This, at best, simulates the steepest descent method, which is known to be slow, and may converge to a nonoptimal point when f is really nondifferentiable. In other words, this algorithm would converge very slowly to a nonoptimal point!

(c) For computing a descent direction--more specifically, for generating new points in $M(x)$ --we are supposed to perform infinite line-searches along each trial direction, with $t \rightarrow 0$. This is forbidden. ||

It appears that these phenomena come from the same imperfection: $M(x)$ is too small, containing only limits of gradients. Suppose, on the contrary, that we replace in (2) the concept " $x_i \rightarrow x$ " by " x_i close enough to x ". More precisely, for $\varepsilon > 0$ define some neighborhood $V_\varepsilon(x)$ (for example the ball of radius ε). Then enlarge $M(x)$ as follows

$$M(x) \subset M_\varepsilon(x) = \{g \mid g = \lim \nabla f(x_i), x_i \rightarrow y, y \in V_\varepsilon(x)\} \quad , \quad (9)$$

which directly follows Goldstein's analysis [3]. This new definition eliminates the above-mentioned phenomena:

(a') $M_\varepsilon(x)$ is never a singleton--unless f is linear, at least in $V_\varepsilon(x)$ (not an interesting case).

(b') If a direction d satisfies $(d, g) < 0 \forall g \in M_\varepsilon(x)$, it can be seen by integration that $f(x + td)$ is a decreasing function of t as long as $x + td \in V_\varepsilon(x)$. A line-search along d will get us out of $V_\varepsilon(x)$, and, from compactness, we will converge in a finite number of steps to some x_N such that $V_\varepsilon(x_N)$ contains the minimum of f .

(c') Constructing $M_\varepsilon(x)$ is easier than constructing $M(x)$ in the sense that we can stop the line-search along a trial direction d_k as soon as $x + td_k \in V_\varepsilon(x)$, i.e. for some finite $t > 0$. ||

Several enlargements of $M(x)$ are possible. One of them, coming from convex analysis, is particularly interesting, despite

the fact that it is difficult to define the corresponding $V_\varepsilon(x)$. We define directly the convex set

$$\partial_\varepsilon f(x) = \{g \mid \forall y \in H, f(y) \geq f(x) + (g, y - x) - \varepsilon\} . \quad (10)$$

(a") This set contains $\partial f(x)$. Also $\partial_\varepsilon f(x) = \partial f(x)$ only in very special situations. Moreover, it is a closed, convex, and (weakly) compact set (because f is Lipschitz).

(b") There exists the following basic relation [11, p.220]:

$$\inf_{t>0} \frac{f(x + td) - f(x) + \varepsilon}{t} = \sup\{(d, g) \mid g \in \partial_\varepsilon f(x)\} . \quad (11)$$

which directly extends (3). It follows that $0 \in \partial_\varepsilon f(x) \iff x$ minimizes f within ε . Also, if d is such that $(d, g) < 0 \forall g \in \partial_\varepsilon f(x)$, then (and only then) it is possible to find $t > 0$ such that $f(x + td) < f(x) - \varepsilon$.

(c") Let x and y be two different points in H . Let $g \in \partial f(y)$ and $\varepsilon \geq 0$. Then $g \in \partial_\varepsilon f(x)$ iff

$$f(y) \geq f(x) + (g, y - x) - \varepsilon . \quad (12)$$

This formula can be looked at from different angles:

(i) $x, y, g \in \partial f(y)$ are given. Then $g \in \partial_\varepsilon f(x)$ for any $\varepsilon \geq f(x) + (g, y - x) - f(y)$ (a known positive number).

(ii) $y, g \in \partial f(y), \varepsilon \geq 0$ are given. Then $g \in \partial_\varepsilon f(x)$ for any x such that $f(x) - (g, x) \leq f(y) - (g, y) + \varepsilon$, i.e. for any x close enough to y .

(iii) x and $\varepsilon \geq 0$ are given. For any y , any $g \in \partial f(y)$ is also in $\partial_\varepsilon f(x)$ provided $f(y) + (g, x - y) \geq f(x) - \varepsilon$. By continuity this is again true whenever y is close enough to x (g is bounded since f is Lipschitz). Observe that the left-hand side of this inequality is the value at x of the approximation of f linearized at y .

In particular, when y is the current point $x + td$ of some line-search, $g \in \partial f(x + td)$ belongs also to $\partial_\epsilon f(x)$ when $f(x + td) - t(g,d) \geq f(x) - \epsilon$, and this is eventually true provided $t \rightarrow 0$. ||

Thus, the introduction of $\partial_\epsilon f(x)$ does not explicitly define $V_\epsilon(x)$, but rather makes precise the concept *close enough to* in terms of objective value units. This is the really good feature of this set: it makes us able to compare movements Δx in H with the common measure of movements Δf in R . Of course, (12) is very useful since we can only generate points g of the form $\nabla f(y)$, and they must be *transported* into sets $\partial_\epsilon f(x)$.

Note that $-\text{Nr } \partial f(x)$ has a good geometrical interpretation: it is the steepest descent direction. Here, $-\text{Nr } \partial_\epsilon f(x)$ has no such simple interpretation. We can only say that there exists $\eta(\epsilon) > 0$ such that $-\text{Nr } \partial_\epsilon f(x)$ points towards the projection of x onto the level line $f(x) - \epsilon - \eta(\epsilon)$ (if such a level line exists, i.e. if $f(x) > \min f + \epsilon$). ||

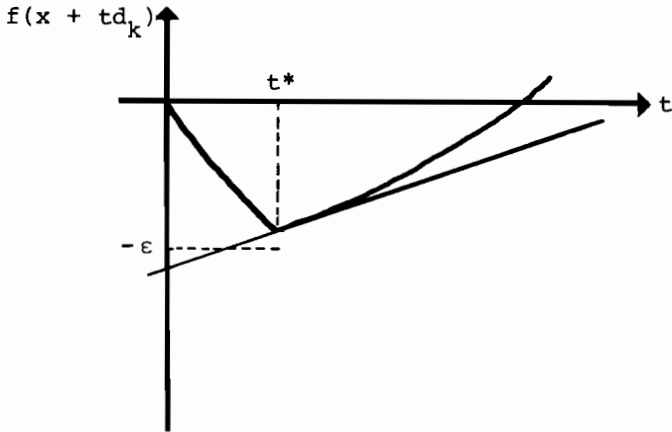
As the first application of this enlargement, we can adapt the algorithm of Section 1 to construct $\partial_\epsilon f(x)$, or equivalently to determine an ϵ -descent direction, i.e. a direction d such that $\inf \{f(x + td) \mid t > 0\} < f(x) - \epsilon$. Let x be fixed, and choose $\epsilon > 0$. Suppose k points g_1, \dots, g_k are already known in $\partial_\epsilon f(x)$. Determine d_k such that $(d_k, g_i) < 0$ $i = 1, \dots, k$. Again it is interesting to choose $d_k = -\text{Nr } \{g_1, \dots, g_k\}$. Make a line-search along d_k . If we can decrease by ϵ , we are done. Therefore, suppose $f(x + td_k) \geq f(x) - \epsilon$ for any t generated by the line-search. Two cases may occur:

(j) $f(x + td_k) \geq f(x) \forall t$. Then we are exactly in the situation described in Section 1. $(g, d_k) \geq 0 \forall t > 0$, $\forall g \in \partial f(x + td_k)$, and we can stop the line-search as soon as $t \leq [f(x + td_k) - f(x) + \epsilon] / (g, d_k)$ which eventually occurs (f is Lipschitz: (g, d_k) cannot go to $+\infty$).

(jj) Slightly more complicated is the case where some t_L is produced during the course of the line-search, such that:

$$f(x) > f(x + t_L d_k) \geq f(x) - \epsilon .$$

Then it might be impossible to produce a positive t such that any $g \in \partial f(x + td_k)$ simultaneously satisfies $g \in \partial_\epsilon f(x)$ and $(g, d_k) \geq 0$.



This sketch is a counter-example: there is a minimizing $t^* > 0$ along d_k . The set L of step sizes generating ϵ -subgradients at x is $[0, t^*[$. The set R of step sizes generating gradients satisfying (7) is $]t^*, +\infty[$ and $L \cap R = \emptyset$. In fact g_{k+1} should be a particular subgradient at the optimal t^* .

In that situation, the gimmick consists in computing g_{k+1} as a convex combination of $g_L \in \partial f(x + t_L d_k)$, $t_L \in L$, and $g_R \in \partial f(x + t_R d_k)$, $t_R \in R$. We choose this combination such that $(g_{k+1}, d_k) \approx 0$; and $g_{k+1} \in \partial_\epsilon f(x)$ if $t_R - t_L$ is small enough. ||

This is the ϵ -descent method given in [4]. It has a curious variant, in which we just neglect the test $f(x + td_k) < f(x) - \epsilon$, and we never move from x . Let us give its schematic description: fix $x \in H$. Let $g_1 \in \partial f(x)$. Set $k = 1$.

- Compute $d_k = -Nr \{g_1, \dots, g_k\}$.
- Minimize $f(x + td_k)$ for $t \geq 0$ and get an optimal $t_k \geq 0$.
- Compute $g_{k+1} \in \partial f(x + t_k d_k)$ such that $(g_{k+1}, d_k) = 0$.

Set $k = k + 1$ and go to 1. ||

Of course, the second and third steps are only schematic. One must simulate them by the development (j), (jj) above. The proof of convergence is fairly illustrative of what is done generally in this kind of algorithm:

Set $\alpha_k = f(x) - f(x + t_k d_k)$, $\epsilon_k = \max \{\alpha_i | i = 1, \dots, k\}$ and denote by x_k the point $x + t_i d_i$ at which this max is attained.

From (12) it can be seen that $g_{i+1} \in \partial_{\alpha_i} f(x)$. Also $\alpha_i \leq \epsilon_k$ $i = 1, \dots, k$. It follows that $g_i \in \partial_{\epsilon_k} f(x)$ $i = 1, \dots, k+1$; therefore

$$-d_{k+1} \in \partial_{\epsilon_k} f(x) \quad .$$

Now pass to the limit: as an increasing bounded sequence $\epsilon_k \rightarrow \epsilon^* = f(x) - \lim f(x_k)$. As in Section 1, $d_{k+1} \rightarrow 0$. There is a general result of uppersemicontinuity of the subdifferential, which says that in this situation, $0 \in \partial_{\epsilon^*} f(x)$. Hence, from (11): $f(x) \leq \min f + \epsilon^*$, which means

$$\lim f(x_k) \leq \min f \quad . \quad ||$$

3. CONSTRUCTING DESCENT ALGORITHMS: BUNDLE METHODS

So far, we have only considered f near a fixed x , strictly locally in Section 1, in a fixed neighborhood in Section 2. Now suppose we have performed several steps of some descent method. A certain number of points have been generated, at which the value of f has been computed together with some subgradient. We symbolize this information by the *bundle* $x_1, \dots, x_k; f_1, \dots, f_k; g_1, \dots, g_k$, where $f_i = f(x_i)$ and $g_i \in \partial f(x_i)$. We denote $G_k = \{g_1, \dots, g_k\}$.

This labeling is therefore not that of iterations. In fact we can store what we want, from just $\{x_k, f_k, g_k\}$ to all the information generated during all the line-searches from the beginning of the algorithm. This is purely a matter of core requirement.

Also it is just for simplifying notations that we suppose $f_i = f(x_i)$ and $g_i \in \partial f(x_i)$. We might suppose that the user is unable to compute exactly the function and its gradient. In fact, all we have said and all we are going to say remains valid when we suppose that the user is given x_i together with some prescribed tolerance ϵ_i , and that he returns f_i and g_i such that

$$f(x_i) - \epsilon_i \leq f_i \leq f(x_i) \quad \text{and} \quad g_i \in \partial_{\epsilon_i} f(x_i) \quad .$$

Then, the only difference with the case $\epsilon_i = 0$ is a minor modification of (12). ||

Now the question is: how do we use this bundle to compute "at best" a new point x_{k+1} , or to generate a new subgradient g_{k+1} ?

From Section 1, the simplest answer is the following. Let us suppose that $G_k \subset \partial f(x_k)$; then we should choose $d_k = -\text{Nr } G_k$. This technique is fully justified in two cases.

- (I) When the algorithm has been stuck from x_1 to x_k : all the points x_i are close together and all the g_i 's are approximately in $\partial f(x_k)$. Of course, this is an emergency situation: constructing algorithms that we hope do not get stuck too often is precisely what we want to do.
- (II) When the situation happens to be
- f is quadratic and
 - all the line-searches have been exact from x_1 to x_k . In that case, $\text{Nr } G_k$ turns out to be the direction of conjugate gradients which is optimal in the sense that it points towards the minimum of f in the set generated by the bundle. Of course, this situation is not very interesting either: the aim of NSO is not to minimize quadratic functions with exact line-searches!

In summary, $d_k = - \text{Nr } G_k$ is a poorly justified choice, which is likely to become worse when k increases. The first idea for overcoming this has been to reset periodically, to force k to be small. This has given variants of what have come to be called *conjugate subgradient* methods.

(a) In [14], one resets as late as possible, i.e. when $\text{Nr } G_k \simeq 0$ (in which case it would be foolish to interpret it as a direction).

(b) In [5] a rather artificial observation was made: there is a simple test that detects situation (II) above, namely

$$(g_i, x_i - x_1) \leq 0 \quad i = 1, \dots, k, \quad (13)$$

which means that $g_i \in \partial_{\varepsilon_i} f(x_1)$, where $\varepsilon_i = f(x_1) - f(x_i)$. (In fact, equality holds in (13) when the situation is as in (II)). Accordingly, one resets whenever (13) is not met.

(c) These two rules for resetting are based on (II) above and do not seem to be very promising. More interesting are rules given in [9], directly based on (I) and on the fact that resetting leads to restart on the gradient, which is clumsy. It is better not to *reset* but to *delete* those g_i 's that appear to be poor approximations of $\partial f(x_k)$. This can be done in either of the two following ways:

- delete those g_i 's for which $|x_i - x_k|$ is too large,
or
- define the numbers

$$\alpha_i = f(x_k) - f(x_i) - (g_i, x_k - x_i) \quad i = 1, \dots, k, \quad (14)$$

and delete those g_i 's for which α_i is too large. (Note that $\alpha_i \geq 0$, $\alpha_k = 0$. All the information is on hand to compute α_i and there is a simple recurrence formula for computing them cheaply.) ||

The latter is very interesting because, as in Section 2, it again gives a means to measure the proximity of x_i to x_k in terms of differences in the objective. From (12) it can be seen that $g_i \in \partial_{\alpha_i} f(x_k)$. Thus, α_i measures how far g_i is from $\partial f(x_k)$. Intuitively, when computing d_k in terms of the g_i 's, the weight of g_i should be smaller when α_i is larger.

Deletion rule (14) can be considered as a weighting of g_i by 1 if $\alpha_i \leq \delta$, 0 if $\alpha_i > \delta$, where δ is some chosen threshold. This is a tough weighting. Clearly, it would be interesting to smooth it. We think that efficient descent methods would be obtained if this question of weighting were properly solved. The next development tries to answer this question, probably imperfectly, as it does not eliminate the need for a threshold.

From the fact that $g_i \in \partial_{\alpha_i} f(x_k)$ $i = 1, \dots, k$ we can write

$$\forall y \in H, \quad f(y) \geq f(x_k) + (g_i, y - x_k) - \alpha_i \quad i = 1, \dots, k,$$

or equivalently, by convex combination:

$$f(y) \geq f(x_k) + (\sum \lambda_i g_i, y - x_k) - \sum \lambda_i \alpha_i$$

$$\forall \lambda: \quad \lambda_i \geq 0, \quad \sum \lambda_i = 1 \quad . \quad (15)$$

Every $y \in H$ can be described as $y = x_k + td$, $t > 0$, $|d| = 1$. In particular, if we choose $y = x_k + td$ as the next iterate, the decrement in the objective cannot be larger than the (we hope positive) number

$$\min \{-t(\sum \lambda_i g_i, d) + \sum \lambda_i \alpha_i \mid \lambda_i \geq 0, \sum \lambda_i = 1\} \quad .$$

It is therefore natural to look for the pair t and d that yields the best possible decrement, i.e. that solves

$$\max_{t>0, |d|=1} \min_{\lambda} -t \left(\sum \lambda_i g_i, d \right) + \sum \lambda_i \alpha_i \quad . \quad (16)$$

We do not see clearly how to solve (16) (which, by the way, has no solution in general: unless $0 \in \text{conv } G_k$, t is infinite). However, parallelling with the general scheme of a descent method, we might fix t at some *guessed* value, solve (16) for d alone, and then forget (16) and perform a line-search along the *optimal* d . In that case, (16) becomes a saddle point problem in λ, d , very similar to (8); its solution is

$$\bar{d} = - \sum \bar{\lambda}_i g_i / \left| \sum \bar{\lambda}_i g_i \right|$$

where $\bar{\lambda}$ solves

$$\min \left| \sum \lambda_i g_i \right| + \frac{1}{t} \sum \lambda_i \alpha_i \quad \lambda_i \geq 0, \quad \sum \lambda_i = 1 \quad . \quad (17)$$

(The case $\sum \bar{\lambda}_i g_i = 0$ needs finer analysis, which is not of interest here.)

Although (17) is a well-posed problem, it is not workable; but there is a trick for transforming it. $1/t$ can interpreted as the (positive) Lagrange multiplier associated with some (inequality) constraint of the form $\sum \lambda_i \alpha_i \leq \epsilon$. To $t > 0$ given, there corresponds the right-hand side $\epsilon = \sum \bar{\lambda}_i \alpha_i$. Since $\alpha_i \geq 0$ and $\alpha_k = 0$, the range for ϵ is $\epsilon \geq 0$. Finally, because the mapping $z \rightarrow \frac{1}{2}z^2$ is monotone increasing, (17) is equivalent to

$$\left\{ \begin{array}{l} \min \frac{1}{2} \left| \sum \lambda_i g_i \right|^2 \\ \sum \lambda_i = 1 \quad \lambda_i \geq 0 \\ \sum \lambda_i \alpha_i \leq \epsilon \quad , \end{array} \right. \quad (18)$$

which is a simple constrained least squares problem. The paper by Mifflin in this volume is devoted to solving it.

In (17), $t = |x_{k+1} - x_k|$ is unknown. We now give an interpretation of the unknown parameter ϵ in (18). It is convenient to denote

$$G_k(\epsilon) = \{g \mid g = \sum \lambda_i g_i, \sum \lambda_i = 1, \lambda_i \geq 0, \sum \lambda_i \alpha_i \leq \epsilon\} .$$

It is a convex polyhedron included in G_k . Using (12), it is not difficult to show that $G_k(\epsilon) \subset \partial_\epsilon f(x_k)$. Denote also $s(\epsilon) = \sum \lambda_i g_i$ the optimal solution in (18). Then, $s(\epsilon)$ defines the best hyperplane separating $G_k(\epsilon)$ from the origin. If ϵ is very small, we can hope that $G_k(\epsilon)$ is a good approximation of $\partial_\epsilon f(x_k)$, so $s(\epsilon)$ will also separate $\partial_\epsilon f(x_k)$ from the origin. This will guarantee a decrease by ϵ in the direction $-s(\epsilon)$, hence the need to choose ϵ as large as possible. On the other hand, if ϵ is too large, $G_k(\epsilon)$ is a poor approximation and $s(\epsilon)$ becomes meaningless in terms of $\partial_\epsilon f(x_k)$.

Thus, the philosophy underlying this development is a construction of $\partial_\epsilon f(x)$ for varying ϵ and x --instead of fixed as in Section 2.

We can call *bundle methods* the class of methods that proceed as follows: at each iteration, consider the bundle of information $g_1, \dots, g_k; \alpha_1, \dots, \alpha_k$. Choose $\epsilon > 0$. Solve (18) for $s(\epsilon)$. Make a step along $-s(\epsilon)$. For these methods to be efficient, several questions should be solved.

Is (18) really the proper problem to solve for computing the direction? We are supposed to find a hyperplane separating $G_k(\epsilon)$ from 0. However, such a hyperplane would be found also by changing the metric, i.e. by defining some positive definite matrix H , considering the objective $\frac{1}{2}(\sum \lambda_i g_i, H \sum \lambda_i g_i)$ in (18) instead of $\frac{1}{2}|\sum \lambda_i g_i|^2$, and taking $s(\epsilon) = H \sum \lambda_i g_i$ as the direction.

Second, we have not been able so far to find satisfying automatic rules for choosing ϵ at each iteration. (An efficient heuristic is $\epsilon = \frac{1}{2}[f(x_k) - \min f]$, but it is not implementable.)

Since ε should be related to $f(x_k) - f(x_{k+1})$, a dialogue between (18) and the line-search might be necessary; this would make the direction depend on the step size, and lead to *gauche* algorithms, with *curve-searches*.

Finally, once the direction is found, we should logically move if we thus get a *sufficient* decrease. Otherwise, as in Section 2, we should add a gradient into the bundle and compute a new direction from the same x_k . How do we make the decision to move, and which gradient do we add into the bundle?

In conjugate subgradient methods, the situation is clear: one has an estimate $|d_k|^2$ of $-f'(x_k, d_k)$. One chooses $0 < m_2 < m_1 < 1$ and $\varepsilon > 0$. We look for $y = x_k + td_k$ and $g_{k+1} \in \partial f(y)$ such that

$$(g_{k+1}, d_k) \geq -m_1 |d_k|^2 .$$

For moving, we require in addition

$$f(y) \leq f(x_k) - m_2 t |d_k|^2 \quad (\text{serious step: } x_{k+1} = y) .$$

If this is impossible, we require

$$f(y) - t(g_{k+1}, d_k) \geq f(x_k) - \varepsilon \quad (\text{null-step; } g_{k+1} \in \partial_\varepsilon f(x_k)) .$$

This ensures that the direction will change at the next iteration, and also that the decrease in the objective is sufficient. In bundle methods, there is no clear reason to choose the same criteria. ||

These questions are still open.

4. RELATIONS WITH OTHER METHODS

(a) When ε approaches 0 in (18), it is clear that $s(\varepsilon)$ goes to g_k . More precisely, a bundle method with $\varepsilon = 0$ would reduce to the algorithm of Section 1.

(b) No less obviously, if ϵ is large enough (say $\epsilon \geq \max \alpha_i$), $s(\epsilon)$ is the direction of conjugate subgradients.

(c) It is also clear that (15) is just a fancy way of writing:

$$f(y) \geq f(x_i) + (g_i, y - x_i) \quad i = 1, \dots, k \quad ,$$

which is the basis for cutting planes. A bundle method is therefore a Boxstep method [7] (or rather Ballstep, since we definitely prefer Euclidean "boxes"), i.e. the method that consists in minimizing the linearization \bar{f} of f inside a ball around x_k , and then extrapolating by some line-search. In such a method, the size of the ball is t in (17), and we are now going to make clear its relation with ϵ in (18).

First of all, $s(\epsilon) = 0$ corresponds to the case where the minimum of \bar{f} is attained in the interior of the ball. In that case, we have an optimality condition $f(x_k) \leq \min f + \epsilon$ (since $s(\epsilon) = 0 \in \partial_\epsilon f(x_k)$); on the other hand, ballstep is then equivalent to the pure cutting plane, and solving the linear program without a normalization constraint also provides an underestimate on $\min f$.

Suppose now the typical situation $s(\epsilon) \neq 0$. Let $u \geq 0$ be the Lagrange multiplier of the extra constraint $\sum \lambda_i \alpha_i \leq \epsilon$ in (18). Then

$$u = - \frac{d}{d\epsilon} \left(\frac{1}{2} |s(\epsilon)|^2 \right) = - |s(\epsilon)| \frac{d}{d\epsilon} (|s(\epsilon)|) \quad .$$

Now consider (17); set $\epsilon = \sum \bar{\lambda}_i \alpha_i$. There also

$$\frac{1}{t} = - \frac{d}{d\epsilon} (|\sum \bar{\lambda}_i g_i|) \quad .$$

From the equivalence between (17) and (18), we know that $\sum \bar{\lambda}_i g_i = s(\epsilon)$. Therefore $\frac{1}{t} = - \frac{d}{d\epsilon} (|s(\epsilon)|)$ and we have the basic relation

$$ut = |s(\epsilon)| \quad . \quad (19)$$

In ballstep, the minimum of the linearization \bar{f} of f is attained at $\bar{x} = x_k + t\bar{d}$ (provided the ball constraint is active) where \bar{d} , of norm 1, is given a posteriori by $\bar{d} = -s(\epsilon)/|s(\epsilon)|$. Therefore, after having solved (18), we can say that, if the ballstep problem had been solved with a ball size $t = |s(\epsilon)|/u$, one would have obtained the point

$$\bar{x} = x_k + t\bar{d} = x_k + \frac{|s(\epsilon)|}{u} \times (-s(\epsilon)/|s(\epsilon)|) = x_k - \frac{1}{u} s(\epsilon) .$$

In particular, when $\epsilon \rightarrow +\infty$, $u \rightarrow 0$. From (19), if $0 \notin \text{conv } G_k$, $t \rightarrow +\infty$. We can interpret this result by saying: when $0 \notin \text{conv } G_k$, the pure cutting plane problem is unbounded. There is one (infinite) solution--namely the limit of solutions of cutting plane problems bounded with the Euclidean norm--which gives the conjugate subgradient direction (since $s(\epsilon) = \text{Nr } G_k$ when $u = 0$).

On the other hand, when $\epsilon \rightarrow 0$, u goes to some finite value, i.e. t does not go to 0. This means that there is a strictly positive t_0 such that the ballstep direction is that of steepest descent whenever the ball size is smaller than t_0 . (This was observed in [7]).

(d) Consider the case where the function to be minimized has the form

$$f(x) = \max_i h_i(x) .$$

In this volume, Pshenichnyi gives an algorithm in which the direction d_k solves the problem

$$\left\{ \begin{array}{l} \min_{\eta, d} \quad \eta + \frac{1}{2}K|d|^2 \\ (\forall h_i(x_k), d) + h_i(x_k) \leq \eta \end{array} \right. \quad \begin{array}{l} \text{i such that} \\ h_i(x_k) \geq f(x_k) - \delta \end{array} \quad (20)$$

for some positive constants K and δ (strictly speaking, he chooses $K = 1$, but it is more convenient for our development to allow for different values).

Since, at the optimum of (20), there is at least one active constraint, we can replace η by $\max_i [(\nabla h_i(x_k), d) + h_i(x_k)]$ and solve instead of (20):

$$\min_d \frac{1}{2} K |d|^2 + \max \{ \nabla h_i(x_k), d \} + h_i(x_k) \mid i \text{ such that } h_i(x_k) \geq f(x_k) - \delta \} .$$

Using the same convexification technique as in (16)-(17), we see that in fact (20) is solved by

$$\bar{d} = -\frac{1}{K} \sum \bar{\lambda}_i \nabla h_i(x_k) \quad \text{where } \bar{\lambda} \text{ solves}$$

$$\min \frac{1}{2} \left| \sum \lambda_i \nabla h_i(x_k) \right|^2 - K \sum \lambda_i h_i(x_k) \quad \lambda_i \geq 0, \quad \sum \lambda_i = 1 \quad . \quad (21)$$

Again K can be considered as the (positive) Lagrange multiplier associated with some (inequality) constraint of the form $\sum \lambda_i h_i(x_k) \geq \beta$; this shows a strong relation with (17) and (18). We now make this relation more precise: any convex function $f(x)$ may be represented as

$$f(x) = \sup_{y, g \in \partial f(y)} f(y) + (g, x - y) \quad .$$

Therefore, any convex minimization problem is equivalent to

$$\begin{cases} \min_{v, x} & v \\ f(y) + (g, x - y) \leq v & \forall y \in H, \quad \forall g \in \partial f(y) \end{cases} \quad , \quad (22)$$

i.e. some sort of linear programming problem; of course this is highly criticizable since (22) contains an infinite number of constraints which, in addition, are not explicitly known. However, one could think of a method of constraint accumulation,

in which one would replace (22) by the relaxed problem

$$\begin{aligned} \min_{v, x} \quad & v \\ f_i + (g_i, x - x_i) & \leq v \quad i = 1, \dots, k \end{aligned}$$

where x_i , f_i , g_i make up the bundle of Section 3.

Analogously with the present situation, we can set

$$h_i(x) = f_i + (g_i, x - x_i) \quad , \quad \nabla h_i(x) = g_i \quad .$$

We then observe that

$$h_i(x_k) = f(x_k) - \alpha_i \quad (\alpha_i \text{ defined in (14)}) \quad .$$

Therefore (21) takes the form

$$\min \frac{1}{2} \left| \sum \lambda_i g_i \right|^2 \quad \sum \lambda_i = 1, \quad \lambda_i \geq 0, \quad \sum \lambda_i [f(x_k) - \alpha_i] \geq \beta \quad ,$$

and the equivalence with (18) is established by taking $\beta = f(x_k) - \epsilon$, so that K in (20) is u in (18).

The role of δ in (20) is to neglect the subgradients for which α_i is too large, i.e., it is exactly the deletion rule Section 3(c) of Mifflin.

Observe in particular that when $K = 0$ in (21), $\epsilon = +\infty$ in (18). Then the direction of Pshenichnyi is that of Demjanov [2]. In other words, Mifflin's version of conjugate subgradients is a variant of Demjanov's method, in which one computes not the values $h_i(x_k)$ that make up $f(x_k)$, but rather the values at x_k of the linearized f at previous x_i (this has been pointed out to us by Bertsekas in a private communication). Although (21) with $K \neq 0$ is apparently only a slight modification of Demjanov (at least when δ is small, since one has $f(x_k) - \delta \leq h_i(x_k) \leq f(x_k)$), taking the differences in $h_i(x_k)$ into account is a refinement.

(e) It is rather amusing to see that such a refinement has existed for quite some time in a slightly different context: Consider an ordinary mathematical programming problem,

$$\begin{aligned} \min f(x) \\ c_j(x) \leq 0 \quad j = 1, \dots, m \end{aligned} \quad (23)$$

Many methods for solving it are based on the principle of feasible directions. In their original form, presented by Zoutendijk [15], they consisted in solving at each iteration the problem

$$\left\{ \begin{array}{l} \min \quad \eta \\ \eta, d \\ (\nabla f(x_k), d) \leq \eta \\ \theta_j (\nabla c_j(x_k), d) \leq \eta \quad j \text{ in some suitable subset of } \{1, \dots, m\} \\ d \text{ in some suitable normalization set.} \end{array} \right. \quad (24)$$

It is certainly numerically necessary to introduce the weights θ_j , since there is no reason to compare variations in the constraints with variations in the objective, even if one admits that the user himself has scaled the space H to strive to homogenize variations in the objective alone. ||

As with (20), this problem can be transformed to look like (21). Since the values of the constraints $c_j(x_k)$ are neglected in (24), we would get (21) with $K = 0$, i.e. some Demjanov or conjugate subgradient method. It would therefore be better to consider in (24) the constraints

$$\theta_j [c_j(x_k) + (\nabla c_j(x_k), d)] \leq \eta \quad .$$

This has been done by Zoutendijk [15] himself, and (with $\theta_j \equiv 1$) by Topkis & Veinott [13], Pironneau & Polak [10] and Mangasarian

[6]. If, as is often true, the $c_j(x_k)$ are slightly negative numbers, the difference is little, but, from the analysis of rates of convergence in [10], the trick is really helpful (as is the use of Euclidean norm, as shown in the same paper). Correspondingly, the algorithm of Pshenichnyi for solving (23) (see his paper in this volume) introduces the same trick. (In addition, it is more satisfactory, since it does not need x_k to be feasible.)

To sum up this Section, (22) establishes a link between non-differentiable optimization and nonlinear programming. A non-smooth problem is a linear program with an infinite number of noncomputable constraints. To solve it, one could take advantage of the large amount of work done in nonlinear programming. Conversely, any efficient method of NSO could help in studying the (unsolved) general problem of nonlinear programming (23); for the latter, the theory of ϵ -subgradients is an apparently new and promising tool.

Most of the methods existing both in NSO and in NLP are essentially equivalent, provided the parameters they generally contain are carefully updated.

Through (12), the theory of ϵ -subgradients seems to shed some light on the fundamental question of scaling the space. It has been observed that the best numerical results are often obtained by quasi-Newton methods and Shor's method of dilatation of the space along differences of gradients [12]. Both these methods (heuristic in NSO) define at each iteration a metric H_k , assumed also to scale the variables. However, the formulae for updating H_k are *off-line* (in particular they do not involve values of the objective). It would probably be important to know what kind of relationship exists between ϵ -subgradients (i.e. general NLP), quasi-Newton methods, and Shor's methods.

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A FEASIBLE DESCENT ALGORITHM FOR LINEARLY CONSTRAINED LEAST SQUARES PROBLEMS

Robert Mifflin

1. INTRODUCTION

Consider the *constrained least squares problem* of finding an n -vector $x = (x_1, x_2, \dots, x_n)^T$ to

$$\text{minimize} \quad \frac{1}{2} |Px - c|^2 = \frac{1}{2} x^T P^T P x - c^T P x + \frac{1}{2} c^T c$$

$$\text{subject to} \quad Ax = b$$

$$x \geq 0$$

where P is a $p \times n$ matrix, c is a p -vector, A is an $m \times n$ matrix, b is an m -vector and a superscript T denotes transposition. The column rank of P may be less than n . We give a numerically stable method for solving this problem based on one given by Wolfe [6] for the special case where $c = 0$, $m = 1$, every component of the row vector A is 1 and $b = 1$. The algorithm solves a sequence of reduced dimension subproblems without nonnegativity constraints. The method is similar in philosophy to one given by Stoer [5], but our procedure for solving the subproblems, which is inspired by the work of Golub and Saunders [2], is different. The algorithm handles equality constraints directly; i.e., we do not use them to eliminate variables and then create inequality constraints from the corresponding nonnegativity restrictions as is suggested by Lawson and Hanson [3] in order to apply their procedure for inequality constrained problems, which involves yet another problem transformation.

We note that general quadratic programming algorithms may be applied to this problem, but, if they do not exploit the

factorized structure of the Hessian matrix $P^T P$ to deal with its possible singularity (or near singularity), such methods may fail.

Our method can also be used to solve strictly convex quadratic programming problems by transforming the objective function as shown in [2, p. 248]. Inequality constraints may be handled by introducing nonnegative slack or surplus variables. To modify the algorithm to deal with variables x_i that are not restricted to be nonnegative, one only needs to work out the implications of expressing such a variable as the difference of two nonnegative variables.

This study was motivated by the fact that Lemarechal's [4] numerical algorithm for minimizing a nondifferentiable function needs a numerically stable subroutine for solving the special constrained least squares problem where

$$n \geq 2, \quad m = 2, \quad c = 0,$$

$$A = \begin{bmatrix} 1 & \alpha_1 & \alpha_2 & \cdots & \alpha_{n-1} \\ 0 & 1 & 1 & \cdots & 1 \end{bmatrix}, \quad b = \begin{bmatrix} \epsilon \\ 1 \end{bmatrix},$$

$$\alpha_i \geq 0 \quad \text{for } i = 1, 2, \dots, n-1, \quad \alpha_1 < \epsilon,$$

and

$$P = \begin{bmatrix} 0 \\ 0 \\ \vdots \\ 0 \end{bmatrix} \begin{matrix} g_1 & g_2 & \cdots & g_{n-1} \end{matrix}.$$

In this application the p -vectors g_1, g_2, \dots, g_{n-1} are generalized gradients of a function of p variables to be minimized, and the solution method employed must not require these vectors to be linearly independent.

2. THE ALGORITHM

For a matrix B we denote its i th row by B_i and its j th column by B^j . Thus, $B^j = B_j^T$. For a vector $y = (y_1, y_2, \dots, y_\ell)^T$, $y > 0$ means $y_i > 0$ for each $i = 1, 2, \dots, \ell$.

For x to be an optimal solution to the above constrained least squares problem it is both necessary and sufficient (because of objective convexity and constraint linearity) that x and some m -vector u satisfy the *optimality conditions*:

$$Ax = b$$

$$x \geq 0$$

$$A^T u + P^T (Px - c) \geq 0$$

$$A_j^T u + P_j^T (Px - c) = 0 \quad \text{or} \quad x_j = 0 \quad \text{for each } j = 1, 2, \dots, n \quad .$$

Throughout, $J \subseteq \{1, 2, \dots, n\}$ denotes a nonempty set of column indices corresponding to a nonvacuous submatrix $\begin{bmatrix} B \\ Q \end{bmatrix}$ of columns of $\begin{bmatrix} A \\ P \end{bmatrix}$; i.e.,

$$j \in J \quad \text{if and only if} \quad \begin{bmatrix} A^j \\ P^j \end{bmatrix} \quad \text{is a column of} \quad \begin{bmatrix} B \\ Q \end{bmatrix} \quad .$$

Each submatrix generated by the algorithm satisfies the assumption that

$$\begin{bmatrix} B \\ Q \end{bmatrix} \quad \text{has full column rank} \quad . \quad (A1)$$

For each such matrix we require a solution (y, u) to the corresponding linear system

$$By = b \quad (1a)$$

$$B^T u + Q^T Q y = Q^T c \quad . \quad (1b)$$

A solution procedure is given in Section 3 for the *nondegenerate case* when

$$B \text{ has full row rank } . \quad (A2)$$

This type of assumption is also made in [5]. For the special problem with $m = 1$ in [6], (A2) is always satisfied. For the case when (A2) does not hold, a degeneracy resolving procedure that takes into account possible nonuniqueness of u has been worked out by the author. This procedure may appear in a subsequent paper.

The following sequence of lemmas serves to motivate the algorithm. All proofs are deferred to the last section of this paper.

Lemma 1

Suppose (A1) holds. If (1) has a solution then y is unique. Furthermore, if (A2) also holds then (1) has a unique solution (y,u) and the solution procedure of Section 3 finds it.

Remark

For the case when (A2) does not hold, the solution procedure can be modified to determine whether or not (1) has a solution, and in the former case to find y and the associated u -set. It can be shown that if (A1) holds and there exists a vector \bar{y} such that $B\bar{y} = b$, then (1) has a solution.

Lemma 2

(y,u) solves (1) if and only if y minimizes $|Qy - c|^2$ subject to $By = b$. In this case

$$|Qy - c|^2 = -b^T u - c^T (Qy - c) . \quad (2)$$

Lemma 3

Suppose (y, u) solves (1). Define a corresponding n -vector x by

$$x_j = \begin{cases} y_j & \text{if } j \in J \\ 0 & \text{if } j \notin J \end{cases}$$

where J corresponds to $\begin{bmatrix} B \\ Q \end{bmatrix}$ as defined above. If

$$y \geq 0$$

and

$$A_j^T u + P_j^T (Qy - c) \geq 0 \quad \text{for all } j \notin J, \quad (3)$$

then x and u satisfy the optimality conditions and, hence, x is an optimal solution to the constrained least squares problem.

Remark

For the special problem of Section 1 with $m = 2$ and $u = (u_1, u_2)^T$, (2) becomes

$$|Qy|^2 = -\epsilon u_1 - u_2$$

and the left-hand side of (3) becomes

$$u_1 \quad \text{if } j = 1$$

or

$$\alpha_{j-1} u_1 + u_2 + g_{j-1}^T Qy \quad \text{if } j \geq 2.$$

Lemma 4

Suppose (A1) holds and (y, u) solves (1). If

$$A_{\ell}^T u + P_{\ell}^T (QY - c) < 0 \quad \text{for some } \ell \notin J \quad (4)$$

then $\begin{bmatrix} B & A^{\ell} \\ Q & P^{\ell} \end{bmatrix}$ has full column rank.

Remark

Note that if B has full row rank then so does $[B \ A^{\ell}]$. Thus, by Lemma 4, if (A1) and (A2) hold for $\begin{bmatrix} B \\ Q \end{bmatrix}$ and B , respectively, then they hold for the augmented matrices $\begin{bmatrix} B & A^{\ell} \\ Q & P^{\ell} \end{bmatrix}$ and $[B \ A^{\ell}]$.

These augmented matrices correspond to the augmented linear system

$$By^+ + A^{\ell} y_{\ell}^+ = b \quad (1a)^+$$

$$B^T u^+ + Q^T (QY^+ + P^{\ell} y_{\ell}^+) = Q^T c \quad (1b)^+$$

$$A_{\ell}^T u^+ + P_{\ell}^T (QY^+ + P^{\ell} y_{\ell}^+) = P_{\ell}^T c \quad (1c)^+$$

Lemma 5

Suppose (A1) and (A2) hold, (y, u) solves (1), and (4) holds. Then (1)⁺ has a solution (y^+, y_{ℓ}^+, u^+) such that

$$y_{\ell}^+ > 0$$

and

$$|QY^+ + P^{\ell} y_{\ell}^+ - c|^2 - |QY - c|^2 = y_{\ell}^+ \left(A_{\ell}^T u + P_{\ell}^T (QY - c) \right) < 0 \quad (5)$$

Lemma 6

Suppose that the assumptions of Lemma 5 hold and that $y > 0$. Let $\bar{\lambda}$ be the largest value of $\lambda \in [0, 1]$ such that

$$\lambda \begin{bmatrix} y^+ \\ y_\ell^+ \end{bmatrix} + (1 - \lambda) \begin{bmatrix} y \\ 0 \end{bmatrix} \geq 0 .$$

Define the $|J| + 1$ vector z by

$$z_j = \begin{cases} \bar{\lambda} y_j^+ + (1 - \bar{\lambda}) y_j & \text{for } j \in J \\ \bar{\lambda} y_\ell^+ & \text{for } j = \ell . \end{cases}$$

Then $\bar{\lambda} > 0$, $z_\ell > 0$, $z \geq 0$,

$$[B \ A^\ell] z = b \quad (6)$$

and

$$|[Q \ P^\ell] z - c|^2 < |Qy - c|^2 . \quad (7)$$

Furthermore, if $y^+ > 0$, then $\bar{\lambda} = 1$ and $z = \begin{bmatrix} y^+ \\ y_\ell^+ \end{bmatrix} > 0$, or if $y^+ \not> 0$ then $z_k = 0$ for some $k \in J$.

To initiate the algorithm we assume that we have an index set J and a corresponding full column rank matrix $\begin{bmatrix} B \\ Q \end{bmatrix}$ such that (1) has a solution with $y > 0$. See Section 5 for an initialization phase that will either provide this initial condition or indicate its infeasibility. Given this initial condition, the algorithm is as follows.

Step 1 Solve (1) for (y, u) .
If $y > 0$ go to Step 2. Otherwise go to Step 3.

Step 2 Test the validity of (3).
If (3) holds stop. Otherwise ((4) holds)
set $\bar{y}_i = y_i$ for each $i \in J$ and $\bar{y}_\ell = 0$, append index ℓ
to J and column $\begin{bmatrix} A^\ell \\ P^\ell \end{bmatrix}$ to $\begin{bmatrix} B \\ Q \end{bmatrix}$, and go to Step 1.

Step 3 Let $\bar{\lambda}$ be the maximum value of $\lambda \in [0,1]$ such that

$$\lambda y + (1 - \lambda)\bar{y} \geq 0$$

and set

$$z = \bar{\lambda}y + (1 - \bar{\lambda})\bar{y} .$$

For each i such that $z_i > 0$, set $\bar{y}_i = z_i$, and for each k such that $z_k = 0$, delete index k from J and the corresponding column $\begin{bmatrix} A^k \\ P^k \end{bmatrix}$ from $\begin{bmatrix} B \\ Q \end{bmatrix}$ and go to Step 1.

We now show that the algorithm is well defined.

Because of the initialization condition the first subproblem (1) to be solved has a solution $y > 0$, and therefore Step 2 is executed before Step 3, so \bar{y} is properly initialized to a non-negative-valued vector.

At each execution of Step 2, if the stop does not occur, a new index ℓ is added to J satisfying (4) and \bar{y} and B are updated such that $\bar{y} \geq 0$ and $B\bar{y} = b$. From Lemma 4 and the remark following Lemma 1 the new system (1) will have a solution.

When entering Step 3 we have J , $y \not\geq 0$ and \bar{y} . By the updating in Steps 2 and 3, $\bar{y} \geq 0$. Therefore, there exists a $\lambda \in [0,1]$ such that $\lambda y + (1 - \lambda)\bar{y} \geq 0$ (for example $\lambda = 0$). From the definitions of $\bar{\lambda}$ and z , there exists $k \in J$ such that $z_k = 0$. Hence, at least index k is deleted from J . A new J and a new $\bar{y} > 0$ are defined. Lemma 7, in particular, ensures that the new J is not empty and that the new system (1) has a solution. This lemma deals with the situation after an execution of Step 3 in terms of data at the most recent execution of Step 2.

Lemma 7

Suppose at Step 3 we have $\begin{bmatrix} B \\ Q \end{bmatrix}$ satisfying (A1), and

$$(H) \left\{ \begin{array}{l} z \geq 0 \quad , \\ Bz = b \quad , \\ z_{\ell^{\circ}} > 0 \quad , \\ z_k = 0 \quad \text{for some } k \neq \ell^{\circ} \quad ; \\ |Qz - c|^2 < |Q^{\circ}y^{\circ} - c|^2 \end{array} \right.$$

where, at the most recent execution of Step 2, J° , B° , Q° and y° were the entering values of J , B , Q and y , and ℓ° was the index added to J° . Let the new values of J , B and Q defined at Step 3 be denoted J^{-} , B^{-} and Q^{-} , respectively. Then the new system (1) has a solution $y = y^{-}$ such that either

$$y^{-} > 0 \quad \text{and} \quad |Q^{-}y^{-} - c|^2 < |Q^{\circ}y^{\circ} - c|^2 \quad (8)$$

or $y^{-} \not> 0$ and the above hypotheses (H) are satisfied with $z = z^{-}$, $B = B^{-}$ and $Q = Q^{-}$ where z^{-} is the value of z determined at the next execution of Step 3.

By recursive application of Lemmas 1 through 7 it is now not difficult to establish finite convergence of the algorithm.

Theorem

Suppose the initial condition holds and (A2) holds for each matrix B generated by the algorithm. Then after a finite number of executions of Steps 1, 2 and 3, the algorithm terminates with an optimal solution to the constrained least squares problem.

Proof

To each J considered at Step 1 there corresponds a unique y solution of (1) and a corresponding objective value $\frac{1}{2}|Qy - c|^2 = \frac{1}{2}|Px - c|^2$ where x is the n -vector defined by appending zeros to y as in Lemma 3.

We now show that each entry to Step 2 has an associated objective value that is strictly lower than the one at the previous entry to Step 2, and, furthermore, that the number of consecutive executions of Step 3 between Step 2 executions is finite.

By Lemma 6, the above statement is clearly true if there are no intermediate Step 3 executions. So suppose Step 3 is entered after some execution of Step 2 which adds index λ^0 to J . By Lemma 6 we have at this first entry to Step 3 that the hypotheses (H) of Lemma 7 are satisfied. From Lemma 7, for all subsequent consecutive executions of Step 3 we have $z_{\lambda^0} > 0$ and some index $k \neq \lambda^0$ is removed from J . Therefore, J never becomes empty, the number of such consecutive Step 3 executions is finite, and furthermore, by Lemma 7, this sequence must terminate with (8) being satisfied. Now (8) implies a return to Step 2 with a strictly improved objective value.

Now, since the number of possible index sets J is finite and all such sets corresponding to entries to Step 2 must be distinct (due to their different corresponding objective values), the algorithm is finite.

3. SOLUTION PROCEDURE FOR (1)

Suppose (A1) and (A2) hold and $\begin{bmatrix} B \\ Q \end{bmatrix}$ has q columns. Let R be an upper triangular nonsingular $q \times q$ matrix such that

$$\begin{bmatrix} B^T & Q^T \end{bmatrix} \begin{bmatrix} B \\ Q \end{bmatrix} = R^T R \quad , \quad (9)$$

and let W be a $q \times m$ matrix such that

$$R^T W = B^T \quad , \quad (10)$$

ω be a q -vector such that

$$R^T \omega = Q^T c \quad , \quad (11)$$

and S be an upper triangular nonsingular $m \times m$ matrix such that

$$W^T W = S^T S \quad . \quad (12)$$

R and S can be found by orthogonal factorization (see, for example, [1,2,3,5]) of $\begin{bmatrix} B \\ Q \end{bmatrix}$ and W , respectively, which is possible by the full rank assumptions (A1) and (A2). See Section 5 for simple determinations of R , W , ω and S in some special cases and see the next section for updating these quantities when $\begin{bmatrix} B \\ Q \end{bmatrix}$ is changed. Having R , W , ω and S , the solution procedure is as follows:

Solve the triangular linear system

$$S^T w = b - W^T \omega$$

for w , and then solve the triangular linear system

$$Sv = w$$

for v so that v satisfies

$$S^T S v = b - W^T \omega \quad . \quad (13)$$

Then set

$$u = b - v \quad (14)$$

and solve the triangular linear system

$$Ry = Wv + \omega \quad (15)$$

for y .

4. UPDATING

To perform the updating of R , W , ω and S we require a transformation that can take a 2-vector

$$\begin{bmatrix} z_1 \\ z_2 \end{bmatrix} \text{ into } \gamma \begin{bmatrix} 1 \\ 0 \end{bmatrix}$$

where

$$\gamma^2 = z_1^2 + z_2^2 .$$

This Euclidean norm-preserving transformation can be accomplished by multiplying $\begin{bmatrix} z_1 \\ z_2 \end{bmatrix}$ on the left by a *Givens matrix*

$$G = \begin{bmatrix} c & s \\ s & -c \end{bmatrix}$$

where

$$\gamma = \text{sign}(z_1) \left(z_1^2 + z_2^2 \right)^{\frac{1}{2}} ,$$

$$c = z_1 / \gamma ,$$

$$s = z_2 / \gamma ,$$

and the convention for $z_1 = 0$ is that $\gamma = z_2$, $c = 0$ and $s = 1$. Note that $G^T = G^{-1}$, i.e. G is an *orthogonal* matrix. See Gill, Golub, Murray and Saunders [1] for details concerning properties and uses of Givens transformations in connection with orthogonal factorization.

Augmentation

When $\begin{bmatrix} B \\ Q \end{bmatrix}$ is replaced at Step 2 by $\begin{bmatrix} B & A^l \\ Q & P^l \end{bmatrix}$, replace

$$R \text{ by } \begin{bmatrix} R & r \\ 0 & \rho \end{bmatrix} ,$$

$$W \text{ by } \begin{bmatrix} W \\ \bar{w} \end{bmatrix} ,$$

$$\omega \text{ by } \begin{bmatrix} \omega \\ \bar{\omega} \end{bmatrix} ,$$

$$S \text{ by } \bar{S} ,$$

where r solves

$$R^T r = B^T A^\ell + Q^T P^\ell ,$$

$$\rho = \left(|A^\ell|^2 + |P^\ell|^2 - |r|^2 \right)^{\frac{1}{2}} ,$$

$$\bar{w} = (A_\ell^T - r^T W) / \rho ,$$

$$\bar{\omega} = (P_\ell^T c - r^T \omega) / \rho ,$$

and \bar{S} is determined as in [1, pp. 529-530] as follows.

Apply appropriate 2×2 Givens transformations sequentially for $i = 1, 2, \dots, m$ to rows i and $m + 1$ of

$$\begin{bmatrix} S \\ \bar{w} \end{bmatrix}$$

to reduce \bar{w} to the zero vector and obtain

$$\begin{bmatrix} \bar{S} \\ 0 \end{bmatrix}$$

where \bar{S} is upper triangular.

The validity of these replacements may be established by making the appropriate multiplications and using the fact that a sequence of Givens transformations is equivalent to multiplication by an orthogonal matrix.

Deletion

When the k th column of the q column matrix $\begin{bmatrix} B \\ Q \end{bmatrix}$ is deleted at Step 3 replace R, W, ω and S by $\bar{R}, \bar{W}, \bar{\omega}$ and \bar{S} , respectively, where the latter quantities are determined as follows.

As in [1, pp. 533-534], apply appropriate 2×2 Givens transformations sequentially for $i = k, k+1, \dots, q-1$ to rows i and $i + 1$ of

$$[R^1 R^2 \dots R^{k-1} R^{k+1} \dots R^q \quad W \quad \omega]$$

to form

$$\begin{bmatrix} \bar{R} & \bar{W} & \bar{\omega} \\ 0 & \underline{w} & \underline{\omega} \end{bmatrix}$$

where

\bar{R} is a $q-1 \times q-1$ upper triangular matrix,

$[\bar{W} \quad \bar{\omega}]$ is $q-1 \times m+1$,

and $(\underline{w} \quad \underline{\omega})$ is $1 \times m+1$.

Then, as in [1, pp. 530-531], solve the triangular linear system

$$S^T s = \underline{w}^T ;$$

for s , set

$$\delta = \left(1 - |s|^2\right)^{\frac{1}{2}} ;$$

and apply appropriate 2×2 Givens transformations sequentially for $i = m, m-1, \dots, 1$ to rows i and $m + 1$ of

$$\begin{bmatrix} s & S \\ \delta & 0 \end{bmatrix}$$

to reduce s to the zero vector and form

$$\begin{bmatrix} 0 & \bar{S} \\ 1 & s^T \end{bmatrix}$$

where \bar{S} is upper triangular, and the expression for the bottom row, as well as the validity of \bar{R} , \bar{W} , $\bar{\omega}$ and \bar{S} , follows from the orthogonality of Givens transformations.

5. INITIALIZATION

For the special problem with $m = 2$ given in Section 1, a starting matrix is

$$\begin{bmatrix} B \\ Q \end{bmatrix} = \begin{bmatrix} 1 & \alpha_1 \\ 0 & 1 \\ 0 & g_1 \end{bmatrix}$$

for which

$$R = \begin{bmatrix} 1 & \alpha_1 \\ 0 & (1 + |g_1|^2)^{\frac{1}{2}} \end{bmatrix}$$

$$W = S = \begin{bmatrix} 1 & 0 \\ 0 & (1 + |g_1|^2)^{-\frac{1}{2}} \end{bmatrix}, \quad \omega = 0$$

and

$$y = \begin{bmatrix} \epsilon^{-\alpha_1} \\ 1 \end{bmatrix} > 0.$$

For the general problem, if a starting matrix is not available we can first solve the initialization problem of

$$\text{minimizing } \left| \begin{bmatrix} -b & A \end{bmatrix} \begin{bmatrix} t \\ x \end{bmatrix} \right|^2$$

$$\begin{aligned} \text{subject to } t &= 1 \\ t, x &\geq 0 . \end{aligned}$$

For this problem a starting matrix is

$$\begin{bmatrix} B \\ Q \end{bmatrix} = \begin{bmatrix} 1 \\ -b \end{bmatrix}$$

for which

$$R = W^{-1} = S^{-1} = \left(1 + |b|^2 \right)^{\frac{1}{2}} , \quad \omega = 0 \quad \text{and} \quad y = 1 .$$

Note that this initialization problem is never degenerate, since all row vectors B generated must contain the coefficient 1 corresponding to variable t . Our algorithm applied to this special initialization problem is essentially that of Lawson and Hanson [3, p. 161] for the problem of minimizing $|Ax - b|^2$ subject to $x \geq 0$.

If the optimal value of the initialization problem is positive then the constrained least squares problem has no feasible solution. Suppose this optimal value is zero. Then the corresponding optimal x satisfies $x \geq 0$ and $Ax = b$. If $x \neq 0$ the set of columns A^j for which $x_j > 0$ will be linearly independent, and together with the corresponding columns P^j will form a starting matrix for the desired problem. If $x = 0$ then $b = 0$, and

$$\begin{bmatrix} B \\ Q \end{bmatrix} = \begin{bmatrix} 1 \\ 0 \\ 0 \end{bmatrix}$$

is a starting matrix for the equivalent problem of

$$\begin{aligned} & \text{minimizing } \left| [0 \ P] \begin{bmatrix} t \\ x \end{bmatrix} - c \right|^2 \\ & \text{subject to } \begin{bmatrix} 1 & 0 \\ 0 & A \end{bmatrix} \begin{bmatrix} t \\ x \end{bmatrix} = \begin{bmatrix} 1 \\ 0 \end{bmatrix} . \\ & t, x \geq 0 . \end{aligned}$$

This starting matrix is, however, degenerate, because $B = \begin{bmatrix} 1 \\ 0 \end{bmatrix}$ does not satisfy (A2).

6. PROOFS OF LEMMAS

Lemma 1

Suppose (A1) holds. Let (y^1, u^1) and (y^2, u^2) be solutions to (1), i.e.,

$$\begin{aligned} By^1 &= b & \text{and} & & By^2 &= b \\ B^T u^1 + Q^T Q y^1 &= Q^T c & & & B^T u^2 + Q^T Q y^2 &= Q^T c . \end{aligned}$$

Subtraction gives

$$B(y^1 - y^2) = 0 \tag{16}$$

and

$$B^T(u^1 - u^2) + Q^T Q(y^1 - y^2) = 0 . \tag{17}$$

Multiplying (17) by $(y^1 - y^2)^T$ and using the transpose of (16) gives

$$(y^1 - y^2)^T Q^T Q (y^1 - y^2) = 0 , \tag{18}$$

and (16) also implies that

$$(y^1 - y^2)^T B^T B (y^1 - y^2) = 0 ,$$

which added to (18) gives

$$(y^1 - y^2)^T [B^T Q^T] \begin{bmatrix} B \\ Q \end{bmatrix} (y^1 - y^2) = 0 \quad .$$

So

$$\begin{bmatrix} B \\ Q \end{bmatrix} (y^1 - y^2) = 0 \quad ,$$

and (A1) implies that

$$y^1 = y^2$$

which, together with (17), implies that

$$B^T (u^1 - u^2) = 0 \quad .$$

Suppose, in addition to (A1), that (A2) holds. Then

$$u^1 = u^2 \quad .$$

Therefore, under assumptions (A1) and (A2), if (1) has a solution it is unique. We now show that (y, u) determined by the procedure of Section 3 is a solution to (1). Note that (15), (10), (12) and (13) imply sequentially that

$$\begin{aligned} By &= BR^{-1}Wv + BR^{-1}\omega \\ &= W^T Wv + W^T \omega \\ &= S^T Sv + W^T \omega \\ &= b - W^T \omega + W^T \omega \quad , \end{aligned}$$

so

$$By = b \quad .$$

(19)

By (14) and (9) we have

$$B^T u + Q^T Q y = B^T (b - v) + (R^T R - B^T B) y .$$

Now (19), (15), (10) and (11) imply sequentially that

$$\begin{aligned} B^T u + Q^T Q y &= B^T b - B^T v + R^T R y - B^T b \\ &= -B^T v + R^T (Wv + \omega) \\ &= -B^T v + B^T v + Q^T c , \end{aligned}$$

so (y, u) solves (1).

Lemma 2

Since $\frac{1}{2}|Qy - c|^2$ is convex in y and By is linear in y , the conditions given by (1) are well known to be both necessary and sufficient for optimality in the problem of minimizing $\frac{1}{2}|Qy - c|^2$ subject to $By = b$. Multiplying on the left (1a) by u^T and (1b) by y^T , and differencing the resulting equations, give

$$u^T B y - y^T B^T u - y^T Q^T Q y = u^T b - y^T Q^T c$$

which, by scalar transposition, is equivalent to

$$y^T Q^T Q y = -b^T u + c^T Q y$$

which is equivalent to (2).

Lemma 3

Combining the definition of x and (1) gives

$$Ax = By = b ,$$

$$Px = Qy \tag{20}$$

and

$$B^T u + Q^T(Px - c) = B^T u + Q^T(Qy - c) = 0 \quad . \quad (21)$$

Also, $y \geq 0$ implies $x \geq 0$. Combining (20), (21) and (3) implies

$$A^T_j u + P^T_j(Px - c) \geq 0 \quad \text{for each } j = 1, 2, \dots, n$$

with strict inequality only if $x_j = 0$.

Thus, x and u satisfy the optimality conditions.

Lemma 4

Suppose, for purposes of a proof by contradiction, that

$$\begin{bmatrix} B & A^\ell \\ Q & P^\ell \end{bmatrix} \begin{bmatrix} \lambda \\ \lambda_\ell \end{bmatrix} = 0 \quad (22)$$

for some nonzero vector $(\lambda, \lambda_\ell)^T$. Note that

$$\lambda_\ell \neq 0 \quad , \quad (23)$$

because, by assumption, $\begin{bmatrix} B \\ Q \end{bmatrix}$ has full column rank. Multiplying (22) on the left by $(u^T, y^T Q^T - c^T)$ gives

$$\left(u^T B + (y^T Q^T - c^T) Q \right) \lambda + \left(u^T A^\ell + (y^T Q^T - c^T) P^\ell \right) \lambda_\ell = 0 \quad .$$

But this is a contradiction, because the first term is zero by the transpose of (1b) and the second term is nonzero by (23) and the transpose of (4).

Lemma 5

By the remark following Lemma 4, Lemma 1 applied to $(1)^+$ implies that $(1)^+$ has a solution (y^+, y_ℓ^+, u^+) . Since

$$[B \ A^\ell] \begin{bmatrix} y \\ 0 \end{bmatrix} = b \quad ,$$

Lemma 2 applied to $(1)^+$ implies that

$$\left| [Q \ P^\ell] \begin{bmatrix} y^+ \\ y_\ell^+ \end{bmatrix} - c \right|^2 \leq \left| [Q \ P^\ell] \begin{bmatrix} y \\ 0 \end{bmatrix} - c \right|^2 = |Qy - c|^2 \quad . \quad (24)$$

By (2) of Lemma 2 applied to $(1)^+$ and (1), we have

$$\begin{aligned} & \left| [Q \ P^\ell] \begin{bmatrix} y^+ \\ y_\ell^+ \end{bmatrix} - c \right|^2 - |Qy - c|^2 = \\ & = -b^T u^+ - c^T (Qy^+ + P^\ell y_\ell^+ - c) + b^T u + c^T (Qy - c) \\ & = -b^T u^+ - c^T Qy^+ - c^T P^\ell y_\ell^+ + b^T u + c^T Qy \quad . \end{aligned} \quad (25)$$

Multiplying the transposes of (1a) and (1b) on the right by u^+ and y^+ , respectively, and $(1a)^+$, $(1b)^+$ and $(1c)^+$ on the left by $-u^T$, $-y^T$ and y_ℓ^+ , respectively, and then adding the resulting equations gives

$$\begin{aligned} & y^T B u^+ + u^T B y^+ + y^T Q^T Q y^+ - u^T B y^+ - u^T A y_\ell^+ - y^T B^T u^+ - \\ & - y^T Q^T Q y^+ - y^T Q^T P^\ell y_\ell^+ = b^T u^+ + c^T Q y^+ - u^T b - y^T Q^T c \end{aligned}$$

which, after cancellation and rearrangement, is equivalent to

$$0 = u^T A^\ell y_\ell^+ + y^T Q^T P^\ell y_\ell^+ + b^T u^+ + c^T Q y^+ - u^T b - y^T Q^T c \quad .$$

Adding this to (25) and simplifying gives

$$\left| [Q \ P^\ell] \begin{bmatrix} y^+ \\ y_\ell^+ \end{bmatrix} - c \right|^2 - |Qy - c|^2 = (u^T A^\ell + y^T Q^T P^\ell - c^T P^\ell) y_\ell^+ \quad (26)$$

which, by transposition, is equivalent to the equality in (5). Now (24), (26) and (4) imply that $y_\ell^+ \geq 0$. Suppose $y_\ell^+ = 0$. Then, by (1a)⁺ and (1b)⁺, (y^+, u^+) solves (1). But assumptions (A1) and (A2) imply, by Lemma 1, that $(y^+, u^+) = (y, u)$. Then, since we suppose $y_\ell^+ = 0$, (4) and (1c)⁺ are contradictory. Therefore, $y_\ell^+ > 0$, which together with (26) and (4) gives the desired inequality in (5) and completes the proof.

Lemma 6

Since $y > 0$ and $y_\ell^+ > 0$, $\bar{\lambda}$ is well defined by

$$\bar{\lambda} = \min \left\{ 1, \min \left[y_j / (y_j - y_j^+) : y_j - y_j^+ > 0, j \in J \right] \right\}$$

and we have $\bar{\lambda} > 0$, $z_\ell = \bar{\lambda} y_\ell^+ > 0$ and $z_j \geq 0$ for all $j \in J$. Moreover, z has at least one zero component if and only if $y^+ \nmid 0$. The definition of z combined with (1a) and (1a)⁺ implies (6) and, combined with the convexity of quadratic functions having positive semidefinite Hessian matrices, implies that

$$\left| [Q P^\ell] z - c \right|^2 \leq \bar{\lambda} \left| [Q P^\ell] \begin{bmatrix} y^+ \\ y_\ell^+ \end{bmatrix} - c \right|^2 + (1 - \bar{\lambda}) \left| Q y - c \right|^2 . \quad (27)$$

Finally, (7) follows from (27) and (5).

Lemma 7

By construction at Step 3, \bar{y} satisfies

$$\bar{y} > 0 \quad \text{and} \quad \begin{bmatrix} B^- \\ Q^- \end{bmatrix} \bar{y} = \begin{bmatrix} B \\ Q \end{bmatrix} z . \quad (28)$$

Also, $J^- \subset J$, so $\begin{bmatrix} B^- \\ Q^- \end{bmatrix}$ has full column rank, because $\begin{bmatrix} B \\ Q \end{bmatrix}$ has full column rank. By (28) and (H) we have

$$B^- \bar{y} = b , \quad (29)$$

so, by the remark following Lemma 2, y^- is well defined and

$$B^- y^- = b \quad . \quad (30)$$

Therefore, by Lemma 2 and (28),

$$|Q^- y^- - c|^2 \leq |Q^- \bar{y} - c|^2 = |Qz - c|^2 \quad . \quad (31)$$

If $y^- > 0$, then (8) follows from (31) and the hypothesis that

$$|Qz - c|^2 < |Q^o y^o - c|^2 \quad . \quad (32)$$

Suppose $y^- \not> 0$. Then, because $\bar{y} > 0$, at the next execution of Step 3, $\bar{\lambda}$ and z^- are well defined and we have

$$\bar{\lambda} > 0 \quad ,$$

$$z^- = \bar{\lambda} y^- + (1 - \bar{\lambda}) \bar{y} \geq 0 \quad ,$$

and

$$z_k^- = 0 \quad \text{for some } k \quad .$$

Also, by (29) and (30), we have

$$B^- z^- = b \quad (33)$$

and, by objective convexity, (31) and (32), we have

$$|Q^- z^- - c|^2 \leq |Q^- \bar{y} - c|^2 = |Qz - c|^2 < |Q^o y^o - c|^2 \quad . \quad (34)$$

If $z_{\ell^o}^- = 0$ then, since $J^- \subset J^o \cup \{\ell^o\}$, (33) and Lemma 2 imply that

$$|Q^o y^o - c|^2 \leq |Q^- z^- - c|^2 \quad ,$$

which contradicts (34). Therefore, $z_{\ell^o}^- > 0$ and the proof is complete.

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SUFFICIENT MINIMIZATION OF PIECEWISE-LINEAR UNIVARIATE FUNCTIONS*

P. Wolfe

PRELIMINARIES

Minimization of a function of a single variable is a fundamental subroutine in most procedures for the minimization of functions of many variables. We now know of important classes of large-scale linear programming problems, having astronomically many variables, which can be recast as problems of minimizing convex, nonlinear but piecewise-linear, functions of reasonable numbers of variables. The "method of conjugate subgradients" is a procedure for solving the latter kind of problem, requiring as a subroutine an efficient procedure for finding a "one-sided" minimum of the given function (in the sense described later) on an arbitrary ray in the space of problem variables. In other words, a procedure is required which will efficiently find a certain kind of approximate minimum, if such exists, of any convex, piecewise-linear function f of a single variable. Further, in the given context of large-scale linear programming, the function f cannot be explicitly described: rather, given any value x of its argument, one can readily calculate $f(x)$ and one support of the graph of f --that is, a "slope", denoted here by $f'(x)$, such that

$$f(y) \geq f(x) + f'(x)(y - x) \quad \text{for all } y$$

--but no more. The kind of function handled here is unusual in one respect: $f'(x)$ is *not* continuous. Consequently, a method designed for the more usual problem, in which $f'(x)$ is continuous, does not work well in this case.

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The algorithm below accomplishes the required end with what we believe to be near the smallest required computational effort.

PROBLEM STATEMENT

The piecewise-linear (more properly, piecewise-affine) function f is defined for all $x \geq 0$ by

$$f(x) = \text{Max} \{a_i x + b_i : i = 1, \dots, I\} \quad , \quad (1)$$

where, for each i , a_i and b_i are real numbers. The quantities a_i and b_i are not explicitly given: rather we suppose that we have a procedure which, for any $x \geq 0$, determines the value $f(x)$ and some a_i, b_i for which $f(x) = a_i x + b_i$. We denote that a_i by $f'(x)$. To simplify notation below we will suppose that $f(0) = 0$, and also that $M = f'(0) < 0$ (since, otherwise $x=0$ is trivially the minimizer).

We know that exact solution of the univariate minimization problem is not needed for success of the overall procedure of which it is a subroutine. Much computational labor can be saved by requiring only that the two following conditions be satisfied by the approximate minimizer x :

$$(a) \quad f(x) \leq M m_2 x \quad ,$$

$$(b) \quad f'(x) \geq M m_1 \quad .$$

(The constants m_1 and m_2 must satisfy the relations $0 < m_2 < m_1 < 1.0$.) Our goal is thus to find $x \geq 0$, $f(x)$, and $f'(x)$ satisfying (a) and (b). The "one-sided" nature of the requirement (b) is distinctive for the kind of problem of interest here, as opposed to the more usual requirement

$$|f'(x)| \leq m_1 |M| \quad ;$$

For efficiency, this problem needs a substantially different algorithm from that used if the latter requirement must be met.

Using the convexity of f it is easy to show that the two functions

$$f(x)/x \quad , \quad f'(x)$$

are nondecreasing; so the set of all x satisfying (a) is an interval of the form $[0, A]$ (where $A=0$ or $A=\infty$ are both possible), and the set of all x satisfying (b) is the form $[B, \infty)$ (with the same possibilities for B). If $A=\infty$ then (a) holds for all x , and our procedure will generate a sequence of values of x tending to infinity, so that $f(x)$ will tend to $-\infty$. Otherwise it is easy to show that $B \leq A$, so that the two intervals overlap. The algorithm below finds a point in that overlap.

ALGORITHM

When f and f' are evaluated for some value of x , the data are recorded in the form of the triple $(x, f(x), f'(x))$. Supposing terminal data (i.e., satisfying both (a) and (b) above) not to have been found, L will denote the most recent triple evaluated for which x lies in the interval $[0, A]$, LL the previous such triple, and R the most recent, if any, for the interval $[B, \infty)$. $Li(x, y, m)$ will denote the line of slope m passing through the point (x, y) ; thus $Li(L)$ is a line of slope $f'(x)$ passing through a point $(x, f(x))$. The "stepping factor" E is some number greater than one; we have found the value $E=6$ useful. Component j of a triple like L will be denoted $L[j]$, as in the APL language. The relations $LL[1] \leq L[1] < R[1]$ will hold throughout the course of the algorithm.

It is supposed that a "starting" value $x > 0$ has been provided for the algorithm. Initially L and LL are set to $(0, 0, f'(0))$ and R to (∞, ∞, ∞) . The algorithm consists in execution of the steps below in order (except where "go to" intervenes).

- 1.1 Evaluate $f(x)$, $f'(x)$.
- 1.2 If (a) does not hold, go to 4.
- 1.3 If (b) holds, terminate; x, f, f' is the solution.

- 2.1 Set LL to L.
- 2.2 Set L to $(x, f(x), f'(x))$.
- 2.3 If $R[1] < \infty$, go to 5.
3. (Extrapolation) Use linear extrapolation on the slopes $LL[3]$, $L[3]$ to determine the value x^* (which may be ∞) at which $f'(x^*)$ should vanish. Set x to $\text{Min}\{x^*, x + E(x - LL[1])\}$. Go to 1.1.
4. Set R to $(x, f(x), f'(x))$.
5. If the point $(L(1), L(2))$ is on the line $Li(R)$, terminate; the solution is $(L[1], L[2], R[3])$.
- 6.1 Set x to the abscissa of the intersection of the lines $Li(L)$, $Li(R)$.
- 6.2 If $R[3] \geq 0$, go to 1.
- 7.1 Set x to the larger of x and the abscissa of the intersection of the line $Li(0, 0, -m_2)$ with the line passing through $(L[1], L[2])$ and $(R[1], R[2])$. Go to 1.

Notes

- Step 3: x^* is determined by the calculation: if $LL[3] \geq L[3]$, set $x^* = \infty$. Otherwise

$$x^* = (f'(x) LL[1] - x LL[3]) / (f'(x) - LL[3]) .$$
- Step 5: The point $(L[1], L[2])$ is on $Li(R)$ when $L[2] = R[2] + R[3] (R[1] - L[1])$. In that case, $Li(L[1], L[2], R[3])$ is a support of f at L just as $Li(L)$ is, and since $R[3] \geq L[3]$, the former is preferable.
- Step 6.1: x is given by the formula

$$(L[2] - R[2] + R[1]R[3] - L[1]L[3]) / (R[3] - L[3]) .$$
- Step 7.1: The desired abscissa is

$$(L[1]R[2] - R[1]L[2]) / (R[2] - L[2] + m_2(R[1] - L[1])) .$$

THE METHOD OF PARAMETRIC DECOMPOSITION IN MATHEMATICAL PROGRAMMING: THE NONCONVEX CASE

Jacques Gauvin

1. INTRODUCTION

We consider a large mathematical program which may be written in the following general formulation:

$$(P_0) \left\{ \begin{array}{l} \max f(x,y,z) \\ \text{subject to } g_i(x,y) \leq 0 \quad , \quad i = 1, \dots, m \\ \qquad \qquad \qquad h_j(x,z) = 0 \quad , \quad j = 1, \dots, p \\ \qquad \qquad \qquad l_k(y,z) \leq 0 \quad , \quad k = 1, \dots, r \end{array} \right.$$

where $x \in R^n$, $y \in R^s$, $z \in R^t$. It is expected here that the objective function f and the constraints g_i, h_j are decomposable in their arguments x , y , and z and that the optimization is easy when y and z are held fixed.

Essentially in the method of parametric decomposition [3] or the method of primal decomposition by resource allocation [6], [15,16] the variables y and z are considered as parameters and are held fixed at the first optimization level where the following problem in x is solved by a standard NLP method:

$$(P_1) \left\{ \begin{array}{l} \max f(x;y,z) \\ \text{subject to } g_i(x,y) \leq 0 \quad , \quad i = 1, \dots, m \\ \qquad \qquad \qquad h_j(x,a) = 0 \quad , \quad j = 1, \dots, p \quad . \end{array} \right.$$

Let $S(y) = \{x \in R^n \mid g_i(x,y) \leq 0 \quad , \quad i = 1, \dots, m\}$

$T(z) = \{x \in R^n \mid h_j(x,z) = 0 \quad , \quad j = 1, \dots, p\}$

and $U(y,z) = S(y) \cap T(z)$ be the feasible set for (P_1) . The optimal value of (P_1) is given by

$$v(y,z) = \begin{cases} \max_{x \in U(y,z)} f(x,y,z) & \text{if } U(y,z) \neq \phi \\ -\infty & \text{if } U(y,z) = \phi \end{cases}$$

which is called the extremal-value function.

$$P(y,z) = \{\bar{x} \in U(y,z) \mid f(\bar{x},y,z) = v(y,z)\}$$

is the set of optimal solutions for (P_1) . $V = \{(y,z) \mid U(y,z) \neq \phi\}$ is the set of feasible parameters for (P_1) .

At the second level, a postoptimization is performed on the feasible parameters

$$(P_2) \left\{ \begin{array}{l} \max v(y,z) \\ \text{subject to } (y,z) \in V \\ l_k(y,z) \leq 0 \quad , \quad k = 1, \dots, r \end{array} \right.$$

In this method of decomposition-coordination by resource allocation, an optimal solution of subproblem (P_1) always gives at least a feasible solution to the original problem (P_0) . Such is not the case in a dual-type method such as that of decomposition-coordination by prices (see Lasdon [9]) where a feasible solution of the original problem will only be attained, in general, at an optimal solution.

Under convexity assumptions, the extremal-value function $v(y,z)$ is concave and subdifferentiable, and some methods have been proposed to solve the problem (see, for example, [15,6,8]. A description of the first two works can be found in [9] Chapter 9).

To our knowledge, it seems that recent methods of nondifferentiable convex optimization [17,10] have not been applied to the problem of parametric decomposition.

It is not the purpose of this paper to survey the results and methods of convex parametric decomposition, but rather to see what can be done in the nonconvex case. To make the presentation more simple, we will assume

$$f(x, y, z) = f(x)$$

$$g_i(x, y) = g_i(x) - y_i$$

$$h_j(x, y) = h_j(x) - z_j \quad .$$

The extension of results to the more general problem can be easily obtained.

First some results from [5] are presented from which a locally Lipschitz property for the extremal-value function is derived. Also an estimation of the generalized gradient of this function is obtained. Maybe these results can be useful for designing a method to solve the nonconvex and nondifferentiable second level optimization problem.

Some results are already available in that direction. F.H. Clarke [2] has given optimality conditions that can be applied to the postoptimization problem. Also some algorithms have been proposed for optimizing nondifferentiable nonconvex functions; see [12] for such a proposal and a review of others.

In the sequel all functions defining program (P_0) are assumed continuously differentiable.

2. A LOCALLY LIPSCHITZ PROPERTY FOR THE EXTREMAL-VALUE FUNCTION

Let \bar{y}, \bar{z} be some feasible parameters for problem (P_0) . For \bar{x} , a local maximum of (P_1) , let $I(\bar{x}; \bar{y}, \bar{z}) = \{i | g_i(\bar{x}) = \bar{y}_i\}$ be the set of active inequality constraints and $K(\bar{x}; \bar{y}, \bar{z})$ be the set of Kuhn-Tucker vectors corresponding to \bar{x} , that is the set of $(u, v) \in \mathbb{R}^m \times \mathbb{R}^p$ such that

$$\left\{ \begin{array}{l} \nabla f(\bar{x}) = \sum_{i=1}^m u_i \nabla g_i(\bar{x}) + \sum_{j=1}^p v_j \nabla h_j(\bar{x}) \\ u_i \geq 0 \quad i = 1, \dots, m \\ u_i (g_i(\bar{x}) - \bar{y}_i) = 0 \quad . \end{array} \right.$$

Let

$$K(\bar{y}, \bar{z}) = \bigcup_{\bar{x} \in P(\bar{y}, \bar{z})} K(\bar{x}; \bar{y}, \bar{z})$$

be the set of all multipliers corresponding to (\bar{y}, \bar{z}) .

The directional derivative of $v(y, z)$ at (\bar{y}, \bar{z}) , for a direction $d = (d^1, d^2) \in \mathbb{R}^m \times \mathbb{R}^p$, $\|d\| = 1$, is

$$v'(\bar{y}, \bar{z}; d) = \lim_{t \downarrow 0} \frac{v(\bar{y} + td^1, \bar{z} + td^2) - v(\bar{y}, \bar{z})}{t} .$$

We will also consider \liminf and \limsup for the right-hand side expression. Examples show that these limits can be infinite if for some $\bar{x} \in P(\bar{y}, \bar{z})$, $K(\bar{x}; \bar{y}, \bar{z})$ is empty or unbounded. To avoid this situation we assume at \bar{x} the Mangasarian-Fromowitz constraint qualification, denoted (CQ1).

$$(CQ1) \left\{ \begin{array}{l} (i) \quad \text{There exists a } w \in \mathbb{R}^n \text{ such that} \\ \quad \quad \quad \nabla g_i(\bar{x}) \cdot w < 0 \quad , \quad i \in I(\bar{x}; \bar{y}, \bar{z}) \\ \quad \quad \quad \nabla h_j(\bar{x}) \cdot w = 0 \quad , \quad j = 1, \dots, p \quad . \\ (ii) \quad \text{The gradients } \{\nabla h_j(\bar{x})\}, j = 1, \dots, p \\ \quad \quad \quad \text{are linearly independent.} \end{array} \right.$$

In the absence of equality constraints, (CQ1) is equivalent to the Cottle constraint qualification: the system

$$\sum_{i \in I(\bar{x}; \bar{y}, \bar{z})} u_i \nabla g_i(\bar{x}) = 0 \quad , \quad u_i \geq 0$$

has no nonzero solution. If the g_i are convex and the h_j affine, (CQ1) is the well-known Slater condition.

This regularity condition has the following nice property (see [4 or 5]).

Lemma 2.1

Let \bar{x} be a local maximum. Then $K(\bar{x}; \bar{y}, \bar{z})$ is a nonempty, convex and *compact* set if and only if (CQ1) is satisfied at \bar{x} .

The presence of equality constraints may cause the set $U(y, z)$ to be empty near (\bar{y}, \bar{z}) . The next lemma gives a condition to rule out this situation [5, Lemma 2.5].

Lemma 2.2

If (CQ1) is satisfied at some $\bar{x} \in P(\bar{y}, \bar{z})$ then $U(y, z)$ is not empty near (\bar{y}, \bar{z}) .

Conditions to have the function $v(y, z)$ continuous are given in the following [5, Theorem 2.6].

Theorem 2.3

If $U(\bar{y}, \bar{z})$ is nonempty and $U(y, z)$ is uniformly compact near (\bar{y}, \bar{z}) then $U(y, z)$ and $v(y, z)$ are upper semicontinuous at (\bar{y}, \bar{z}) . Furthermore if (CQ1) holds at some $\bar{x} \in P(\bar{y}, \bar{z})$ then $v(y, z)$ is also lower semicontinuous at (\bar{y}, \bar{z}) .

The (CQ1) regularity condition has the advantage of being preserved in a neighborhood of (\bar{y}, \bar{z}) [5, Corollary 2.10].

Theorem 2.4

If in Theorem 2.3, (CQ1) holds at every $\bar{x} \in P(\bar{y}, \bar{z})$, then there exists a $\delta > 0$, such that for all (y, z) satisfying

$\|(y, z) - (\bar{y}, \bar{z})\| \leq \delta$, (CQ1) holds also at each $x \in P(y, z)$, the point-to-set map $K(y, z)$ is upper semi-continuous at (\bar{y}, \bar{z}) and $K(y, z)$ is uniformly compact near (\bar{y}, \bar{z}) . More precisely, this result means that for any sequence $\{(y^n, z^n)\}$, $(y^n, z^n) \rightarrow (\bar{y}, \bar{z})$, there exist $(u^n, v^n) \in K(y^n, z^n)$, a subsequence $\{(u^m, v^m)\}$ and a $(\bar{u}, \bar{v}) \in K(\bar{y}, \bar{z})$ such that $(u^m, v^m) \rightarrow (\bar{u}, \bar{v})$.

It should be noted that $K(y, z)$ is not necessarily lower semicontinuous at (\bar{y}, \bar{z}) under the assumptions of Theorem 2.4.

The next result gives bounds for the potential directional derivatives of $v(y, z)$. It does not require any second-order assumptions [5, Theorem 3.3].

Theorem 2.5

Suppose that $U(\bar{y}, \bar{z})$ is nonempty, $u(y, z)$ is uniformly compact near (\bar{y}, \bar{z}) and (CQ1) holds at some $\bar{x} \in P(\bar{y}, \bar{z})$, then for any direction d

$$\begin{aligned} & \min_{(u, v) \in K(\bar{x}; \bar{y}, \bar{z})} \{u \cdot d^1 + v \cdot d^2\} \\ & \leq \liminf_{t \rightarrow 0} \frac{v(\bar{y} + td^1, \bar{z} + td^2) - v(\bar{y}, \bar{z})}{t} . \end{aligned}$$

Furthermore, if we assume that (CQ1) holds at every $\bar{x} \in P(\bar{y}, \bar{z})$ then

$$\begin{aligned} & \max_{\bar{x} \in P(\bar{y}, \bar{z})} \min_{(u, v) \in K(\bar{x}; \bar{y}, \bar{z})} \{u \cdot d^1 + v \cdot d^2\} \\ & \leq \liminf_{t \rightarrow 0} \frac{v(\bar{y} + td^1, \bar{z} + td^2) - v(\bar{y}, \bar{z})}{t} \\ & \leq \limsup_{t \rightarrow 0} \frac{v(\bar{y} + td^1, \bar{z} + td^2) - v(\bar{y}, \bar{z})}{t} \end{aligned}$$

$$\leq \max_{\bar{x} \in P(\bar{y}, \bar{z})} \max_{(u, v) \in K(\bar{x}; \bar{y}, \bar{z})} \{u \cdot d^1 + v \cdot d^2\} .$$

The bounds given in Theorem 2.5 are sharp: there are examples for which the directional derivatives of $v(y, z)$ exist at (\bar{y}, \bar{z}) with the upper bound attained for a direction d_1 , the lower bound attained for some other d_2 , and a value strictly in between for a different d_3 (see [5, example 3.1]).

If in Theorem 2.5, we replace (CQ1) by the following more restrictive regularity condition.

$$(CQ2) \left\{ \begin{array}{l} \text{The gradients } \{\nabla g_i(\bar{x}), \nabla h_j(\bar{x})\}, i \in I(\bar{x}; \bar{y}, \bar{z}), j = 1, \dots, p, \\ \text{are linearly independent.} \end{array} \right.$$

then the directional derivative exists and is given by

$$v'(\bar{x}, \bar{y}; d) = \max_{\bar{x} \in P(\bar{y}, \bar{z})} \{\bar{u} \cdot d^1 + \bar{v} \cdot d^2\} \quad (2.6)$$

where (\bar{u}, \bar{v}) is the unique multiplier vector corresponding to \bar{x} .

Under convexity assumptions, we can obtain the following corollary of Theorem 2.5 [5, Corollary 3.5]; [7].

Corollary 2.7

Suppose the functions $f, \{g_i\}, i = 1, \dots, m$ are convex and $\{h_j\}$ are affine. If $U(\bar{y}, \bar{z})$ is nonempty, $U(y, z)$ is uniformly compact near (\bar{y}, \bar{z}) and (CQ1) (which is then equivalent to the Slater condition) is satisfied for each $\bar{x} \in P(\bar{y}, \bar{z})$, then $v(y, z)$ has a directional derivative for any direction d at (\bar{y}, \bar{z}) and

$$v'(\bar{y}, \bar{z}; d) = \max_{\bar{x} \in P(\bar{y}, \bar{z})} \min_{(u, v) \in K(\bar{x}; \bar{y}, \bar{z})} \{u \cdot d^1 + v \cdot d^2\} .$$

From Theorem 2.4 and Theorem 2.5, it is possible to obtain a locally Lipschitz property for the extremal-value function.

Theorem 2.8

Suppose that $U(\bar{y}, \bar{z})$ is nonempty, $U(y, z)$ is uniformly compact near (\bar{y}, \bar{z}) and (CQ1) holds at every $\bar{x} \in P(\bar{y}, \bar{z})$; then there exists a δ -neighborhood $N_\delta(\bar{y}, \bar{z})$ and a finite K such that for any $(y_1, z_1), (y_2, z_2) \in N_\delta(\bar{y}, \bar{z})$

$$|v(y_2, z_2) - v(y_1, z_1)| \leq K \|(y_2, z_2) - (y_1, z_1)\| .$$

Proof

From Theorem 2.4, the regularity condition (CQ1) remains valid at every $(y, z) \in N_\varepsilon(\bar{y}, \bar{z})$, for some $\varepsilon > 0$, with the set of multipliers $K(y, z)$ uniformly bounded. Therefore, for some δ , $0 < \delta < \varepsilon$, Theorem 2.5 is valid at any $(y, z) \in N_\delta(\bar{y}, \bar{z})$, and, for any direction $d = (d^1, d^2)$, $\|d\| = 1$, there exists some finite K_1 and K_2 such that

$$\begin{aligned} K_1 &\leq \liminf_{t \rightarrow 0} [v(y+td^1, z+td^2) - v(y, z)]/t \\ &\leq \limsup_{t \rightarrow 0} [v(y+td^1, z+td^2) - v(y, z)]/t \leq K_2 \end{aligned} \quad (2.9)$$

For $(y^2, z^2), (y^1, z^1) \in N_\delta(\bar{y}, \bar{z})$, write $(y^2, z^2) - (y^1, z^1) = \lambda(d^1, d^2)$ where $\|(d^1, d^2)\| = 1$ and $\lambda = \|(y^2, z^2) - (y^1, z^1)\|$. From (2.9) it follows that the function $v(y^1+td^1, z^1+td^2)$ is Lipschitz continuous on the ray d , for $t \in [0, \lambda]$, hence absolutely continuous on that ray. Therefore its derivative exists almost everywhere and we can write

$$v(y^1+\lambda d^1, z^1+\lambda d^2) - v(y^1, z^1) = \int_0^\lambda \frac{d}{dt} v(y^1+td^1, z^1+td^2) dt .$$

This derivative, when it exists, is equal to the right derivative, which is the directional derivative $v'(y^1+td^1, z^1+td^2; d)$. From (2.9) we then obtain

$$K_1 \lambda \leq v(y^2, z^2) - v(y^1, z^1) \leq K_2 \lambda$$

and the result follows with $K = \max\{|K_1|, |K_2|\} \cdot \|\cdot\|$

Stern and Topkis [14] have obtained the previous result for a program without equality constraints

$$\begin{cases} \max f(x) \\ \text{subject to } g_i(x) \leq y_i, & i = 1, \dots, m \end{cases}$$

assuming that the functions $g_i(x)$ have continuous second derivatives.

F.H. Clarke [2] also considers the problem

$$(P_s) \begin{cases} \min g_0(x) \\ \text{subject to } g_i(x) \leq s_i, & i = 1, \dots, m \end{cases}$$

where he defines (P_s) to be "calm" if the extremal-value function

$$\phi(s) = \inf\{g_0(x) \mid g_i(x) \leq s_i, \quad i = 1, \dots, m\}$$

is finite, and if

$$\liminf_{s' \rightarrow s} \frac{[\phi(s') - \phi(s)]}{\|s' - s\|} > -\infty$$

The program (P_s) is also defined to be "normal" if the Kuhn-Tucker conditions hold at every optimal point \bar{x} of (P_s) . In fact, Clarke does not assume that the functions $g_i(x)$ are differentiable and he gives some generalized Kuhn-Tucker conditions which reduce to the usual ones when differentiability is assumed. Then he shows that if (P_s) is "calm" then (P_s) is normal. The converse of this result is not valid unless (P_s) is assumed "normal" with bounded multipliers. He also shows that on a neighborhood S of

0 where $\phi(s)$ is finite, $P(s)$ is "calm" and "normal" for almost all s in S . The following example shows that the "almost all s " is meaningful.

Example 2.10

$$(P_s) \quad \begin{cases} \min g_0(x) = -x \\ \text{subject to } g_1(x) = x^3 \leq s \end{cases}$$

Then $\phi(s) = -s^{1/3}$ and (P_s) is either "normal" or "calm" at $s = 0$.

3. THE GENERALIZED GRADIENT OF THE EXTREMAL-VALUE FUNCTION

For the previous section, we have conditions to guarantee the extremal-value function $v(y,z)$ to be Lipschitz in some ball about (\bar{y}, \bar{z}) . Following Clarke [1] the gradient $\nabla v(y,z)$ then exists almost everywhere in that ball (Rademacher's theorem).

Definition 3.1 [1]

The *generalized gradient* of $v(y,z)$ at (\bar{y}, \bar{z}) denoted $\partial v(\bar{y}, \bar{z})$, is the convex hull of the set of limits of the form $\lim \nabla v(\bar{y}+h_n, \bar{z}+k_n)$, where $(h_n, k_n) \rightarrow (0,0)$.

$\partial v(\bar{y}, \bar{z})$ is a nonempty convex compact set. The *generalized directional derivative* of $v(y,z)$ at (\bar{y}, \bar{z}) for the direction $d = (d^1, d^2) \in \mathbb{R}^m \times \mathbb{R}^p$ is

$$v^0(\bar{y}, \bar{z}; d) = \limsup_{\substack{(y,z) \rightarrow (\bar{y}, \bar{z}) \\ t \rightarrow 0}} [v(y+td^1, z+td^2) - v(y,z)]/t .$$

Then

$$v^0(\bar{y}, \bar{z}; d) = \max\{g^1 \cdot d^1 + g^2 \cdot d^2 \mid (g^1, g^2) \in \partial v(\bar{y}, \bar{z})\} ,$$

that is, $v^0(\bar{y}, \bar{z}; \cdot)$ is the *support function* of $\partial v(\bar{y}, \bar{z})$.

We can obtain, under the assumptions of Theorem 2.8, the following estimation for the generalized gradient of the extremal-value function.

Theorem 3.2

Suppose that $U(\bar{y}, \bar{z})$ is nonempty, $U(y, z)$ is uniformly compact near (\bar{y}, \bar{z}) and (CQ1) holds at every optimal point $\bar{x} \in P(\bar{y}, \bar{z})$. Then the generalized gradient of $v(y, z)$ at (\bar{y}, \bar{z}) is contained in the convex hull of all the Kuhn-Tucker multipliers corresponding to the optimal points; that is $\partial v(\bar{y}, \bar{z}) \subseteq \text{co } K(\bar{y}, \bar{z})$.

Proof

Take a sequence $\{(y^n, z^n)\}$, $(y^n, z^n) \rightarrow (\bar{y}, \bar{z})$ where $\nabla v(y^n, z^n)$ exists. For any direction $d = (d^1, d^2) \in \mathbb{R}^m \times \mathbb{R}^p$ we have, by Theorem 2.4, that (CQ1) still holds in some neighborhood of (\bar{y}, \bar{z}) , and we have, by Theorem 2.5,

$$\begin{aligned} \nabla v(y^n, z^n) \cdot d &= \lim_{t \rightarrow 0} [v(y^n + td^1, z^n + td^2) - v(y^n, z^n)] / t \\ &\leq \max_{(u, v) \in K(y^n, z^n)} [u \cdot d^1 + v \cdot d^2] \\ &= [u^n \cdot d^1 + v^n \cdot d^2] \end{aligned}$$

for some $(u^n, v^n) \in K(y^n, z^n)$. Again from Theorem 2.4, there exists a subsequence $\{(u^m, v^m)\}$, a $(\bar{u}, \bar{v}) \in K(\bar{y}, \bar{z})$ such that $(u^m, v^m) \rightarrow (\bar{u}, \bar{v})$. Taking the limit on both sides we obtain by [13, Theorem 32.2],

$$\begin{aligned} [\lim \nabla v(y^n, z^n)] \cdot d &\leq \bar{u} \cdot d^1 + \bar{v} \cdot d^2 \\ &\leq \max_{(u, v) \in K(\bar{y}, \bar{z})} [u \cdot d^1 + v \cdot d^2] \\ &= \max_{(u, v) \in \text{co } K(\bar{y}, \bar{z})} [u \cdot d^1 + v \cdot d^2] \end{aligned}$$

Since this result holds for every direction d , we have, by [13, Theorem 13.1],

$$\lim \nabla v(y^n, z^n) \in \text{co } K(\bar{y}, \bar{z}) \quad .$$

This with Definition 3.1 gives the result. \parallel

If the directional derivative $v'(\bar{y}, \bar{z}; d)$ exists and is equal to $v^0(y, z; d)$ for every direction d , then $v(y, z)$ is said to be *quasidifferentiable* at (\bar{y}, \bar{z}) (see [11]).

Corollary 3.3

If in Theorem 3.2 the regularity condition (CQ1) is replaced by (CQ2) (see (2.6)), then $\partial v(\bar{y}, \bar{z}) = \text{co } K(\bar{y}, \bar{z})$ and $v(y, z)$ is quasidifferentiable at (\bar{y}, \bar{z}) .

Proof

For any $(u, v) \in K(\bar{y}, \bar{z})$, we have

$$\begin{aligned} u \cdot d^1 + v \cdot d^2 &\leq \max_{(u, v) \in K(\bar{y}, \bar{z})} [u \cdot d^1 + v \cdot d^2] \\ &= v'(\bar{y}, \bar{z}; d) \quad \text{by (2.6)} \\ &\leq v^0(\bar{y}, \bar{z}; d) \quad ; \end{aligned}$$

hence, by [13, Theorem 13.1], $(u, v) \in \partial v(\bar{y}, \bar{z})$ and $K(\bar{y}, \bar{z}) \subseteq \partial v(\bar{y}, \bar{z})$.

By Theorem 3.2, $\partial v(\bar{y}, \bar{z}) = \text{co } K(\bar{y}, \bar{z})$. Therefore we have, for every direction d ,

$$\begin{aligned} v^0(\bar{y}, \bar{z}; d) &= \max_{(u, v) \in \text{co } K(\bar{y}, \bar{z})} [u \cdot d^1 + v \cdot d^2] \\ &= \max_{(u, v) \in K(\bar{y}, \bar{z})} [u \cdot d^1 + v \cdot d^2] = v'(\bar{y}, \bar{z}; d) \end{aligned}$$

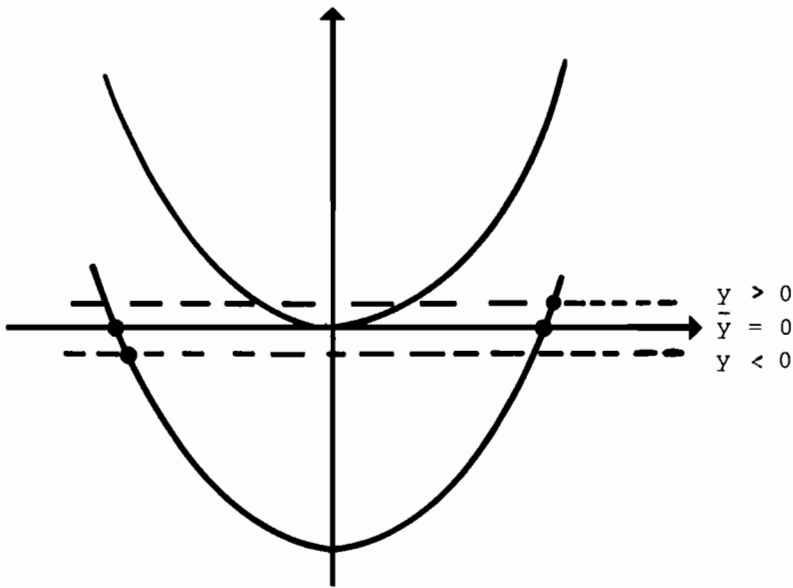
by [13, Theorem 32.2] and (2.6). \parallel

The next example illustrates the previous result.

Example 3.4

$$\max f = \begin{cases} x^3 & \text{if } x \geq 0 \\ x^2 & \text{if } x \leq 0 \end{cases}$$

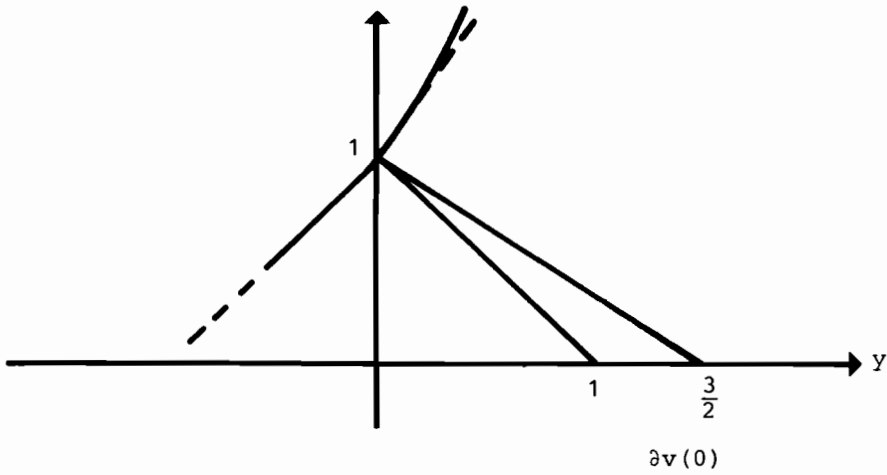
$$\text{subject to } g = x^2 - 1 \leq y$$



At $\bar{y} = 0$, the maximum occurs at $\bar{x}_1 = 1$ with multiplier $\bar{u}_1 = 3/2$ and at $\bar{x}_2 = -1$ with multiplier $\bar{u}_2 = 1$ where (CQ2) is satisfied at both points.

For $y > 0$, the maximizer is $+\sqrt{1+y}$, and for $y < 0$, the maximizer is $-\sqrt{1+y}$; hence the extremal-value function is

$$v(y) = \begin{cases} (1+y)^{3/2} & \text{if } y > 0 \\ (1+y) & \text{if } y < 0 \end{cases}$$



Therefore

$$v'(0; 1) = \frac{3}{2} = v^0(0; 1)$$

$$v'(0; -1) = -1 = v^0(0; -1)$$

and

$$\partial v(y) = \begin{cases} \frac{3}{2} & \text{if } y > 0 \\ [1, \frac{3}{2}] & \text{if } y = 0 \\ 1 & \text{if } y < 0 \end{cases}$$

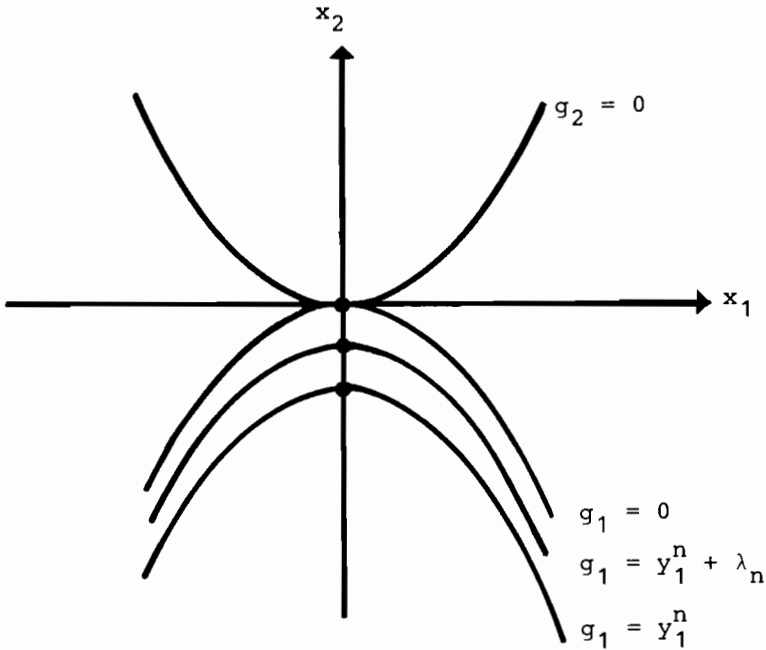
The next example shows that under assumption (CQ1) the extremal-value function is not necessary quasidifferentiable.

Example 3.5

$$\max f = x_2$$

$$\text{subject to } g_1 = x_2 + x_1^2 \leq y_1$$

$$g_2 = x_2 - x_1^2 \leq y_2$$



At $\bar{y} = (0,0)$, the maximum occurs at $\bar{x} = (0,0)$ where (CQ1) holds with the set of multipliers

$$K(\bar{x};0) = \{(u_1, u_2) \mid u_1 + u_2 = 1, u_1 \geq 0, u_2 \geq 0\} .$$

For the direction $d = (1,0)$, $v'(0;d) = \frac{1}{2}$. Now take a sequence $\{y^n, \lambda_n\}$ where $y^n = (y_1^n, 0)$, $y_1^n < 0$, $y_1^n \rightarrow 0$, $\lambda_n \neq 0$, $\lambda_n < |y_1^n|$; then

$$[v(y^n + \lambda_n d) - v(y^n)]/\lambda_n = [y_1^n + \lambda_n - y_1^n]/\lambda_n = 1 \leq v^0(0;d) ,$$

and therefore $v^0(0;d) \neq v'(0;d)$ and $v(y)$ is not quasidifferentiable at $\bar{y} = 0$.

The next results characterize the gradients of $v(y,z)$ when they are assumed to exist.

Corollary 3.6

If the assumptions of Theorem 3.2 are satisfied and if the gradient

$$\nabla v(\bar{y}, \bar{z}) = \left[\frac{\partial v}{\partial \bar{y}_i}, \frac{\partial v}{\partial \bar{z}_j} \right], \quad i = 1, \dots, m; \quad j = 1, \dots, p,$$

exists, then

$$\begin{aligned} \max_{\bar{x} \in P(\bar{y}, \bar{z})} \min_{(u, v) \in K(\bar{x}; \bar{y}, \bar{z})} \begin{bmatrix} u_i \\ v_j \end{bmatrix} &\leq \begin{bmatrix} \frac{\partial v}{\partial \bar{y}_i} \\ \frac{\partial v}{\partial \bar{z}_j} \end{bmatrix} \\ &\leq \min_{\bar{x} \in P(\bar{y}, \bar{z})} \max_{(u, v) \in K(\bar{x}; \bar{y}, \bar{z})} \begin{bmatrix} u_i \\ v_j \end{bmatrix} \\ & \quad i = 1, \dots, m; \quad j = 1, \dots, p. \end{aligned}$$

Proof

For any direction d we have, by Theorem 2.5,

$$\begin{aligned} \nabla v(\bar{y}, \bar{z}) \cdot d &= v'(\bar{y}, \bar{z}; d) \\ &\geq \min_{\bar{x} \in P(\bar{y}, \bar{z})} \max_{(u, v) \in K(\bar{x}; \bar{y}, \bar{z})} [u \cdot d^1 + v \cdot d^2]. \end{aligned}$$

Hence for the direction $-d$,

$$\nabla v(\bar{y}, \bar{z}) \cdot (-d) \geq \max \min - [u \cdot d^1 + v \cdot d^2],$$

which implies

$$\nabla v(\bar{y}, \bar{z}) \cdot d \leq \min_{\bar{x} \in P(\bar{y}, \bar{z})} \max_{(u, v) \in K(\bar{x}; \bar{y}, \bar{z})} [u \cdot d^1 + v \cdot d^2].$$

These two inequalities taken with the directions d equal to the unit vectors in $R^m \times R^p$ give the results. ||

Corollary 3.7

If in Corollary 3.6 the regularity condition (CQ1) is replaced by (CQ2), then the set of multipliers is a singleton, that is $K(\bar{y}, \bar{z}) = \{(\bar{u}, \bar{v})\}$ and $\forall v(\bar{y}, \bar{z}) = (\bar{u}, \bar{v})$.

Proof

Under (CQ2), for each $\bar{x} \in P(\bar{y}, \bar{z})$, $K(\bar{x}; \bar{y}, \bar{z})$ is a singleton. From Corollary 3.6, we have

$$\max_{\bar{x} \in P(\bar{y}, \bar{z})} u_i \leq \frac{\partial v}{\partial y_i} \leq \min_{\bar{x} \in P(\bar{y}, \bar{z})} u_i, \quad i = 1, \dots, m,$$

which implies the result. (The same holds for $\frac{\partial v}{\partial z_j}$, $j = 1, \dots, p$.)

Recently R. Mifflin [11,12] has introduced the notion of semismooth and weakly upper semismooth functions. The definitions are the following.

Definition 3.8

A function $F: \mathbb{R}^n \rightarrow \mathbb{R}$ is *semismooth* at $\bar{x} \in \mathbb{R}^n$ if

- (i) F is Lipschitz on a ball about \bar{x} and
- (ii) for each direction $d \in \mathbb{R}^n$ and for any sequences $\{t_k\} \subset \mathbb{R}_+$, $\{\theta_k\} \subset \mathbb{R}^n$ and $\{g_k\} \subset \mathbb{R}^n$ such that $\{t_k\} \downarrow 0$, $\{\theta_k/t_k\} \rightarrow 0$ and $g_k \in \partial F(\bar{x} + t_k d + \theta_k)$ the sequence $\{g_k \cdot d\}$ has exactly one accumulation point.

Definition 3.9

$F: \mathbb{R}^n \rightarrow \mathbb{R}$ is *weakly upper semismooth* at \bar{x} if the following conditions hold:

- (i) F is Lipschitz on a ball about \bar{x}
- (ii) for each direction $d \in \mathbb{R}^n$

$$\liminf_{k \rightarrow \infty} g_k \cdot d \geq \limsup_{t \downarrow 0} [F(\bar{x} + td) - F(\bar{x})]/t$$

where $\{g_k\} \subset R^n$ is any sequence such that $g_k \in \partial F(\bar{x} + t_k d)$
and $\{t_k\} \subset R_+$ is any sequence such that $\{t_k\} \downarrow 0$.

In Example 3.5 it can be easily seen that the extremal-value function is weakly upper semismooth and even semismooth. At the present time, it is not known whether these properties are generally satisfied or not by the extremal-value function under the assumption of Theorem 3.2.

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A SET OF NONSMOOTH OPTIMIZATION TEST PROBLEMS

Test Problem 1: MAXQUAD

The objective function to be minimized is

$$f(x) = \max \{ (A_k x, x) - (b_k, x) \mid k = 1, \dots, 5 \}$$

where

$$x \in R^{10} .$$

The (symmetric, diagonal dominating) matrices A_k , together with the vectors b_k , are given by the fancy formulae:

$$A_k(i, j) = e^{i/j} \cos(i, j) \sin(k) \quad i < j$$

$$b_k(i) = e^{i/k} \sin(i.k) ,$$

and the diagonals of A_k are

$$A_k(i, i) = \frac{i}{10} \sin(k) + \sum_{j \neq i} A_k(i, j) .$$

The following naive FORTRAN program fills up the data for $K \text{ MAX} = 5, N = 10$.

```

do 30 k=1,kmax
ak=float(k)
do 10 i=1,n
ai=float(i)
do 10 j=i,n
aj=float(j)
a(k,i,j)=exp(ai/aj)*cos(ai*aj)*sin(ak)
a(k,j,i)=a(k,i,j)
10 continue
do 20 i=1,n
ai=float(i)
f(k,i)=exp(ai/ak)*sin(ai*ak)
a(k,i,i)=abs(sin(ak))*float(i)/float(n)
do 20 j=1,n
if(j.ne.i) a(k,i,i)=a(k,i,i)+abs(a(k,i,j))
20 continue
30 continue

```

For this particular problem, the optimal solution is

$$f^* = -0.8414$$

$$x^* = (-0.1263, -0.0346, -0.0067, 0.2668, 0.0673, \\ 0.2786, 0.0744, 0.1387, 0.0839, 0.0385) .$$

The following FORTRAN program computes the value of the objective, VAL, and a subgradient, G, at a given point, x.

```

k0=0
do 50 k=1,kmax
fi=0.
do 10 i=1,n
z=x(i)
fi=fi-f(k,i)*z
do 10 j=1,n
10 fi=fi+a(k,i,j)*x(j)*z
if(k0.eq.0) go to 20
if(fi.le.val) go to 50
20 k0=k
val=fi
50 continue
do 70 i=1,n
z=-f(k0,i)
do 60 j=1,n
60 z=z+2.*a(k0,i,j)*x(j)
70 g(i)=z

```

The standard starting point is $x_i = 1$ $i = 1, \dots, n$, for which $f = 5337$. Another interesting starting point is $x = 0$, since f has a kink at this point ($f_k(0) = 0$ $k = 1, \dots, k \text{ max}$).

Test Problem 2: SHELL DUAL

This is the second problem of Colville, where all the constraints are penalized with an ℓ_1 penalty function (A.R. Colville: *A Comparative Study of Nonlinear Programming Codes*, IBM, New York, Scientific Report 320.2949 (1968)).

There are two groups of variables: x_i , $i = 1, \dots, k=10$ and y_j , $j = 1, \dots, m=5$.

The original problem from Colville is

$$\min 2 \sum_{j=1}^m d_j y_j^3 + (Cy, y) - (b, x)$$

$$(Ax)_j - 2(Cy)_j \leq e_j + 3d_j y_j^2 \quad j = 1, \dots, m$$

$$x \geq 0, \quad y \geq 0.$$

Here we define the functions

$$P_j(x, y) = (Ax)_j - 2(Cy)_j - 3d_j y_j^2 - e_j$$

$$Q(x, y) = \sum_{i=1}^k \min(0, x_i) + \sum_{j=1}^m \min(0, y_j),$$

and the function to be minimized is

$$f(x, y) = 2 \left| \sum_{j=1}^m d_j y_j^3 \right| + (Cy, y) - (b, x) \\ + 100 \left\{ \sum_{j=1}^m \max(0, P_j(x, y)) - Q(x, y) \right\}.$$

The vectors d , e , b and matrices A and C are given at the end of this problem.

A program which computes the function $f(x,y)$ and its gradient is also given. In this program, the data m , k , d , b , e , A , C , together with $\text{PENAL} = 100$ are supposed to be passed in a `COMMON` block.

The variables x and y form a single vector X where

$$X(I) = y_i \quad i = 1, \dots, m$$

$$X(M+I) = x_i \quad i = 1, \dots, k \quad .$$

The optimal point is $f^* = 32.3488$

$$y^* = (0.3, 0.3335, 0.4, 0.4283, 0.224)$$

$$x^* = (0., 0., 5.1741, 0., 3.0611, \\ 11.8396, 0., 0., 0.1039, 0.) \quad .$$

The starting point is $f = 2400$.

$$y_j = 0.0001 \quad j = 1, \dots, 5$$

$$x_i = 0.0001 \quad i \neq 7$$

$$x_7 = 60 \quad .$$

Data for SHELL DUAL

d = 4. 8. 10. 6. 2.

e = -15. -27. -36. -18. -12.

Matrix a					b
-16.	2.	0.	1.	0.	-40.
0.	-2.	0.	4.	2.	-2.
-3.5	0.	2.	0.	0.	-0.25
0.	-2.	0.	-4.	-1.	-4.
0.	-9.	-2.	1.	-2.8	-4.
2.	0.	-4.	0.	0.	-1.
-1.	-1.	-1.	-1.	-1.	-40.
-1.	-2.	-3.	-2.	-1.	-60.
1.	2.	3.	4.	5.	5.
1.	1.	1.	1.	1.	1.

Symmetric Matrix c

30.	-20.	-10.	32.	-10.
-20.	39.	-6.	-31.	32.
-10.	-6.	10.	-6.	-10.
32.	-31.	-6.	39.	-20.
-10.	32.	-10.	-20.	30.

Program for computing function and gradient in SHELL DUAL

```

      z=0.
      do 10 j=1,m
10    z=z+d(j)*x(j)*x(j)*x(j)
      if (z.lt.0.) go to 30
      val=2.*z
      do 20 j=1,m
20    g(j)=6.*d(j)*x(j)*x(j)
      go to 50
30    val=-2.*z
      do 40 j=1,m
40    g(j)=-6.*d(j)*x(j)*x(j)
c
50    do 70 j=1,m
      z=0.
      do 60 i=1,m
60    z=z+c(i,j)*x(i)
      val=val+z*x(j)
70    g(j)=g(j)+2.*z
c
      do 80 i=1,k
      il=m+i
      val=val-b(i)*x(il)
80    g(il)=-b(i)
c
c          compute the constraints
c
      do 200 j=1,m
      z=-3.*d(j)*x(j)*x(j)-e(j)
      do 120 i=1,k
      il=m+i
120   z=z+a(i,j)*x(il)
      do 140 i=1,m
140   z=z-2.*c(i,j)*x(i)
      if(z.le.0.) go to 200
      val=val+penal*z
      g(j)=g(j)-6.*penal*d(j)*x(j)
      do 160 i=1,k
      il=m+i
160   g(il)=g(il)+penal*a(i,j)
      do 180 i=1,m
180   g(i)=g(i)-2.*penal*c(i,j)
200   continue
c
c          now the nonnegativity constraints
c
      do 320 i=1,n
      if (x(i).ge.0.) go to 320
      val=val-penal*x(i)
      g(i)=g(i)-penal
320   continue

```

Test Problem 3: EQUIL

The following nonconvex problem is a min-max formulation of an economic equilibrium problem in Scarf (The Approximation of Fixed Points of a Continuous Mapping, *SIAM Journal on Applied Mathematics*, 15 (1967), 1328-1343).

$$\text{Minimize } \{ \max [f_i(x) : i = 1, 2, \dots, N] : \sum_{j=1}^N x_j = 1, \\ x_j \geq 0, j = 1, 2, \dots, N \}$$

where

$$x \in \mathbb{R}^N$$

$$f_i(x) = \sum_{\ell=1}^{NA} f_i^\ell(x) \quad \text{for } i = 1, 2, \dots, N,$$

$$f_i^\ell(x) = (A_{\ell i} \sum_{k=1}^N W_{\ell k} x_k) / (x_i^\ell \sum_{k=1}^N A_{\ell k} x_k^{1-B_\ell}) - W_{\ell i}$$

$$\text{for } \ell = 1, 2, \dots, NA.$$

The input data N , NA , $W_{\ell k}$, $A_{\ell k}$ and B are given below.

Data for EQUIL

N = 8 NA = 5

Matrix W

3.	1.	.1	.1	5.	.1	.1	6.
.1	10.	.1	.1	5.	.1	.1	.1
.1	9.	10.	.1	4.	.1	7.	.1
.1	.1	.1	10.	.1	3.	.1	.1
.1	.1	.1	.1	.1	.1	.1	11.

Matrix A

1.	1.	1.	1.	1.	1.	1.	1.
2.	.8	1.	.5	1.	1.	1.	1.
1.	1.2	.8	1.2	1.6	2.	.6	.1
2.	.1	.6	2.	1.	1.	1.	2.
1.2	1.2	.8	1.	1.2	.1	3.	4.

Vector b = [.5 1.2 .0 2. 1.5]

An interesting starting point is

$$x_j = .125 \text{ for all } j$$

where $f(x) = \max_i f_i(x)$ has the value 9.7878. The optimal value of f is zero and occurs at a strictly positive point near

$$x = (.27, .03, .06, .09, .07, .31, .10, .07) .$$

The following FORTRAN SUBROUTINE CALCUL (x,G,VAL) requires the data NA, N, A, W and B to be available in the common blocks labelled ENTIER and REEL. It is useful when

$$(C) \quad \sum_{j=1}^N x_j = 1 .$$

If $x_j > 0$ for all j it returns the value of f in VAL and returns a projected (with respect to (C)) generalized gradient in G. Otherwise, it returns an arbitrarily large number in VAL and no useful G.

```

SUBROUTINE CALCUL(X,G,VAL)
COMMON /ENTIER/ NA,N
COMMON /REEL/ A(5,10),W(5,10),B(5)
COMMON /EXCESS/ ED(10),R(5,10),D(5,10),XB(5,10)
DIMENSION X(10),G(10)
VAL=1.E20
DO 10 I=1,N
ED(I)=0.
G(I)=0.
10 CONTINUE
DO 100 L=1,NA
SUMR=0.
SUMD=0.
DO 20 K=1,N
IF(X(K).LE.0.)GO TO 400
XB(L,K)=X(K)**B(L)
SUMD=SUMD+A(L,K)*X(K)/XB(L,K)
20 SUMR=SUMR+W(L,K)*X(K)
DO 80 I=1,N
D(L,I)=XB(L,I)*SUMD
R(L,I)=A(L,I)*SUMR/D(L,I)
80 ED(I)=ED(I)+R(L,I)-W(L,I)
100 CONTINUE
IMAX=1
VAL=ED(1)
DO 200 I=2,N
IF(ED(I).LE.VAL) GO TO 200
VAL=ED(I)
IMAX=I
200 CONTINUE
SUM=0.
DO 300 J=1,N
DO 240 L=1,NA
TLJ=A(L,J)*(1.-P(L))
IF(J.EQ.IMAX) GO TO 230
T=TLJ*XP(L,IMAX)/XB(L,J)
GO TO 240
230 T=TLJ+B(L)*D(L,J)/X(J)
240 G(J)=G(J)+(A(L,IMAX)*W(L,J)-T*R(L,IMAX))/D(L,IMAX)
SUM=SUM+G(J)
300 CONTINUE
SUM=SUM/FLOAT(N)
DO 350 J=1,N
350 G(J)=G(J)-SUM
400 RETURN
END

```

Test Problem 4: TR48

This problem is the dual of a transportation problem with 48 sources and 48 destinations. It was communicated to us by J.L. Goffin and the cost data is from M. Held, R.M. Karp: A Dynamic Programming Approach to Sequencing Problems, *J. Soc. Indust. Appl. Math.*, 10, 1 (1962), 196-210.

The objective to be minimized is

$$f(x) = -\left\{ \sum_{i=1}^n s_i x_i + \sum_{j=1}^n d_j \min_{i=1}^n (a_{ij} - x_i) \right\}$$

where $n = 48$. See the following pages for the statements for computing the function and its gradient, for the data, and for the initial and optimal points.

The initialization $x = 0$ is standard. The point "initial 2" is given by J.L. Goffin. It has been computed by subgradient optimization.

Statements for computing function and gradient in TR48

```

f=0.
do 10 i=1,n
g(i)=s(i)
10 f=f+s(i)*x(i)
do 50 j=1,n
xmax=1.e30
do 40 i=1,n
z=a(i,j)-x(i)
if(z.gt.xmax) go to 40
xmax=z
k=i
40 continue
g(k)=g(k)-d(j)
f=f+d(j)*xmax
50 continue
f=-f
do 70 i=1,n
70 g(i)=-g(i)

```

Sources and destinations

vector s	22, 53, 64, 15, 66, 37, 16, 23, 67, 18, 52, 69, 17, 29, 50, 13, 95, 34, 59, 36, 22, 94, 28, 34, 36, 38, 55, 77, 45, 34, 32, 58, 30, 88, 74, 59, 93, 54, 89, 30, 79, 46, 35, 41, 99, 52, 76, 93.
vector d	61, 67, 24, 84, 13, 86, 89, 46, 48, 50, 74, 75, 88, 40, 29, 45, 32, 21, 61, 21, 51, 14, 89, 79, 38, 20, 97, 19, 10, 73, 59, 92, 52, 66, 89, 65, 63, 47, 7, 61, 87, 19, 36, 43, 9, 12, 8, 67.

Statements for reading the symmetric cost matrix A

```

1 format((16f5.0))
read 1, ((a(i-1,j)),j=i,n),i=2,n)
do 10 i=1,n
a(i,i)=100000.
do 10 j=1,i
a(i,j)=a(j,i)
10 continue

```

(data for A on next page)

Initializations and Optimal Point for TR48

	initial 1	initial 2	optimal
f(x)	-464816.	-638524.94	-638565.
1	0.	11.19	144.
2	0.	127.20	257.
3	0.	-129.70	0.
4	0.	344.50	483.
5	0.	-40.72	89.
6	0.	-295.30	-165.
7	0.	-202.30	-72.
8	0.	-382.30	-252.
9	0.	-217.70	-88.
10	0.	-307.70	-178.
11	0.	178.10	311.
12	0.	-4.36	126.
13	0.	-123.30	7.
14	0.	-265.30	-135.
15	0.	28.28	158.
16	0.	70.57	209.
17	0.	-31.81	101.
18	0.	-222.30	-92.
19	0.	96.19	229.
20	0.	-52.79	80.
21	0.	-34.71	95.
22	0.	-59.16	71.
23	0.	-373.70	-244.
24	0.	-28.35	102.
25	0.	-141.70	-12.
26	0.	2.28	132.
27	0.	198.50	337.
28	0.	-69.16	61.
29	0.	-26.35	104.
30	0.	-88.72	41.
31	0.	130.80	261.
32	0.	-12.35	118.
33	0.	-30.70	99.
34	0.	-376.30	-246.
35	0.	23.18	156.
36	0.	-400.30	-270.
37	0.	197.10	330.
38	0.	-260.30	-130.
39	0.	813.50	952.
40	0.	-191.70	-62.
41	0.	31.29	161.
42	0.	345.50	484.
43	0.	-7.72	122.
44	0.	335.50	474.
45	0.	947.50	1086.
46	0.	722.50	861.
47	0.	-300.30	-170.
48	0.	73.20	206.

This problem seems very difficult. Another problem, called A48, and simpler, consists in defining

$$s_i = d_i = 1 \quad i = 1, \dots, n \quad .$$

The optimal value for this problem is -9870.

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NONSMOOTH OPTIMIZATION

Proceedings of a IIASA Workshop,

March 28 - April 8, 1977

Claude Lemarechal and Robert Mifflin, Editors

Optimization, a central methodological tool of systems analysis, is used in many of IIASA's research areas, including the Energy Systems and Food and Agriculture Programs. IIASA's activity in the field of optimization is strongly connected with nonsmooth or nondifferentiable extreme problems, which consist of searching for conditional or unconditional minima of functions that, due to their complicated internal structure, have no continuous derivatives. Particularly significant for these kinds of extreme problems in systems analysis is the strong link between nonsmooth or nondifferentiable optimization and the decomposition approach to large-scale programming.

This volume contains the report of the IIASA workshop held from March 28 to April 8, 1977, entitled Nondifferentiable Optimization. However, the title was changed to Nonsmooth Optimization for publication of this volume as we are concerned not only with optimization without derivatives, but also with problems having functions for which gradients exist almost everywhere but are not continuous, so that the usual gradient-based methods fail.

Because of the small number of participants and the unusual length of the workshop, a substantial exchange of information was possible. As a result, details of the main developments in nonsmooth optimization are summarized in this volume, which might also be considered a guide for inexperienced users. Eight papers are presented: three on subgradient optimization, four on descent methods, and one on applicability. The report also includes a set of nonsmooth optimization test problems and a comprehensive bibliography.
