## UNIVERSITY OF WATERLOO

## Nonlinear Optimization (CO 367)

E. de Klerk, C. Roos, and T. Terlaky

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## Chapter 0

## Introduction

Nothing takes place in the world whose meaning is not that of some maximum or minimum.
L. Euler (1707-1783)

### 0.1 The general nonlinear optimization problem

The general nonlinear optimization (NLO) problem can be written as follows:

$$
\begin{align*}
(N L O) & \inf \\
\text { s.t. } & f(x)  \tag{1}\\
& h_{i}(x)=0, \quad i \in I=\{1, \cdots, p\} \\
& g_{j}(x) \leq 0, \quad j \in J=\{1, \cdots, m\} \\
& x \in \mathcal{C},
\end{align*}
$$

where $x \in \mathbb{R}^{n}, \mathcal{C} \subseteq \mathbb{R}^{n}$ is a certain set and $f, h_{1}, \cdots, h_{p}, g_{1}, \cdots, g_{m}$ are functions defined on $\mathcal{C}$ (or on an open set that contains the set $\mathcal{C}$ ). The set of feasible solutions will be denoted by $\mathcal{F}$, hence

$$
\mathcal{F}=\left\{x \in \mathcal{C} \mid h_{i}(x)=0, \quad i=1, \cdots, p \text { and } g_{j}(x) \leq 0, \quad j=1, \cdots, m\right\}
$$

We will use the following standard terminology:

- The function $f$ is called the objective function of (NLO) and $\mathcal{F}$ is called the feasible set (or feasible region);
- If $\mathcal{F}=\emptyset$ then we say that problem (NLO) is infeasible;
- If $f$ is not bounded below on $\mathcal{F}$ then we say that problem (NLO) is unbounded;
- If the infimum of $f$ over $\mathcal{F}$ is attained at $\bar{x} \in \mathcal{F}$ then we call $\bar{x}$ an optimal solution (or minimum or minimizer) of (NLO) and $f(\bar{x})$ the optimal (objective) value of (NLO).

Example 0.1 As an example, consider minimization of the 'humpback function'(see Figure 1):

$$
\min x_{1}^{2}\left(4-2.1 x_{1}^{2}+\frac{1}{3} x_{1}^{4}\right)+x_{1} x_{2}+x_{2}^{2}\left(-4+4 x_{2}^{2}\right)
$$

subject to the constraints $-2 \leq x_{1} \leq 2$ and $-1 \leq x_{2} \leq 1$. Note that the feasible set here is simply the


Figure 1: Example of a nonlinear optimization problem with the 'humpback' function as objective function. The contours of the objective function are shown.
rectangle:

$$
\mathcal{F}=\left\{\left(x_{1}, x_{2}\right):-2 \leq x_{1} \leq 2,-1 \leq x_{2} \leq 1\right\}
$$

This NLO problem has two optimal solutions, namely $(0.0898,-0.717)$ and $(-0.0898,0.717)$, as one can (more or less) verify by looking at the contours of the objective function in Figure 1.

## Notation

Matrices will be denoted by capital letters $(A, B, P, \ldots)$, vectors by small Latin letters and components of vectors and matrices by the indexed letters [e.g. $z=\left(z_{1}, \ldots, z_{n}\right)$, $\left.A=\left(a_{i j}\right)_{i=1}^{m} \underset{j=1}{n}\right]$. For the purpose of matrix-vector multiplication, vectors in $\mathbb{R}^{n}$ will always be viewed as $n \times 1$ matrices (column vectors). Index sets will be denoted by $I, J$ and $K$.
We now define some classes of NLO problems. Recall that $f: \mathbb{R}^{n} \rightarrow \mathbb{R}$ is called a quadratic function if there is a square matrix $Q \in \mathbb{R}^{n \times n}$, a vector $c \in \mathbb{R}^{n}$ and a number $\gamma \in \mathbb{R}$ such that

$$
f(x)=\frac{1}{2} x^{T} Q x+c^{T} x+\gamma \text { for all } x \in \mathbb{R}^{n}
$$

If $Q=0$ then $f$ is called affine and if $Q=0$ and $\gamma=0$ then $f$ is called linear. We will abuse this terminology a bit by sometimes referring to affine functions as linear.

## Classification of nonlinear optimization problems

We now list a few important classes of optimization problems, with reference to the general problem (1):

Linear Optimization (LO): The functions $f, h_{1}, \cdots, h_{p}, g_{1}, \cdots, g_{m}$ are affine and the set $\mathcal{C}$ either equals to $\mathbb{R}^{n}$ or to the nonnegative orthant $\mathbb{R}_{+}^{n}$ of $\mathbb{R}^{n}$. Linear optimization is often called linear programming. The reader is probably familiar with the simplex algorithm for solving LO problems.

Unconstrained Optimization: The index sets $I$ and $J$ are empty and $\mathcal{C}=\mathbb{R}^{n}$.
Quadratic Optimization (QO): The objective function $f$ is quadratic, and all the constraint functions $h_{1}, \cdots, h_{p}, g_{1}, \cdots, g_{m}$ are affine and the set $\mathcal{C}$ is $\mathbb{R}^{n}$ or the positive orthant $\mathbb{R}_{+}^{n}$ of $\mathbb{R}^{n}$.

Quadratically Constrained Quadratic Optimization: Same as QO, except that the functions $g_{1}, \cdots, g_{m}$ are quadratic.

### 0.2 A short history of nonlinear optimization

As far as we know, Euclid's book Elements was the first mathematical textbook in the history of mankind (4th century B.C.). It contains the following optimization problem.
In a given triangle $A B C$ inscribe a parallelogram $A D E F$ such that $E F \| A B$ and $D E \| A C$ and such that the area of this parallelogram is maximal.


Figure 2: Illustration of Euclid's problem.

Let $H$ denote the height of the triangle, and let $b$ indicate the length of the edge $A C$. Every inscribed parallelogram of the required form is uniquely determined by choosing a vertex $F$ at a distance $x<b$ from $A$ on the edge $A C$ (see Figure 2).

Exercise 0.1 Show that Euclid's problem can be formulated as the quadratic optimization problem (QO):

$$
\begin{equation*}
\max _{0<x<b} \frac{H x(b-x)}{b} . \tag{2}
\end{equation*}
$$

Euclid could show that the maximum is obtained when $x=\frac{1}{2} b$, by using geometric reasoning. A unified methodology for solving nonlinear optimization problems would have to wait until the development of calculus in the 17 th century. Indeed, in any modern text on calculus we learn to solve problems like (2) by setting the derivative of the objective function $f(x):=\frac{H x(b-x)}{b}$ to zero, and solving the resulting equation to obtain $x=\frac{1}{2} b$.

This modern methodology is due to Fermat (1601-1665). Because of this work, Lagrange (1736 - 1813) stated clearly that he considered Fermat to be the inventor of calculus (as opposed to Newton (1643 - 1727) and Liebnitz (1646-1716) who were later locked in a bitter struggle for this honour). Lagrange himself is famous for extending the method of Fermat to solve (equality) constrained optimization problems by forming a function now known as the Lagrangian, and applying Fermat's method to the Lagrangian. In the words of Lagrange himself:

One can state the following general principle. If one is looking for the maximum or minimum of some function of many variables subject to the condition that these variables are related by a constraint given by one or more equations, then one should add to the function whose extremum is sought the functions that yield the constraint equations each multiplied by undetermined multipliers and seek the maximum or minimum of the resulting sum as if the variables were independent. The resulting equations, combined with the constraint equations, will serve to determine all unknowns.

To better understand what Lagrange meant, consider the general NLO problem with only equality constraints $\left(J=\emptyset\right.$ and $\mathcal{C}=\mathbb{R}^{n}$ in (1)):

$$
\begin{array}{ll}
\text { inf } & f(x) \\
\text { s.t. } & h_{i}(x)=0, \quad i \in I=\{1, \cdots, p\} \\
& x \in \mathbb{R}^{n} .
\end{array}
$$

Now define the Lagrangian function

$$
L(x, y):=f(x)+\sum_{i=1}^{p} y_{i} h_{i}(x)
$$

The new variables $y_{i}$ are called (Lagrange) multipliers, and are the undetermined multipliers Lagrange referred to. ${ }^{1}$ Now apply Fermat's method to find the minimum of the function $L(x, y)$, 'as if the variables $x$ and $y$ were independent'. In other words, solve the system of nonlinear equations defined by setting the gradient of $L$ to zero, and retaining the feasibility conditions $h_{i}(x)=0(i \in I)$ :

$$
\begin{equation*}
\nabla L(x, y)=0, h_{i}(x)=0(i \in I) \tag{3}
\end{equation*}
$$

If $x^{*}$ is an optimal solution of NLO then there now exists a vector $y^{*} \in \mathbb{R}^{p}$ such that $\left(x^{*}, y^{*}\right)$ is a solution of the nonlinear equations (3). We can therefore solve the nonlinear system (3) and the $x$-part of one of the solutions of (3) will yield the optimal solution of (NLO) (provided it exists). Note that it is difficult to know beforehand whether an optimal solution of (NLO) exists.

This brings us to the problem of solving a system of nonlinear equations. Isaac Newton lent his name to perhaps the most widely known algorithm for this problem. In conjunction with Fermat and Lagrange's methods, this yielded one of the first optimization algorithms. It is interesting to note that even today, Newton's optimization algorithm is the most widely used and studied algorithm for nonlinear optimization. The most recent optimization algorithms, namely interior point algorithms, use this method as their 'engine'.

The study of nonlinear optimization in the time of Fermat, Lagrange, Euler and Newton was driven by the realization that many physical principles in nature can be explained via optimization (extremum) principles. For example, the well known principle of Fermat for the refraction of light may be stated as:
in an inhomogeneous medium, light travels from one point to another along the path requiring the shortest time.
Similarly, it was known that many problems in (celestial) mechanics could be formulated as extremum problems.

We return to the problem of deciding whether (NLO) has an optimal solution at all. In the 19th century, Karl Weierstrass (1815-1897) proved the famous result known to any student of analysis - that a continuous function attains its infimum and supremum on a compact set. This gave a practical sufficient condition for the existence of optimal solutions.
In modern times, nonlinear optimization is used in optimal engineering design, finance, statistics and many other fields. It has been said that we live in the age of optimization, where everything has to be better and faster than before. Think of designing a car with minimal air resistance, a bridge of minimal weight that still meets essential specifications, a stock portfolio where the risk is minimal and the expected return high; the list is endless. If you can make a mathematical model of your decision problem, then you can optimize it!

[^0]
## Outline of this course

This short history of nonlinear optimization is of course far from complete and has served only to introduce some of the most important topics that will be studied in this course. In Chapters 1 and 2 we will study the methods of Fermat and Lagrange. Socalled duality theory based on the methodology of Lagrange will follow in Chapter 3. Then we will turn our attention to optimization algorithms in the remaining chapters. First we will study classical methods like the (reduced) gradient method and Newton's method (Chapter 4), before turning to the modern application of Newton's method in interior point algorithms (Chapter 6). Finally, we conclude with a chapter on special classes of structured nonlinear optimization problems that can be solved very efficiently by interior point algorithms.
In the rest of this chapter we give a few more examples of historical and practical problems to give some idea of the field of nonlinear optimization.

### 0.3 Some historical examples

These examples of historical nonlinear optimization problems are taken from the wonderful book Stories about Maxima and Minima by V.M. Tikhomirov (AMS, 1990). For more background and details the reader is referred to that book. We only state the problems here - the solutions will be presented later in this book in the form of exercises, once the reader has studied optimality conditions. Of course, the mathematical tools we will employ were not available to the originators of the problems. For the (ingenious) original historical solutions the reader is again referred to the book by Tikhomirov.

### 0.3.1 Tartaglia's problem

Niccolo Tartaglia (1500-1557) posed the following problem:
How do you divide the number 8 into two parts such that the result of multiplying the product of the parts by their difference will be maximal?
If we denote the unknown parts by $x_{1}$ and $x_{2}$, we can restate the problem as the nonlinear optimization problem:

$$
\begin{equation*}
\max x_{1} x_{2}\left(x_{1}-x_{2}\right) \tag{4}
\end{equation*}
$$

subject to the constraints

$$
x_{1}+x_{2}=8, x_{1} \geq 0, x_{2} \geq 0
$$

Tartaglia knew that the correct answer is $x_{1}=4+(4 / \sqrt{3}), x_{2}=4-(4 / \sqrt{3})$. How can one prove that this is correct? (Solution via Exercise 2.4.)

### 0.3.2 Kepler's problem

The famous astronomer Johannes Kepler was so intrigued by the geometry of wine barrels that he wrote a book about it in 1615: New solid geometry of wine barrels. In this work he considers the following problem (among others).

Given a sphere, inscribe a cylinder of maximal volume.
Kepler knew the cylinder of maximal volume is such that the ratio of its base diameter to the height is $\sqrt{2}$. (And of course the diagonal of the cylinder has the same length as the diameter of the sphere.) How can one show that this is correct? (Solution via Exercises 0.2 and 2.2.)

Exercise 0.2 Formulate Kepler's problem as a nonlinear optimization problem (NLO). $\triangleleft$

### 0.3.3 Steiner's problem

In the plane of a given triangle, find a point such that the sum of its distances from the vertices of the triangle is minimal.
This problem was included in the first book on optimization, namely On maximal and minimal values by Viviani in 1659.
If we denote the vertices of the triangle by the three given vectors $a \in \mathbb{R}^{2}, b \in \mathbb{R}^{2}$ and $c \in \mathbb{R}^{2}$, and let $x=\left[x_{1} x_{2}\right]^{T}$ denote the vector with the (unknown) coordinates of the point we are looking for, then we can formulate Steiner's problem as the following nonlinear optimization problem.

$$
\min _{x \in \mathbb{R}^{2}}\|x-a\|+\|x-b\|+\|x-c\|
$$

The solution is known as the Torricelli point. How can one find the Torricelli point for any given triangle? (Solution via Exercise 2.3.)

### 0.4 Quadratic optimization examples

In countless applications we wish to solve a linear system of the form $A x=b$. If this system is over-determined (more equations than variables), then we can still obtain the so-called least squares solution by solving the NLO problem:

$$
\begin{equation*}
\min _{x}\|A x-b\|^{2} \tag{5}
\end{equation*}
$$

Since

$$
\|A x-b\|^{2}=(A x-b)^{T}(A x-b)=x^{T} A^{T} A x-2 b^{T} A x+b^{T} b,
$$

it follows that problem (5) is a quadratic optimization (QO) problem.
Below we give examples of the least squares and other quadratic optimization problems.

### 0.4.1 The concrete mixing problem: least square estimation

In civil engineering, different sorts of concrete are needed for different purposes. One of the important characteristics of the concrete are its sand-and-gravel composition, i.e. what percentages of the stones in the sand-and-gravel mixture belong to a certain stone size categories. For each sort concrete the civil engineers can give an ideal sand-andgravel composition that ensures the desired strength with minimal cement content.

Unfortunately, in the sand-and-gravel mines, such ideal composition can not be found in general. The solution is to mix different sand-and-gravel mixtures in order to approximate the desired quality as closely as possible.

## Mathematical model

Let us assume that we have $n$ different stone size categories. The ideal mixture for our actual purpose is given by the vector $c=\left(c_{1}, c_{2}, \cdots, c_{n}\right)^{T}$, where $0 \leq c_{i} \leq 1$ for all $i=1, \cdots, n$ and $\sum_{i=1}^{n} c_{i}=1$. The components $c_{i}$ indicate what fraction of the sand-and-gravel mixture belongs to the $i$-th stone size category. Further, let assume that we can get sand-and-gravel mixtures from $m$ different mines, and the stone composition at each mine $j=1, \cdots, m$ is given by the vectors $A_{j}=\left(a_{1 j}, \cdots, a_{n j}\right)^{T}$, where $0 \leq a_{i j} \leq 1$ for all $i=1, \cdots, n$ and $\sum_{i=1}^{n} a_{i j}=1$. The goal is to find the best possible approximation of the ideal mixture by using the material offered by the $m$ mines.

Let $x=\left(x_{1}, \cdots, x_{m}\right)$ be a the vector of unknown percentages in the mixture, i.e.

$$
\sum_{j=1}^{m} x_{j}=1, \quad x_{j} \geq 0
$$

The resulting mixture composition

$$
z=\sum_{i=1}^{m} A_{j} x_{j}
$$

should be as close as possible to the ideal one, i.e. we have to minimize

$$
\|z-c\|^{2}=(z-c)^{T}(z-c)=\sum_{j=1}^{n}\left(z_{i}-c_{i}\right)^{2}
$$

This optimization problem is a linearly constrained QO problem. We can further simplify this problem by eliminating variable $z$. Then, introducing matrix $A=\left(A_{1}, \cdots, A_{m}\right)$ composed from the vectors $A_{j}$ as its columns, we have the following simple QO problem:

$$
\begin{aligned}
\min (A x-c)^{T}(A x-c) & \\
e^{T} x= & 1 \\
x \geq & 0 .
\end{aligned}
$$

Exercise 0.3 In the above concrete mixing problem the deviation of the mixture from the targeted ideal composition is measured by the Euclidean distance of the vectors $z=A x$ and $c$. The distance of two vectors can be measured alternatively by e.g. the $\|\cdot\|_{1}$ or by the $\|\cdot\|_{\infty}$ norms. Restate the mixing problem by using these norms and show that this way, in both cases, pure linear optimization problems can be obtained.

### 0.4.2 Portfolio analysis (mean-variance models)

An important application of the QO problem is the computation of an efficient frontier for mean-variance models, introduced by Markowitz [31]. Given assets with expected return $r_{i}$ and covariances $v_{i j}$, the problem is to find portfolios of the assets that have minimal variance given a level of total return, and maximal return given a level of total variance. Mathematically, if $x_{i}$ is the proportion invested in asset $i$ then the basic mean-variance problem is

$$
\min _{x}\left\{\frac{1}{2} x^{T} V x: e^{T} x=1, r^{T} x=\lambda, D x=d, x \geq 0\right\}
$$

where $e$ is an all-one vector, and $D x=d$ may represent additional constraints on the portfolios to be chosen (for instance those related to volatility of the portfolio). This problem can be viewed as a parametric QO problem, with parameter $\lambda$ representing the total return of investment. The so-called efficient frontier is then just the optimal value function.

Example 0.2 Consider the following MV-model,

$$
\min _{x}\left\{x^{T} V x: e^{T} x=1, r^{T} x=\lambda, x \geq 0\right\}
$$

where

$$
\begin{gathered}
V=\left[\begin{array}{ccccc}
0.82 & -0.23 & 0.155 & -0.013 & -0.314 \\
-0.23 & 0.484 & 0.346 & 0.197 & 0.592 \\
0.155 & 0.346 & 0.298 & 0.143 & 0.419 \\
-0.013 & 0.197 & 0.143 & 0.172 & 0.362 \\
-0.314 & 0.592 & 0.419 & 0.362 & 0.916
\end{array}\right], \\
r=\left(\begin{array}{lllll}
1.78 & 0.37 & 0.237 & 0.315 & 0.49
\end{array}\right)^{T} .
\end{gathered}
$$

One can check (e.g. by using MATLAB) that for $\lambda>1.780$ or $\lambda<0.237$ the QO problem is infeasible. For the values $\lambda \in[0.237,1.780]$ the QO problem has an optimal solution.

Exercise 0.4 A mathematical description of this and related portfolio problems is given at:
http://www-fp.mcs.anl.gov/otc/Guide/CaseStudies/port/formulations.html
Choose your own stock portfolio at the website:
http://www-fp.mcs.anl.gov/otc/Guide/CaseStudies/port/demo.html
and solve this problem remotely via internet to obtain the optimal way of dividing your capital between the stocks you have chosen. In doing this you are free to set the level of risk you are prepared to take. Give the mathematical description of the problem you have solved and report on your results.

## Chapter 1

## Convex Analysis

If the nonlinear optimization problem (NLO) has a convex objective function and the feasible set is a convex set, then the underlying mathematical structure is much richer than in the general case. For example, one can formulate necessary and sufficient conditions for the existence of optimal solutions in this case. It is therefore important to study convexity in some detail.

### 1.1 Convex sets and convex functions

Given two points $x^{1}$ and $x^{2}$ in $\mathbb{R}^{n}$, any point on the line connecting them is called a convex combinations of $x^{1}$ and $x^{2}$. Formally we have the following definition.

Definition 1.1 Let two points $x^{1}, x^{2} \in \mathbb{R}^{n}$ and $0 \leq \lambda \leq 1$ be given. Then the point

$$
x=\lambda x^{1}+(1-\lambda) x^{2}
$$

is a convex combination of the two points $x^{1}, x^{2}$.
The set $\mathcal{C} \subset \mathbb{R}^{n}$ is called convex, if all convex combinations of any two points $x^{1}, x^{2} \in$ $\mathcal{C}$ are again in $\mathcal{C}$.

In other words, the line segment connecting two arbitrary points of a convex set is contained in the set.

Figure 1.1 and Figure 1.2 show examples of convex and nonconvex sets in the plane.


Figure 1.1: Convex sets


Figure 1.2: Non convex sets

Exercise 1.1 We can define the convex combination of $k$ points as follows. Let the points $x^{1}, \cdots, x^{k} \in \mathbb{R}^{n}$ and $0 \leq \lambda_{1}, \cdots, \lambda_{k}$ with $\sum_{i=1}^{k} \lambda_{i}=1$ be given. Then the vector

$$
x=\sum_{i=1}^{k} \lambda_{i} x^{i}
$$

is a convex combination of the given points.
Prove that the set $\mathcal{C}$ is convex if and only if for any $k \geq 2$ all convex combinations of any $k$ points from $\mathcal{C}$ are also in $\mathcal{C}$.

The intersection of (possibly infinitely many) convex sets is again a convex set.
Theorem 1.2 Let $C_{i}(i=1, \ldots)$ be a collection of convex sets. The set

$$
C:=\cap_{i=1}^{\infty} C_{i}
$$

ic convex.

Exercise 1.2 Prove Theorem 1.2.

Another fundamental property of a convex set is that its closure is again a convex set.

Theorem 1.3 Let $C \subset \mathbb{R}^{n}$ be a convex set and let cl $(C)$ denote its closure. Then $\operatorname{cl}(C)$ is a convex set.

We now turn our attention to so-called convex functions. A parabola $f(x)=a x^{2}+$ $b x+c$ with $a>0$ is a familiar example of a convex function. Intuitively, it is easier to characterize minima of convex functions than minima of more general functions, and for this reason we will study convex functions in some detail.

Definition 1.4 A function $f: \mathcal{C} \rightarrow R$ defined on a convex set $\mathcal{C}$ is called convex if for all $x^{1}, x^{2} \in \mathcal{C}$ and $0 \leq \lambda \leq 1$ one has

$$
f\left(\lambda x^{1}+(1-\lambda) x^{2}\right) \leq \lambda f\left(x^{1}\right)+(1-\lambda) f\left(x^{2}\right)
$$

Exercise 1.3 Let $f: \mathbb{R}^{n} \mapsto \mathbb{R}$ be defined by $f(x)=\|x\|$ for some norm $\|\cdot\|$. Prove that $f$ is a convex function.

Exercise 1.4 Show that the following univariate functions are not convex:

$$
f(x)=\sin x, \quad f(x)=e^{-x^{2}}, \quad f(x)=x^{3} .
$$

Definition 1.5 The epigraph of a function $f: \mathcal{C} \rightarrow R$, where $\mathcal{C} \subset \mathbb{R}^{n}$, is the $(n+1)$ dimensional set

$$
\{(x, \tau): f(x) \leq \tau, x \in \mathcal{C}, \tau \in \mathbb{R}\}
$$

Figure 1.3 illustrates the above definition.

$\mathcal{C}$

Figure 1.3: The epigraph of a convex function $f$.

Exercise 1.5 Prove that the function $f: \mathcal{C} \rightarrow R$ defined on the convex set $\mathcal{C}$ is convex if and only if the epigraph of $f$ is a convex set.

We also will need the concept of a strictly convex function. These are convex functions with the nice property that - if a minimum of the function exists - then this minimum is unique.

Definition 1.6 $A$ (convex) function $f: \mathcal{C} \rightarrow \mathbb{R}$, defined on a convex set $\mathcal{C}$, is called strictly convex if for all $x^{1}, x^{2} \in \mathcal{C}$ and $0<\lambda<1$ one has

$$
f\left(\lambda x^{1}+(1-\lambda) x^{2}\right)<\lambda f\left(x^{1}\right)+(1-\lambda) f\left(x^{2}\right) .
$$

We have seen in Exercise 1.5 that a function is convex if and only if its epigraph is convex.
Also, the next exercise shows that a quadratic function is convex if and only if the matrix $Q$ in its definition is positive-semidefinite (PSD).

Exercise 1.6 Let a symmetric matrix $Q \in \mathbb{R}^{n \times n}$, a vector $c \in \mathbb{R}^{n}$ and a number $\gamma \in \mathbb{R}$ be given. Prove that the quadratic function

$$
\frac{1}{2} x^{T} Q x+c^{T} x+\gamma
$$

is convex on $\mathbb{R}^{n}$ if and only if the matrix $Q$ is $P S D$, and strictly convex if and only if $Q$ is positive definite.

Exercise 1.7 Decide whether the following quadratic functions are convex or not. (Hint: use the result from the previous exercise.)

$$
f(x)=x_{1}^{2}+2 x_{1} x_{2}+x_{2}^{2}+5 x_{1}-x_{2}+\frac{1}{2}, \quad f(x)=x_{1}^{2}+x_{2}^{2}+x_{3}^{2}-2 x_{1} x_{2}-2 x_{1} x_{3}-2 x_{2} x_{3} .
$$

If we multiply a convex function by -1 , then we get a so-called concave function.
Definition 1.7 A function $f: \mathcal{C} \rightarrow \mathbb{R}$, defined on a convex set $\mathcal{C}$, is called (strictly) concave if the function $-f$ is (strictly) convex.

Note that we can change the problem of maximizing a concave function into the problem of minimizing a convex function.

Now we review some further properties of convex sets and convex functions that are necessary to understand and analyze convex optimization problems. First we review some elementary properties of convex sets.

### 1.2 More on convex sets

### 1.2.1 Convex hulls and extremal sets

For any set $\mathcal{S} \subset \mathbb{R}^{n}$ we define the smallest convex set that contains it, the so-called convex hull of $\mathcal{S}$, as follows.

Definition 1.8 Let $\mathcal{S} \subset \mathbb{R}^{n}$ be an arbitrary set. The set

$$
\operatorname{conv}(\mathcal{S}):=\left\{x \mid x=\sum_{i=1}^{k} \lambda_{i} x^{i}, x^{i} \in \mathcal{S}, i=1, \cdots, k ; \lambda_{i} \in[0,1], \sum_{i=1}^{k} \lambda_{i}=1, k \geq 1\right\}
$$

is called the convex hull of the set $\mathcal{S}$.
Observe that $\operatorname{conv}(\mathcal{S})$ is generated by taking all possible convex combinations of points from $\mathcal{S}$.
We now define some important convex subsets of a given convex set $\mathcal{C}$, namely the so-called extremal sets, that play an important role in convex analysis. Informally, an extremal set $\mathcal{E} \subset \mathcal{C}$ is a convex subset of $\mathcal{C}$ with the following property: if any point on the line segment connecting two points $x^{1} \in \mathcal{C}$ and $x^{2} \in \mathcal{C}$ lies in $\mathcal{E}$, then the two points $x^{1}$ and $x^{2}$ must also lie in $\mathcal{E}$. The faces of a polytope are familiar examples of extreme sets of the polytope.

Definition 1.9 The convex set $\mathcal{E} \subseteq \mathcal{C}$ is an extremal set of the convex set $\mathcal{C}$ if, for all $x^{1}, x^{2} \in \mathcal{C}$ and $0<\lambda<1$, one has $\lambda x^{1}+(1-\lambda) x^{2} \in \mathcal{E}$ only if $x^{1}, x^{2} \in \mathcal{E}$.

An extremal set consisting of only one point is called an extremal point. Observe that extremal sets are convex by definition, and the convex set $\mathcal{C}$ itself is always an extremal set of $\mathcal{C}$. It is easy to verify the following result.

Lemma 1.10 If $\mathcal{E}^{1} \subseteq \mathcal{C}$ is an extremal set of the convex set $\mathcal{C}$ and $\mathcal{E}^{2} \subseteq \mathcal{E}^{1}$ is an extremal set of $\mathcal{E}^{1}$ then $\mathcal{E}^{2}$ is an extremal set of $\mathcal{C}$.
*Proof: Let $x, y \in \mathcal{C}, 0<\lambda<1$ and $z_{\lambda}=\lambda x+(1-\lambda) y \in \mathcal{E}^{2}$. Due to $\mathcal{E}^{2} \subseteq \mathcal{E}^{1}$ we have $z_{\lambda} \in \mathcal{E}^{1}$, moreover $x, y \in \mathcal{E}^{1}$ because $\mathcal{E}^{1}$ is an extremal set of $\mathcal{C}$. Finally, because $\mathcal{E}^{2}$ is an extremal set of $\mathcal{E}^{1}$, $x, y \in \mathcal{E}^{1}$ and $z_{\lambda} \in \mathcal{E}^{2}$ we conclude that $x, y \in \mathcal{E}^{2}$ and thus $\mathcal{E}^{2}$ is an extremal set of $\mathcal{C}$.

Example 1.11 Let $\mathcal{C}$ be the cube $\left\{x \in \mathbb{R}^{3} \mid 0 \leq x \leq 1\right\}$, then the vertices are extremal points, the edges are 1-dimensional extremal sets, the faces are 2-dimensional extremal sets, and the whole cube is a 3 -dimensional extremal set of itself.


Example 1.12 Let $\mathcal{C}$ be the cylinder $\left\{x \in \mathbb{R}^{3} \mid x_{1}^{2}+x_{2}^{2} \leq 1,0 \leq x_{3} \leq 1\right\}$, then

- the points on the circles $\left\{x \in \mathbb{R}^{3} \mid x_{1}^{2}+x_{2}^{2}=1, x_{3}=1\right\}$ and $\left\{x \in \mathbb{R}^{3} \mid x_{1}^{2}+x_{2}^{2}=1, x_{3}=0\right\}$ are the extremal points,
- the lines $\left\{x \in \mathbb{R}^{3} \mid x_{1}=a, x_{2}=b, 0 \leq x_{3} \leq 1\right\}$, with $a \in[-1,1]$ and $b=\sqrt{1-a^{2}}$ or $b=-\sqrt{1-a^{2}}$, are the 1-dimensional extremal sets,
- the faces $\left\{x \in \mathbb{R}^{3} \mid x_{1}^{2}+x_{2}^{2} \leq 1, x_{3}=1\right\}$ and $\left\{x \in \mathbb{R}^{3} \mid x_{1}^{2}+x_{2}^{2} \leq 1, x_{3}=0\right\}$ are the 2-dimensional extremal sets, and
- the cylinder itself is the only 3 -dimensional extremal set.


Example 1.13 Let $f(x)=x^{2}$ and let $\mathcal{C}$ be the epigraph of $f$, then all points $\left(x_{1}, x_{2}\right)$ such that $x_{2}=x_{1}^{2}$ are extremal points. The epigraph itself is the only two dimensional extremal set.


Lemma 1.14 Let $\mathcal{C}$ be a closed convex set. Then all extremal sets of $\mathcal{C}$ are closed.
In the above examples we have pointed out extremal sets of different dimension without giving a formal definition of what the dimension of a convex set is. To this end, recall from linear algebra that if $\mathcal{L}$ is a (linear) subspace of $\mathbb{R}^{n}$ and $a \in \mathbb{R}^{n}$ then $a+\mathcal{L}$ is called an affine subspace of $\mathbb{R}^{n}$. By definition, the dimension of $a+\mathcal{L}$ is the dimension of $\mathcal{L}$.

Definition 1.15 The smallest affine space $a+\mathcal{L}$ containing a convex set $\mathcal{C} \subseteq \mathbb{R}^{n}$ is the so-called affine hull of $\mathcal{C}$ and denoted by $\operatorname{aff}(\mathcal{C})$. The dimension of $\mathcal{C}$ is defined as the dimension of $\operatorname{aff}(\mathcal{C})$.

Given two points $x^{1} \in \mathcal{C}$ and $x^{2} \in \mathcal{C}$, we call any point that lies on the (infinite) line that passes through $x^{1}$ and $x^{2}$ an affine combination of $x^{1}$ and $x^{2}$. Formally we have the following definition.

Definition 1.16 Let two points $x^{1}, x^{2} \in \mathbb{R}^{n}$ and $\lambda \in \mathbb{R}$ be given. Then the point

$$
x=\lambda x^{1}+(1-\lambda) x^{2}
$$

is an affine combination of the two points $x^{1}, x^{2}$.
Observe that in defining the affine combination we do not require that the coefficients $\lambda$ and $1-\lambda$ are from the interval $[0,1]$, while this was required in the definition of the convex combination of points.

Exercise 1.8 Let $\mathcal{C} \subset \mathbb{R}^{n}$ be defined by

$$
\mathcal{C}=\left\{x \mid \sum_{i=1}^{n} x_{i}=1, x \geq 0\right\}
$$

The set $\mathcal{C}$ is usually called the standard simplex in $\mathbb{R}^{n}$.

1. Give the extreme points of $\mathcal{C}$; Motivate your answer.
2. Show that $\mathcal{C}=\operatorname{conv}\left\{e^{1}, \ldots, e^{n}\right\}$, where $e^{i}$ is the $i$ th standard unit vector.
3. What is $\operatorname{aff}(\mathcal{C})$ in this case?
4. What is the dimension of $\mathcal{C}$ ? Motivate your answer.

Exercise 1.9 Let $\mathcal{C} \subseteq \mathbb{R}^{n}$ be a given convex set and $k \geq 2$ a given integer. Prove that

$$
\operatorname{aff}(\mathcal{C})=\left\{z \mid z=\sum_{i=1}^{k} \lambda_{i} x^{i}, \sum_{i=1}^{k} \lambda_{i}=1, \lambda_{i} \in \mathbb{R}, x^{i} \in \mathcal{C}, \forall i\right\}
$$

Exercise 1.10 Let $\mathcal{E}$ be an extremal set of the convex set $\mathcal{C}$. Prove that $\mathcal{E}=\operatorname{aff}(\mathcal{E}) \cap \mathcal{C}$. (Hint: Use Exercise 1.9 with $k=2$.)

Lemma 1.17 Let $\mathcal{E}^{2} \subset \mathcal{E}^{1} \subseteq \mathcal{C}$ be two extremal sets of the convex set $\mathcal{C}$. Then $\operatorname{dim}\left(\mathcal{E}^{2}\right)<\operatorname{dim}\left(\mathcal{E}^{1}\right)$.
${ }^{*}$ Proof: Because $\mathcal{E}^{2} \subset \mathcal{E}^{1}$ we have $\operatorname{aff}\left(\mathcal{E}^{2}\right) \subseteq \operatorname{aff}\left(\mathcal{E}^{1}\right)$. Further, by Exercise 1.10,

$$
\mathcal{E}^{2}=\operatorname{aff}\left(\mathcal{E}^{2}\right) \cap \mathcal{E}^{1} .
$$

If we assume to the contrary that $\operatorname{dim}\left(\mathcal{E}^{2}\right)=\operatorname{dim}\left(\mathcal{E}^{1}\right)$ then we have $\operatorname{aff}\left(\mathcal{E}^{2}\right)=\operatorname{aff}\left(\mathcal{E}^{1}\right)$ and so

$$
\mathcal{E}^{2}=\operatorname{aff}\left(\mathcal{E}^{2}\right) \cap \mathcal{E}^{1}=\operatorname{aff}\left(\mathcal{E}^{1}\right) \cap \mathcal{E}^{1}=\mathcal{E}^{1}
$$

contradicting the assumption $\mathcal{E}^{2} \subset \mathcal{E}^{1}$.

Lemma 1.18 Let $\mathcal{C}$ be a nonempty compact (closed and bounded) convex set. Then $\mathcal{C}$ has at least one extremal point.
*Proof: Let $\mathcal{F} \subseteq \mathcal{C}$ be the set of points in $\mathcal{C}$ which are furthest from the origin. The set of such points is not empty, because $\mathcal{C}$ is bounded and closed and the norm function is continuous. We claim that any point $z \in \mathcal{F}$ is an extremal point of $\mathcal{C}$.
Let us assume to the contrary that $z \in \mathcal{F}$ is not an extremal point. Then there exist points $x, y \in \mathcal{C}$, both different from $z$, and a $\lambda \in(0,1)$ such that

$$
z=\lambda x+(1-\lambda) y
$$

Further, we have $\|x\| \leq\|z\|$ and $\|y\| \leq\|z\|$ because $z$ maximizes the norm of the points over $\mathcal{C}$. Thus by the triangle inequality

$$
\|z\| \leq \lambda\|x\|+(1-\lambda)\|y\| \leq\|z\|
$$

which implies that $\|z\|=\|x\|=\|y\|$, i.e. all the three point $x, y, z$ are on the surface of the $n$ dimensional sphere with radius $\|z\|$ and centered at the origin. This is a contradiction, because these three different points lie on the same line as well. The lemma is proved.

Observe, that the above proof does not require the use of the origin as reference point. We could choose any point $u \in \mathbb{R}^{n}$ and prove that the furthest point $z \in \mathcal{C}$ from $u$ is an extremal point of $\mathcal{C}$.
The following theorem shows that a compact convex set is completely determined by its extremal points.

Theorem 1.19 (Krein-Milman Theorem) Let $\mathcal{C}$ be a compact convex set. Then $\mathcal{C}$ is the convex hull of its extreme points.

Exercise 1.11 Let $f$ be a continuous, concave function defined on a compact convex set $\mathcal{C}$. Show that the minimum of $f$ is attained at an extreme point of $\mathcal{C}$. (Hint: Use the KreinMilman Theorem.)

### 1.2.2 Convex cones

In what follows we define and prove some elementary properties of convex cones.

Definition 1.20 The $\operatorname{set} \mathcal{C} \subset \mathbb{R}^{n}$ is a convex cone if it is a convex set and for all $x \in \mathcal{C}$ and $0 \leq \lambda$ one has $\lambda x \in \mathcal{C}$.

## Example 1.21

- The set $\mathcal{C}=\left\{\left(x_{1}, x_{2}\right) \in \mathbb{R}^{2} \mid x_{2} \geq 2 x_{1}, x_{2} \geq-\frac{1}{2} x_{1}\right\}$ is a convex cone in $\mathbb{R}^{2}$.
- The set $\mathcal{C}^{\prime}=\left\{\left(x_{1}, x_{2}, x_{3}\right) \in \mathbb{R}^{3} \mid x_{1}^{2}+x_{2}^{2} \leq x_{3}^{2}, x_{3} \geq 0\right\}$ is a convex cone in $\mathbb{R}^{3}$.


Definition 1.22 A convex cone is called pointed if it does not contain any subspace except the origin.

A pointed convex cone could be defined equivalently as a convex cone that does not contain any line.

Lemma 1.23 $A$ convex cone $\mathcal{C}$ is pointed if and only if the origin 0 is an extremal point of $\mathcal{C}$.
*Proof: If the convex cone $\mathcal{C}$ is not pointed, then it contains a nontrivial subspace, in particular, it contains a one dimensional subspace, i.e. a line $\mathcal{L}$ going through the origin. Let $0 \neq x \in \mathcal{L}$, and then we have $-x \in \mathcal{L}$ as well. From here we have $0=\frac{1}{2} x+\frac{1}{2}(-x) \in \mathcal{C}$, i.e. 0 is not an extremal point.

If the convex cone $\mathcal{C}$ is pointed, then it does not contain any subspace, except the origin 0 . In that case we show that 0 is an extremal point of $\mathcal{C}$. If we assume to the contrary that there exists
$0 \neq x^{1}, x^{2} \in \mathcal{C}$ and a $\lambda \in(0,1)$ such that $0=\lambda x^{1}+(1-\lambda) x^{2}$, then we derive $x^{1}=-\frac{1-\lambda}{\lambda} x^{2}$. This implies that the line through $x^{1}$, the origin 0 and $x^{2}$ is in $\mathcal{C}$, contradicting the assumption that $\mathcal{C}$ is pointed.

Example 1.24 If a convex cone $\mathcal{C} \in \mathbb{R}^{2}$ is not pointed, then it is either

- a line through the origin,
- a halfspace, or
- $\mathbb{R}^{2}$.

Example 1.25 Let $V_{1}, V_{2}$ be two planes through the origin in $\mathbb{R}^{3}$, given by the following equations,

$$
\begin{aligned}
V_{1}: & =\left\{x \in \mathbb{R}^{3} \mid x_{3}=a_{1} x_{1}+a_{2} x_{2}\right\} \\
V_{2}: & =\left\{x \in \mathbb{R}^{3} \mid x_{3}=b_{1} x_{1}+b_{2} x_{2}\right\}
\end{aligned}
$$

then the convex set

$$
\mathcal{C}=\left\{x \in \mathbb{R}^{3} \mid x_{3} \geq a_{1} x_{1}+a_{2} x_{2}, x_{3} \leq b_{1} x_{1}+b_{2} x_{2}\right\}
$$

is a non-pointed cone.

$V_{1}: x_{3}=2 x_{1}-x_{2}$ $V_{2}: x_{3}=x_{1}+3 x_{2}$

Every convex cone $\mathcal{C}$ has an associated dual cone. By definition, every vector in the dual cone has a nonnegative inner product with every vector in $\mathcal{C}$.

Definition 1.26 Let $\mathcal{C} \subseteq \mathbb{R}^{n}$ be a convex cone. The dual cone $\mathcal{C}^{*}$ is defined by

$$
\mathcal{C}^{*}:=\left\{z \in \mathbb{R}^{n} \mid x^{T} z \geq 0 \text { for all } x \in \mathcal{C}\right\}
$$

In the literature the dual cone $\mathcal{C}^{*}$ is frequently referred to as the polar or positive polar of the cone $\mathcal{C}$.

Exercise 1.12 Prove that $\left(\mathbb{R}_{+}^{n}\right)^{*}=\mathbb{R}_{+}^{n}$, i.e. the nonnegative orthant is a self-dual cone. $\triangleleft$

Exercise 1.13 Let $\mathcal{S}_{n}$ denote the set of $n \times n$, symmetric positive semidefinite matrices.
(i) Prove that $\mathcal{S}_{n}$ is a convex cone.
(ii) Prove that $\left(\mathcal{S}_{n}\right)^{*}=\mathcal{S}_{n}$, i.e. $\mathcal{S}_{n}$ is a self-dual cone.

Exercise 1.14 Prove that the dual cone $\mathcal{C}^{*}$ is a closed convex cone.

An important, deep theorem $[38,42]$ says that the dual of the dual cone $\left(\mathcal{C}^{*}\right)^{*}$ is the closure $\overline{\mathcal{C}}$ of the cone $\mathcal{C}$.

An important cone in the study of convex optimization is the so-called recession cone. Given an (unbounded) convex feasible set $\mathcal{F}$ and some $x \in \mathcal{F}$, the recession cone of $\mathcal{F}$ consists of all the directions one can travel in without ever leaving $\mathcal{F}$, when starting from $x$. Surprisingly, the recession cone does not depend on the choice of $x$.

Lemma 1.27 Let us assume that the convex set $\mathcal{C}$ is closed and not bounded. Then
(i) for each $x \in \mathcal{C}$ there is a non-zero vector $z \in \mathbb{R}^{n}$ such that $x+\lambda z \in \mathcal{C}$ for all $\lambda \geq 0$, i.e. the set $R(x)=\{z \mid x+\lambda z \in \mathcal{C}, \lambda \geq 0\}$ is not empty;
(ii) the set $R(x)$ is a closed convex cone (the so-called recession cone at $x$ );
(iii) the cone $R(x)=\mathcal{R}$ is independent of $x$, thus it is 'the' recession cone of the convex set $\mathcal{C}$; $^{1}$
(iv) $\mathcal{R}$ is a pointed cone if and only if $\mathcal{C}$ has at least one extremal point.
*Proof: $\quad(i)$ Let $x \in \mathcal{C}$ be given. Because $\mathcal{C}$ is unbounded, then there is a sequence of points $x^{1}, \cdots, x^{k}, \cdots$ such that $\left\|x^{k}-x\right\| \rightarrow \infty$. Then the vectors

$$
y^{k}=\frac{x^{k}-x}{\left\|x^{k}-x\right\|}
$$

are in the unit sphere. The unit sphere is a closed convex, i.e. compact set, hence there exists an accumulation point $\bar{y}$ of the sequence $y^{k}$. We claim that $\bar{y} \in R(x)$. To prove this we take any $\lambda>0$ and prove that $x+\lambda \bar{y} \in \mathcal{C}$. This claim follows from the following three observations: 1 . If we omit all the points from the sequence $y^{k}$ for which $\left\|x-x^{k}\right\|<\lambda$ then $\bar{y}$ is still an accumulation point of the remaining sequence $y^{k_{i}}$. 2. Due to the convexity of $\mathcal{C}$ the points

$$
x+\lambda y^{k_{i}}=x+\frac{\lambda}{\left\|x^{k_{i}}-x\right\|}\left(x^{k_{i}}-x\right)=\left(1-\frac{\lambda}{\left\|x^{k_{i}}-x\right\|}\right) x+\frac{\lambda}{\left\|x^{k_{i}}-x\right\|} x^{k_{i}}
$$

are in $\mathcal{C}$. 3. Because $\mathcal{C}$ is closed, the accumulation point $x+\lambda \bar{y}$ of the sequence $x+\lambda y^{k_{i}} \in \mathcal{C}$ is also in $\mathcal{C}$. The proof of the first statement is complete.

[^1](ii) The set $R(x)$ is a cone, because $z \in R(x)$ imply $\mu z \in R(x)$. The convexity of $R(x)$ easily follows from the convexity of $\mathcal{C}$. Finally, if $\left.z^{i} \in R(x)\right)$ for all $i=1,2, \cdots$ and $\bar{z}=\lim _{i \rightarrow \infty} z^{i}$, then for each $\lambda \geq 0$ the closedness of $\mathcal{C}$ and $x+\lambda z^{i} \in \mathcal{C}$ imply that
$$
\lim _{i \rightarrow \infty}\left(x+\lambda z^{i}\right)=x+\lambda \bar{z} \in \mathcal{C},
$$
hence $\bar{z} \in R(x)$ proving that $R(x)$ is closed.
(iii) Let $x^{1}, x^{2} \in \mathcal{C}$. We have to show that $z \in R\left(x^{1}\right)$ imply $z \in R\left(x^{2}\right)$. Let us assume to the contrary that $z \notin R\left(x^{2}\right)$, i.e. there is an $\alpha>0$ such that $x^{2}+\alpha z \notin \mathcal{C}$. Due to $z \in R\left(x^{1}\right)$ we have $x^{1}+(\lambda+\alpha) z \in \mathcal{C}$ for all $\alpha, \lambda \geq 0$. Using the convexity of $\mathcal{C}$ we have that the point
$$
x_{\lambda}^{2}=x^{2}+\frac{\alpha}{\lambda+\alpha}\left(x^{1}-x^{2}+(\lambda+\alpha) z\right)=x^{2}\left(1-\frac{\alpha}{\lambda+\alpha}\right)+\frac{\alpha}{\lambda+\alpha}\left(x^{1}+(\lambda+\alpha) z\right)
$$
is in $\mathcal{C}$. Further the limit point
$$
\lim _{\lambda \rightarrow \infty} x_{\lambda}^{2}=x^{2}+\alpha z,
$$
due to the closedness of $\mathcal{C}$, is also in $\mathcal{C}$.
(iv) We leave the proof of this part as an exercise.

Exercise 1.15 Prove part (iv) of Lemma 1.27.

Corollary 1.28 The nonempty closed convex set $\mathcal{C}$ is bounded if and only if its recession cone $\mathcal{R}$ consists of the zero vector alone.
*Proof: If $\mathcal{C}$ is bounded, then it contains no half line, thus for each $x \in \mathcal{C}$ the set $R(x)=\{0\}$, i.e. $\mathcal{R}=\{0\}$.

The other part of the proof follows form item (i) of Lemma 1.27.

Example 1.29 Let $\mathcal{C}$ be the epigraph of $f(x)=\frac{1}{x}$. Then every point on the curve $x_{2}=\frac{1}{x_{1}}$ is an extreme point of $\mathcal{C}$. For an arbitrary point $x=\left(x_{1}, x_{2}\right)$ the recession cone is given by

$$
R(x)=\left\{z \in \mathbb{R}^{2} \mid z_{1}, z_{2} \geq 0\right\}
$$

Hence, $R=R(x)$ is independent of $x$, and $R$ is a pointed cone of $\mathcal{C}$.


Lemma 1.30 If the convex set $\mathcal{C}$ is closed and has an extremal point, then each extremal set of $\mathcal{C}$ has at least one extremal point as well.
*Proof: Let us assume to the contrary that an extremal set $\mathcal{E} \subset \mathcal{C}$ has no extremal point. Then by item (iv) of Lemma 1.27 the recession cone of $\mathcal{E}$ is not pointed, i.e. it contains a line. By statement (iii) of the same lemma, this line is contained in the recession cone of $\mathcal{C}$ as well. Applying (iv) of Lemma 1.27 again we conclude that $\mathcal{C}$ cannot have an extremal point. This is a contradiction, the lemma is proved.

Lemma 1.31 Let $\mathcal{C}$ be a convex set and $\mathcal{R}$ be its recession cone. If $\mathcal{E}$ is an extremal set of $\mathcal{C}$ the recession cone $\mathcal{R}_{\mathcal{E}}$ of $\mathcal{E}$ is an extremal set of $\mathcal{R}$.
${ }^{*}$ Proof: $\quad$ Clearly $\mathcal{R}_{\mathcal{E}} \subseteq \mathcal{R}$. Let us assume that $\mathcal{R}_{\mathcal{E}}$ is not an extremal set of $\mathcal{R}$. Then there are vectors $z^{1}, z^{2} \in \mathcal{R}, z^{1} \notin \mathcal{R}_{\mathcal{E}}$ and a $\lambda \in(0,1)$ such that $z=\lambda z^{1}+(1-\lambda) z^{2} \in \mathcal{R}_{\mathcal{E}}$. Finally, for a certain $\alpha>0$ and $x \in \mathcal{E}$ we have

$$
x^{1}=x+\alpha z^{1} \in \mathcal{C} \backslash \mathcal{E}, \quad x^{2}=x+\alpha z^{2} \in \mathcal{C}
$$

and

$$
\lambda x^{1}+(1-\lambda) x^{2}=x+\alpha z \in \mathcal{E}
$$

contradicting the extremality of $\mathcal{E}$.

### 1.2.3 The relative interior of a convex set

The standard simplex in $\mathbb{R}^{3}$ was defined as the set

$$
\left\{x \in \mathbb{R}^{3} \mid \sum_{i=1}^{3} x_{i}=1, x \geq 0\right\}
$$

The interior of this convex set is empty, but intuitively the points that do not lie on the 'boundary' of the simplex do form a 'sort of interior'. This leads us to a generalized concept of the interior of a convex set, namely the relative interior. If the convex set is full-dimension (i.e. $\mathcal{C} \in \mathbb{R}^{n}$ has dimension $n$ ), then the concepts of interior and relative interior coincide.

Definition 1.32 Let a convex set $\mathcal{C}$ be given. The point $x \in \mathcal{C}$ is in the relative interior of $\mathcal{C}$ if for all $\bar{x} \in \mathcal{C}$ there exists $\tilde{x} \in \mathcal{C}$ and $0<\lambda<1$ such that $x=\lambda \bar{x}+(1-\lambda) \tilde{x}$. The set of relative interior points of the set $\mathcal{C}$ will be denoted by $\mathcal{C}^{0}$.

The relative interior $\mathcal{C}^{0}$ of a convex set $\mathcal{C}$ is obviously a subset of the convex set. We will show that the relative interior $\mathcal{C}^{0}$ is a relatively open (i.e. it coincides with its relative interior) convex set.

Example 1.33 Let $\mathcal{C}=\left\{x \in \mathbb{R}^{3} \mid x_{1}^{2}+x_{2}^{2} \leq 1, x_{3}=1\right\}$ and $\mathcal{L}=\left\{x \in \mathbb{R}^{3} \mid x_{3}=0\right\}$, then $\mathcal{C} \subset$ aff $(\mathcal{C})=$ $(0,0,1)+\mathcal{L}$. Hence, $\operatorname{dim} \mathcal{C}=2$ and $\mathcal{C}^{0}=\left\{x \in \mathbb{R}^{3} \mid x_{1}^{2}+x_{2}^{2}<1, x_{3}=1\right\}$.


Lemma $1.34 \operatorname{Let} \mathcal{C} \subset \mathbb{R}^{n}$ be a convex set. Then for each $x \in \mathcal{C}^{0}, y \in \mathcal{C}$ and $0<\lambda \leq 1$ we have

$$
z=\lambda x+(1-\lambda) y \in \mathcal{C}^{0} \subseteq \mathcal{C}
$$

*Proof: Let $u \in \mathcal{C}$ be an arbitrary point. Then we have to show that for each $u \in \mathcal{C}$ there is an $\bar{u} \in \mathcal{C}$ and a $0<\rho<1$ such that $z=\rho \bar{u}+(1-\rho) u$. The proof is constructive.

Because $x \in \mathcal{C}^{0}$, by Definition 1.32 there is an $0<\alpha<1$ such that the point

$$
v:=\frac{1}{\alpha} x+\left(1-\frac{1}{\alpha}\right) u
$$

is in $\mathcal{C}$. Let

$$
\bar{u}=\vartheta v+(1-\vartheta) y, \quad \text { where } \quad \vartheta=\frac{\lambda \alpha}{\lambda \alpha+1-\lambda} .
$$

Due to the convexity of $\mathcal{C}$ we have $\bar{u} \in \mathcal{C}$. Finally, let us define $\rho=\lambda \alpha+1-\lambda$. Then one can easily verify that $0<\rho<1$ and

$$
z=\lambda x+(1-\lambda) y=\rho \bar{u}+(1-\rho) u
$$

Figure 1.4 illustrates the above construction.


Figure 1.4: If $x \in \mathcal{C}^{0}$ and $y \in \mathcal{C}$, then the point $z$ is in $\mathcal{C}^{0}$.

A direct consequence of the above lemma is the following.

Corollary 1.35 The relative interior $\mathcal{C}^{0}$ of a convex set $\mathcal{C} \subset \mathbb{R}^{n}$ is convex.
Lemma 1.36 Let $\mathcal{C}$ be a convex set. Then $\left(\mathcal{C}^{0}\right)^{0}=\mathcal{C}^{0}$. Moreover, if $\mathcal{C}$ is nonempty then its relative interior $\mathcal{C}^{0}$ is nonempty as well.
*Proof: The proof of this lemma is quite technical. For a proof the reader is referred to the excellent books of Rockafellar [38] and Stoer and Witzgall [42].

### 1.3 More on convex functions

Now we turn our attention to convex functions.

### 1.3.1 Basic properties of convex functions

In this section we have collected some useful facts about convex functions.
Lemma 1.37 Let $f$ be a convex function defined on the convex set $\mathcal{C}$. Then $f$ is continuous on the relative interior $\mathcal{C}^{0}$ of $\mathcal{C}$.
*Proof: Let $p \in \mathcal{C}^{0}$ be an arbitrary point. Without loss of generality we may assume that $\mathcal{C}$ is full dimensional, $p$ is the origin and $f(p)=0$.
Let us first consider the one dimensional case. Because 0 is in the interior of the domain $\mathcal{C}$ of $f$ we have a $v>0$ such that $v \in \mathcal{C}$ and $-v \in \mathcal{C}$ as well. Let us consider the two linear functions

$$
\ell_{1}(x):=x \frac{f(v)}{v} \quad \text { and } \quad \ell_{2}(x):=x \frac{f(-v)}{-v} .
$$

One easily checks that the convexity of $f$ implies the following relations:

- $\ell_{1}(x) \geq f(x)$ if $x \in[0, v]$;
- $\ell_{1}(x) \leq f(x)$ if $x \in[-v, 0]$;
- $\ell_{2}(x) \geq f(x)$ if $x \in[-v, 0]$;
- $\ell_{2}(x) \leq f(x)$ if $x \in[0, v]$.

Then by defining

$$
h(x):=\max \left\{\ell_{1}(x), \ell_{2}(x)\right\} \quad \text { and } \quad g(x):=\min \left\{\ell_{1}(x), \ell_{2}(x)\right\}
$$

on the interval $[-v, v]$ we have

$$
g(x) \leq f(x) \leq h(x) .
$$

The linear functions $\ell_{1}(x)$ and $\ell_{2}(x)$ are obviously continuous, thus the functions $h(x)$ and $g(x)$ are continuous as well. By observing the relations $f(0)=h(0)=g(0)=0$ it follows that the function $f(x)$ is continuous at the point 0 .

We use an analogous construction for the $n$-dimensional case. Let us assume again that 0 is an interior point of $\mathcal{C}$ and $f(0)=0$. Let $v^{1}, \cdots, v^{n}, v^{n+1}$ be vectors such that the convex set

$$
\left\{x \mid x=\sum_{i=1}^{n+1} \lambda_{i} v^{i}, \quad \lambda_{i} \in[0,1], \sum_{i=1}^{n+1} \lambda_{i}=1\right\}
$$

equals the space $\mathbb{R}^{n}$. For all $i=1, \cdots, n+1$ let the linear functions (hyperplanes) $L_{i}(x)$ be defined by $n+1$ of their values: $L_{i}(0)=0$ and $L_{i}\left(v^{j}\right)=f\left(v^{j}\right)$ for all $j \neq i$. Let us further define

$$
h(x):=\max \left\{L_{1}(x), \cdots, L_{n+1}(x)\right\} \quad \text { and } \quad g(x):=\min \left\{L_{1}(x), \cdots, L_{n+1}(x)\right\}
$$

Then one easily proves that the functions $g(x)$ and $h(x)$ are continuous, $f(0)=h(0)=g(0)=0$ and

$$
g(x) \leq f(x) \leq h(x)
$$

Thus the function $f(x)$ is continuous at the point 0 .

Exercise 1.16 Prove that the functions $g(x)$ and $h(x)$, defined in the proof above, are continuous, $f(0)=h(0)=g(0)=0$ and

$$
g(x) \leq f(x) \leq h(x)
$$

Note that $f$ can be discontinuous on the relative boundary $\mathcal{C} \backslash \mathcal{C}^{0}$.
Example 1.38 The function

$$
f(x)= \begin{cases}x^{2} & \text { for }-1<x<1 \\ x^{2}+1 & \text { otherwise }\end{cases}
$$

is not continuous on $\mathbb{R}$, and it is also not convex. If $f$ is only defined on $\mathcal{C}=\{x \in \mathbb{R} \mid-1 \leq x \leq 1\}$ then $f$ is still not continuous, but it is continuous on $\mathcal{C}^{0}$ and convex on $\mathcal{C}$.



The following result, called Jensen's inequality, is simply a generalization of the inequality $f\left(\lambda x^{1}+(1-\lambda) x^{2}\right) \leq \lambda f\left(x^{1}\right)+(1-\lambda) f\left(x^{2}\right)$.

Lemma 1.39 (Jensen inequality) Let $f$ be a convex function defined on a convex set $\mathcal{C} \subseteq \mathbb{R}^{n}$. Let the points $x^{1}, \cdots, x^{k} \in \mathcal{C}$ and $\lambda_{1}, \cdots, \lambda_{k} \geq 0$ with $\sum_{i=1}^{k} \lambda_{i}=1$ be given . Then

$$
f\left(\sum_{i=1}^{k} \lambda_{i} x^{i}\right) \leq \sum_{i=1}^{k} \lambda_{i} f\left(x^{i}\right) .
$$

*Proof: The proof is by induction on $k$. If $k=2$ then the statement is true by Definition 1.4. Let us assume that the statement holds for a given $k \geq 2$, then we prove that it also holds for $k+1$.

Let the points $x^{1}, \cdots, x^{k}, x^{k+1} \in \mathcal{C}$ and $\lambda_{1}, \cdots, \lambda_{k}, \lambda_{k+1} \geq 0$ with $\sum_{i=1}^{k+1} \lambda_{i}=1$ be given. If at most $k$ of the $\lambda_{i}, 1 \leq i \leq k+1$ coefficients are nonzero then, by leaving out the points $x^{i}$ with zero coefficients, the inequality directly follows from the inductive assumption. Now let us consider the case when all the coefficients $\lambda_{i}$ are nonzero. Then by convexity of the set $\mathcal{C}$ we have that

$$
\tilde{x}=\sum_{i=1}^{k} \frac{\lambda_{i}}{\sum_{j=1}^{k} \lambda_{j}} x^{i} \in \mathcal{C}
$$

Further

$$
\begin{aligned}
f\left(\sum_{i=1}^{k+1} \lambda_{i} x^{i}\right) & =f\left(\sum_{j=1}^{k} \lambda_{j} \sum_{i=1}^{k} \frac{\lambda_{i}}{\sum_{j=1}^{k} \lambda_{j}} x^{i}+\lambda_{k+1} x^{k+1}\right) \\
& =f\left(\left[\sum_{j=1}^{k} \lambda_{j}\right] \tilde{x}+\lambda_{k+1} x^{k+1}\right) \\
& \leq\left[\sum_{j=1}^{k} \lambda_{j}\right] f(\tilde{x})+\lambda_{k+1} f\left(x^{k+1}\right) \\
& \leq\left[\sum_{j=1}^{k} \lambda_{j}\right]\left(\sum_{i=1}^{k} \frac{\lambda_{i}}{\sum_{j=1}^{k} \lambda_{j}} f\left(x^{i}\right)\right)+\lambda_{k+1} f\left(x^{k+1}\right) \\
& =\sum_{i=1}^{k+1} \lambda_{i} f\left(x^{i}\right)
\end{aligned}
$$

where the first inequality follows from the convexity of the function $f$ (Definition 1.4) and, at the second inequality, the inductive assumption was used. The proof is complete.

The reader can easily prove the following two lemmas by applying the definitions. We leave the proofs as exercises.

Lemma 1.40 Let $f^{1}, \cdots, f^{k}$ be convex functions defined on a convex set $\mathcal{C} \subseteq \mathbb{R}^{n}$. Then

- for all $\lambda_{1}, \cdots, \lambda_{k} \geq 0$ the function

$$
f(x)=\sum_{i=1}^{k} \lambda_{i} f^{i}(x)
$$

is convex;

- the function

$$
f(x)=\max _{1 \leq i \leq k} f^{i}(x)
$$

is convex.

Definition 1.41 The function $h: \mathbb{R} \rightarrow \mathbb{R}$ is called

- monotonically non-decreasing if for all $t_{1}<t_{2} \in \mathbb{R}$ one has $h\left(t_{1}\right) \leq h\left(t_{2}\right)$;
- strictly monotonically increasing if for all $t_{1}<t_{2} \in \mathbb{R}$ one has $h\left(t_{1}\right)<h\left(t_{2}\right)$.

Lemma 1.42 Let $f$ be a convex function on the convex set $\mathcal{C} \subseteq \mathbb{R}^{n}$ and $h: \mathbb{R} \rightarrow \mathbb{R}$ be a convex monotonically non-decreasing function. Then the composite function $h(f(x))$ : $\mathcal{C} \rightarrow \mathbb{R}$ is convex.

Exercise 1.17 Prove Lemma 1.40 and Lemma 1.42.

Exercise 1.18 Assume that the function $h$ in Lemma 1.42 is not monotonically non-decreasing. Give a concrete example that in this case the statement of the lemma fails.

Definition 1.43 Let a convex function $f: \mathcal{C} \rightarrow \mathbb{R}$ defined on the convex set $\mathcal{C}$ be given. Let $\alpha \in \mathbb{R}$ be an arbitrary number. The set $\mathcal{D}_{\alpha}=\{x \in \mathcal{C} \mid f(x) \leq \alpha\}$ is called a level set of the function $f$.

Lemma 1.44 If $f$ is a convex function on the convex set $\mathcal{C}$ then for all $\alpha \in \mathbb{R}$ the level set $\mathcal{D}_{\alpha}$ is a (possibly empty) convex set.
${ }^{*}$ Proof: Let $x, y \in \mathcal{D}_{\alpha}$ and $0 \leq \lambda \leq 1$. Then we have $f(x) \leq \alpha, f(y) \leq \alpha$ and we may write

$$
f(\lambda x+(1-\lambda) y) \leq \lambda f(x)+(1-\lambda) f(y) \leq \lambda \alpha+(1-\lambda) \alpha=\alpha .
$$

Here the first inequality followed from the convexity of the function $f$. The lemma is proved.

### 1.3.2 On the derivatives of convex functions

The first and second order derivatives of (sufficiently differentiable) convex functions have some interesting and useful properties which we review in this section. The reader may recall from elementary calculus that a univariate function $f$ on $\mathbb{R}$ is convex if $f^{\prime \prime}(x) \geq 0$ for all $x \in \mathbb{R}$ (assuming sufficient differentiability), and that such a function attains its minimum at some $\bar{x} \in \mathbb{R}$ if and only if $f^{\prime}(\bar{x})=0$. We will work towards generalizing these results to multivariate convex functions.
Recall that the gradient $\nabla f$ of the function $f$ is defined as the vector formed by the partial derivatives $\frac{\partial f}{\partial x_{i}}$ of $f$. Further we introduce the concept of directional derivative.

Definition 1.45 Let $x \in \mathbb{R}^{n}$ and a direction (vector) $s \in \mathbb{R}^{n}$ be given. The directional derivative $\delta f(x, s)$ of the function $f$, at point $x$, in the direction $s$, is defined as

$$
\delta f(x, s)=\lim _{\lambda \rightarrow 0} \frac{f(x+\lambda s)-f(x)}{\lambda}
$$

if the above limit exists.
If the function $f$ is continuously differentiable then $\frac{\partial f}{\partial x_{i}}=\delta f\left(x, e^{i}\right)$ where $e^{i}$ is the $i-$ th unit vector. This implies the following result.

Lemma 1.46 If the function $f$ is continuously differentiable then for all $s \in \mathbb{R}^{n}$ we have

$$
\delta f(x, s)=\nabla f(x)^{T} s
$$

The Hesse matrix (or Hessian) $\nabla^{2} f(x)$ of the function $f$ at a point $x \in \mathcal{C}$ is composed of the second order partial derivatives of $f$ as

$$
\left(\nabla^{2} f(x)\right)_{i j}=\frac{\partial^{2} f(x)}{\partial x_{i} \partial x_{j}} \quad \text { for all } \quad i, j=1, \cdots, n
$$

Lemma 1.47 Let $f$ be a function defined on a convex set $\mathcal{C} \subseteq \mathbb{R}^{n}$. The function $f$ is convex if and only if the function $\phi(\lambda)=f(x+\lambda s)$ is convex on the interval $[0,1]$ for all $x \in \mathcal{C}$ and $x+s \in \mathcal{C}$.
${ }^{*}$ Proof: Let us assume that $f$ is a convex function. Then we prove that $\phi(\lambda)$ is convex on the interval $[0,1]$. Let $\lambda_{1}, \lambda_{2} \in[0,1]$ and $0 \leq \alpha \leq 1$. Then one has

$$
\begin{aligned}
\phi\left(\alpha \lambda_{1}+(1-\alpha) \lambda_{2}\right) & =f\left(x+\left[\alpha \lambda_{1}+(1-\alpha) \lambda_{2}\right] s\right) \\
& =f\left(\alpha\left[x+\lambda_{1} s\right]+(1-\alpha)\left[x+\lambda_{2} s\right]\right) \\
& \leq \alpha f\left(x+\lambda_{1} s\right)+(1-\alpha) f\left(x+\lambda_{2} s\right) \\
& =\alpha \phi\left(\lambda_{1}\right)+(1-\alpha) \phi\left(\lambda_{2}\right)
\end{aligned}
$$

proving the first part of the lemma.
On the other hand, if $\phi(\lambda)$ is convex on the interval $[0,1]$ for each $x, x+s \in \mathcal{C}$ then the convexity of the function $f$ can be proved as follows. For given $y, x \in \mathcal{C}$ let us define $s:=y-x$. Then we write

$$
\begin{aligned}
f(\alpha y+(1-\alpha) x) & =f(x+\alpha(y-x))=\phi(\alpha)=\phi(\alpha 1+(1-\alpha) 0) \\
& \leq \alpha \phi(1)+(1-\alpha) \phi(0)=\alpha f(y)+(1-\alpha) f(x)
\end{aligned}
$$

The proof of the lemma is complete.

Example 1.48 Let $f(x)=x_{1}^{2}+x_{2}^{2}$ and let $E_{f}$ be the epigraph of $f$. For every $s \in \mathbb{R}^{2}$, we can define the half-plane $V_{s} \subset \mathbb{R}^{3}$ as $\left\{(x, y) \in \mathbb{R}^{2} \times \mathbb{R} \mid x=\mu s, \mu>0\right\}$. Now, for $x=(0,0)$ the epigraph of $\phi(\lambda)=f(x+\lambda s)=f(\lambda s)$ equals $V_{s} \cap E_{f}$, which is a convex set. Hence, $\phi(\lambda)$ is convex.


Lemma 1.49 Let $f$ be a continuously differentiable function on the open convex set $\mathcal{C} \subseteq \mathbb{R}^{n}$. Then the following statements are equivalent.

1. The function $f$ is convex on $\mathcal{C}$.
2. For any two vectors $x, \bar{x} \in \mathcal{C}$ one has

$$
\nabla f(x)^{T}(\bar{x}-x) \leq f(\bar{x})-f(x) \leq \nabla f(\bar{x})^{T}(\bar{x}-x)
$$

3. For any $x \in \mathcal{C}$, and any $s \in \mathbb{R}^{n}$ such that $x+s \in \mathcal{C}$, the function $\phi(\lambda)=f(x+\lambda s)$ is continuously differentiable on the open interval $(0,1)$ and $\phi^{\prime}(\lambda)=s^{T} \nabla f(x+\lambda s)$, which is a monotonically non-decreasing function.
*Proof: First we prove that 1 implies 2. Let $0 \leq \lambda \leq 1$ and $x, \bar{x} \in \mathcal{C}$. Then the convexity of $f$ implies

$$
f(\lambda \bar{x}+(1-\lambda) x) \leq \lambda f(\bar{x})+(1-\lambda) f(x) .
$$

This can be rewritten as

$$
\frac{f(x+\lambda(\bar{x}-x))-f(x)}{\lambda} \leq f(\bar{x})-f(x) .
$$

Taking the limit as $\lambda \rightarrow 0$ and applying Lemma 1.46 the left-hand-side inequality of 2 follows. As one interchanges the role $x$ and $\bar{x}$, the right-hand-side inequality is obtained analogously.
Now we prove that 2 implies 3. Let $x, x+s \in \mathcal{C}$ and $0 \leq \lambda_{1}, \lambda_{2} \leq 1$. When we apply the inequalities of 2 with the points $x+\lambda_{1} s$ and $x+\lambda_{2} s$ the following relations are obtained.

$$
\left(\lambda_{2}-\lambda_{1}\right) \nabla f\left(x+\lambda_{1} s\right)^{T} s \leq f\left(x+\lambda_{2} s\right)-f\left(x+\lambda_{1} s\right) \leq\left(\lambda_{2}-\lambda_{1}\right) \nabla f\left(x+\lambda_{2} s\right)^{T} s,
$$

hence

$$
\left(\lambda_{2}-\lambda_{1}\right) \phi^{\prime}\left(\lambda_{1}\right) \leq \phi\left(\lambda_{2}\right)-\phi\left(\lambda_{1}\right) \leq\left(\lambda_{2}-\lambda_{1}\right) \phi^{\prime}\left(\lambda_{2}\right) .
$$

Assuming $\lambda_{1}<\lambda_{2}$ we have

$$
\phi^{\prime}\left(\lambda_{1}\right) \leq \frac{\phi\left(\lambda_{2}\right)-\phi\left(\lambda_{1}\right)}{\lambda_{2}-\lambda_{1}} \leq \phi^{\prime}\left(\lambda_{2}\right),
$$

proving that the function $\phi^{\prime}(\lambda)$ is monotonically non-decreasing.
Finally we prove that 3 implies 1 . We only have to prove that $\phi(\lambda)$ is convex if $\phi^{\prime}(\lambda)$ is monotonically non-decreasing. Let us take $0<\lambda_{1}<\lambda_{2}<1$ where $\phi\left(\lambda_{1}\right) \leq \phi\left(\lambda_{2}\right)$. Then for $0 \leq \alpha \leq 1$ we may write

$$
\begin{aligned}
(1-\alpha) \phi\left(\lambda_{1}\right) & +\alpha \phi\left(\lambda_{2}\right)-\phi\left((1-\alpha) \lambda_{1}+\alpha \lambda_{2}\right) \\
& =\alpha\left[\phi\left(\lambda_{2}\right)-\phi\left(\lambda_{1}\right)\right]-\left[\phi\left((1-\alpha) \lambda_{1}+\alpha \lambda_{2}\right)-\phi\left(\lambda_{1}\right)\right] \\
& =\alpha\left(\lambda_{2}-\lambda_{1}\right)\left(\int_{0}^{1} \phi^{\prime}\left(\lambda_{1}+t\left(\lambda_{2}-\lambda_{1}\right)\right) \mathrm{d} t-\int_{0}^{1} \phi^{\prime}\left(\lambda_{1}+t \alpha\left(\lambda_{2}-\lambda_{1}\right)\right) \mathrm{d} t\right) \\
& \geq 0 .
\end{aligned}
$$

The expression for the derivative of $\phi$ is left as an exercise in calculus (Exercise 1.19). The proof of the Lemma is complete.
Figure 1.5 illustrates the inequalities at statement 2 of the lemma.
Exercise 1.19 Let $f: \mathbb{R}^{n} \mapsto \mathbb{R}$ be twice continuously differentiable and let $x \in \mathbb{R}^{n}$ and $s \in \mathbb{R}^{n}$ be given. Define $\phi: \mathbb{R} \mapsto \mathbb{R}$ via $\phi(\lambda)=f(x+\lambda s)$. Prove that

$$
\phi^{\prime}(\lambda)=s^{T} \nabla f(x+\lambda s)
$$

and

$$
\phi^{\prime \prime}(\lambda)=s^{T} \nabla^{2} f(x+\lambda s) s
$$



Figure 1.5: Inequalities derived for the gradient of a convex function $f$.

Lemma 1.50 Let $f$ be a twice continuously differentiable function on the open convex set $\mathcal{C} \subseteq \mathbb{R}^{n}$. The function $f$ is convex if and only if its Hesse matrix $\nabla^{2} f(x)$ is PSD for all $x \in \mathcal{C}$.

Proof: Let us take an arbitrary $x \in \mathcal{C}$ and $s \in \mathbb{R}^{n}$ such that $x+s \in \mathcal{C}$, and define $\phi(\lambda)=f(x+\lambda s)$. If $f$ is convex, then $\phi^{\prime}(\lambda)$ is monotonically non-decreasing. This implies that $\phi^{\prime \prime}(\lambda)$ is nonnegative for each $x \in \mathcal{C}$ and $0 \leq \lambda \leq 1$. Thus

$$
s^{T} \nabla^{2} f(x) s=\phi^{\prime \prime}(0) \geq 0
$$

proving the positive semidefiniteness of the Hessian $\nabla^{2} f(x)$.
On the other hand, if the Hessian $\nabla^{2} f(x)$ is positive semidefinite for each $x \in \mathcal{C}$, then

$$
s^{T} \nabla^{2} f(x+\lambda s) s=\phi^{\prime \prime}(\lambda) \geq 0
$$

i.e. $\phi^{\prime}(\lambda)$ is monotonically non-decreasing proving the convexity of $f$ by Lemma 1.49. The theorem is proved.

Exercise 1.20 Prove the following statement analogously as Lemma 1.50 was proved.
Let $f$ be a twice continuously differentiable function on the open convex set $\mathcal{C}$. Then $f$ is strictly convex if its Hesse matrix $\nabla^{2} f(x)$ is positive definite (PD) for all $x \in \mathcal{C}$. $\triangleleft$

Exercise 1.21 Give an example of a twice continuously differentiable strictly convex function $f$ where $\nabla^{2} f(x)$ is not positive definite (PD) for all $x$ in the domain of $f$.

## Chapter 2

## Optimality conditions

We consider two cases separately. First optimality conditions for unconstrained optimization are considered. Then optimality conditions for some special constrained optimization problems are derived.

### 2.1 Optimality conditions for unconstrained optimization

Consider the problem

$$
\begin{equation*}
\text { minimize } \quad f(x), \tag{2.1}
\end{equation*}
$$

where $x \in \mathbb{R}^{n}$ and $f: \mathbb{R}^{n} \rightarrow \mathbb{R}$ is a differentiable function. First we define local and global minima of the above problem.

Definition 2.1 Let a function $f: \mathbb{R}^{n} \rightarrow \mathbb{R}$ be given.
A point $\bar{x} \in \mathbb{R}^{n}$ is a local minimum of the function $f$ if there is an $\epsilon>0$ such that $f(\bar{x}) \leq f(x)$ for all $x \in \mathbb{R}^{n}$ when $\|\bar{x}-x\| \leq \epsilon$.
A point $\bar{x} \in \mathbb{R}^{n}$ is a strict local minimum of the function $f$ if there is an $\epsilon>0$ such that $f(\bar{x})<f(x)$ for all $x \in \mathbb{R}^{n}(x \neq \bar{x})$ when $\|\bar{x}-x\| \leq \epsilon$.
A point $\bar{x} \in \mathbb{R}^{n}$ is a global minimum of the function $f$ if $f(\bar{x}) \leq f(x)$ for all $x \in \mathbb{R}^{n}$.
A point $\bar{x} \in \mathbb{R}^{n}$ is a strict global minimum of the function $f$ if $f(\bar{x})<f(x)$ for all $x \in \mathbb{R}^{n}(x \neq \bar{x})$.

Convex functions possess the appealing property that a local minimum is global.
Example 2.2 Consider the convex function $f_{1}: \mathbb{R} \rightarrow \mathbb{R}$ defined as follows.

$$
f_{1}(x)= \begin{cases}-x+1 & \text { if } x<0 \\ 1 & \text { if } 0 \leq x \leq 1 \\ x & \text { if } x>1\end{cases}
$$



The point $\bar{x}=0$ is a global minimum of the function $f_{1}$ because $f_{1}(\bar{x}) \leq f_{1}(x)$ for all $x \in \mathbb{R}$. Because $\bar{x}=0$ is a global minimum, it follows immediately that it is also a local minimum of the function $f_{1}$. The point $\bar{x}=0$ is neither a strict local nor a strict global minimum point of $f_{1}$ because for any $\epsilon>0$ we can find an $x$ for which $f_{1}(\bar{x})=f_{1}(x)$ applies with $\|\bar{x}-x\| \leq \epsilon$.

Now let us consider the non-convex function $f_{1}: \mathbb{R} \rightarrow \mathbb{R}$ defined as

$$
f_{2}(x)= \begin{cases}-x & \text { if } x<2 \\ 2 & \text { if }-2 \leq x \leq-1 \\ -x+1 & \text { if }-1<x<0 \\ 1 & \text { if } 0 \leq x \leq 1 \\ x & \text { if } x>1\end{cases}
$$



The point $\bar{x}=0$ is a global minimum of the function $f_{2}$ because $f_{2}(\bar{x}) \leq f_{2}(x)$ for all $x \in \mathbb{R}$. Because it is a global minimum it is at the same a local minimum as well. The point $\bar{x}=0$ is neither a strict local, nor a strict global minimum of the function $f_{2}$ because for any $\epsilon>0$ we can find an $x$ for which $f_{2}(\bar{x})=f_{2}(x)$ applies with $\|\bar{x}-x\| \leq \epsilon$.

The point $x^{*}=-2$ is also a local minimum of the function $f_{2}$ because $f_{2}\left(x^{*}\right) \leq f_{2}(x)$ for all $x \in \mathbb{R}$ when $\left\|x^{*}-x\right\| \leq \epsilon$, with $0<\epsilon<1$. It is not a strict local minimum because $f_{2}\left(x^{*}\right) \not \leq f_{2}(x)$ for all $x \in \mathbb{R}$ when $\left\|x^{*}-x\right\|<\epsilon$, with $\epsilon>0$. The point $x^{*}=-2$ is not a global minimum of $f_{2}$ because $f_{2}(-2)>f_{2}(0)$.

Example 2.3 Consider the convex function $f_{1}(x)=x^{2}$ where $x \in \mathbb{R}$.


The point $\bar{x}=0$ is a strict local minimum of the function $f_{1}$ because $f_{1}(\bar{x})<f_{1}(x)$ for all $x \in \mathbb{R}$ when $\|\bar{x}-x\|<\epsilon$, with $\epsilon>0$. The point $\bar{x}=0$ is also a strict global minimum of the function $f_{1}$ because $f_{1}(\bar{x})<f_{1}(x)$ for all $x \in \mathbb{R}$.

Consider the non-convex function $f_{2}: \mathbb{R} \rightarrow \mathbb{R}$ defined as

$$
f_{2}(x)= \begin{cases}(x+3)^{2}+3 & \text { if } x<-2 \\ x^{2} & \text { if } x \geq-2\end{cases}
$$



The point $\bar{x}=0$ is a strict local minimum as well as an strict global minimum for the function $f_{2}$, because $f_{2}(\bar{x})<f_{2}(x)$ for all $x \in \mathbb{R}$ when $\|\bar{x}-x\|<\epsilon$, with $\epsilon>0$, and for all $x \in \mathbb{R}$. The point $x^{*}=-3$ is a strict local minimum because $f_{2}\left(x^{*}\right)<f_{2}(x)$ for all $x \in \mathbb{R}$ when $\left\|x^{*}-x\right\|<\epsilon$, with $0<\epsilon<1$. The point $x^{*}=-3$ is not a strict global minimum, because $f_{2}(-3)>f_{2}(0)$.

Lemma 2.4 Any (strict) local minimum of a convex function $f$ is a (strict) global minimum of $f$ as well.

Exercise 2.1 Prove Lemma 2.4.

Now we arrive at the famous result of Fermat that says that a necessary condition for $\bar{x}$ to be a minimum of a continuously differentiable function $f$ is that $\nabla f(\bar{x})=0$.

Theorem 2.5 (Fermat) Let $f$ be continuously differentiable. If the point $\bar{x} \in \mathbb{R}^{n}$ is a local minimum of the function $f$ then $\nabla f(\bar{x})=0$.

Proof: As $\bar{x}$ is a minimum, one has

$$
f(\bar{x}) \leq f(\bar{x}+\lambda s) \text { for all } s \in \mathbb{R}^{n} \text { and } \lambda \in R .
$$

By bringing $f(\bar{x})$ to the right hand side and dividing by $\lambda>0$ we have

$$
0 \leq \frac{f(\bar{x}+\lambda s)-f(\bar{x})}{\lambda} .
$$

Taking the limit as $\lambda \rightarrow 0$ results in

$$
0 \leq \delta f(\bar{x}, s)=\nabla f(\bar{x})^{T} s \quad \text { for all } s \in \mathbb{R}^{n}
$$

As $s \in \mathbb{R}^{n}$ is arbitrary we conclude that $\nabla f(\bar{x})=0$.

Remark: In the above theorem it is enough to assume that the partial derivatives of $f$ exist. The same proof applies if we choose $e^{i}$, the standard unit vectors instead of the arbitrary direction $s$.

Exercise 2.2 Consider Kepler's problem as formulated in Exercise 0.2.

1. Show that Kepler's problem can be written as the problem of minimizing a nonlinear univariate function on an open interval.
2. Show that the solution given by Kepler is indeed optimal by using Theorem 2.5.

Observe that the above theorem contains only a one sided implication. It does not say anything about a minimum of $f$ if $\nabla f(\bar{x})=0$. Such points are not necessarily minimum points. These points are called stationary points. Think of the stationary (inflection) point $x=0$ of the univariate function $f(x)=x^{3}$. In other words, Fermat's result only gave a necessary condition for a minimum, namely $\nabla f(\bar{x})=0$. We will now see that this is also a sufficient condition if $f$ is convex.

Theorem 2.6 Let $f$ be a continuously differentiable convex function. The point $\bar{x} \in \mathbb{R}^{n}$ is a minimum of the function $f$ if and only if $\nabla f(\bar{x})=0$.

Proof: As $\bar{x}$ is a minimum of $f$ then by Theorem 2.5 we have $\nabla f(\bar{x})=0$. On the other hand, if $f$ is a convex function and $\nabla f(\bar{x})=0$ then

$$
f(x)-f(\bar{x}) \geq \nabla f(\bar{x})^{T}(x-\bar{x})=0 \quad \text { for all } \quad x \in \mathbb{R}^{n}
$$

hence the theorem is proved.

Exercise 2.3 We return to Steiner's problem (see Section 0.3.3) of finding the Torricelli point of a given triangle, that was defined as the solution of the optimization problem

$$
\begin{equation*}
\min _{x \in \mathbb{R}^{2}}\|x-a\|+\|x-b\|+\|x-c\| \tag{2.2}
\end{equation*}
$$

where $a, b$, and $c$ are given vectors in $\mathbb{R}^{2}$ that form the vertices of the given triangle.

1. Show that the objective function is convex.
2. Give necessary and sufficient conditions for a minimum of (2.2). (In other words, give the equations that determine the Torricelli point. You may assume that all three angles of the triangle are smaller that $\frac{2 \pi}{3}$.)
3. Find the Torricelli point of the triangle with vertices $(0,0),(3,0)$ and $(1,2)$.

If $f$ is a twice continuously differentiable (not necessarily convex) function then second order sufficient conditions for local minima are derived as follows. Let $\nabla^{2} f(x)$ denote the Hesse matrix of $f$ at the point $x$.

Theorem 2.7 Let $f$ be a twice continuously differentiable function. If at a point $\bar{x} \in$ $\mathbb{R}^{n}$ it holds that $\nabla f(\bar{x})=0$ and $\nabla^{2} f(x)$ is positive semidefinite in an $\epsilon$-neighborhood $(\epsilon>0)$ of $\bar{x}$ then the point $\bar{x}$ is a local minimum of the function $f$. (If we assume positive definiteness then we get a strict local minimum.)

Proof: Taking the Taylor series expansion of $f$ at $\bar{x}$ we have
$f(x)=f(\bar{x}+(x-\bar{x}))=f(\bar{x})+\nabla f(\bar{x})^{T}(x-\bar{x})+\frac{1}{2}(x-\bar{x})^{T} \nabla^{2} f(\bar{x}+\alpha(x-\bar{x}))(x-\bar{x})$
for some $0 \leq \alpha \leq 1$. Using the assumptions we have the result $f(x) \geq f(\bar{x})$ as $x$ is in the neighborhood of $\bar{x}$ where the Hesse matrix is positive semidefinite.

Corollary 2.8 Let $f$ be a twice continuously differentiable function. If at $\bar{x} \in \mathbb{R}^{n}$ the gradient $\nabla f(\bar{x})=0$ and the Hessian $\nabla^{2} f(\bar{x})$ is positive definite then the point $\bar{x}$ is a strict local minimum of the function $f$.

Proof: Since $f$ is twice continuously differentiable, it follows from the positive definiteness of the Hesse matrix at $\bar{x}$ that it is positive definite in a neighborhood of $\bar{x}$. Hence the claim follows from theorem 2.7.

### 2.2 Optimality conditions for constrained optimization

The following theorem generalizes the optimality conditions for a convex function on $\mathbb{R}^{n}$ (Theorem 2.6), by replacing $\mathbb{R}^{n}$ by any relatively open convex set $\mathcal{C} \subseteq \mathbb{R}^{n}$.

Theorem 2.9 Let us consider the optimization problem $\min \{f(x): x \in \mathcal{C}\}$ where $\mathcal{C}$ is a relatively open convex set and $f$ is a convex differentiable function. The point $\bar{x}$ is an optimal solution of this problem if and only if $\nabla f(\bar{x})^{T} s=0$ for all $s \in \mathcal{L}$, where $\mathcal{L}$ denotes the linear subspace with $\operatorname{aff}(\mathcal{C})=x+\mathcal{L}$ for any $x \in \mathcal{C}$. Here $\operatorname{aff}(\mathcal{C})$ denotes the affine hull of $\mathcal{C}$.

Proof: Let $s \in \mathcal{L}$ and $\lambda \in \mathbb{R}$. If $\bar{x}$ is a minimum, one has

$$
f(\bar{x}) \leq f(\bar{x}+\lambda s) \text { if } \bar{x}+\lambda s \in \mathcal{C}
$$

Note that $\bar{x}+\lambda s \in \operatorname{aff}(\mathcal{C})$ since $s \in \mathcal{L}$, and $\bar{x}+\lambda s \in \mathcal{C}$ if $\lambda$ is sufficiently small, since $\mathcal{C}$ is a relatively open set.

By bringing $f(\bar{x})$ to the right hand side and dividing by $\lambda>0$ we have

$$
0 \leq \frac{f(\bar{x}+\lambda s)-f(\bar{x})}{\lambda} \quad \text { for all } \quad s \in \mathcal{L}
$$

if $\lambda>0$ is sufficiently small. Taking the limit as $\lambda \downarrow 0$ results in

$$
0 \leq \delta f(\bar{x}, s)=\nabla f(\bar{x})^{T} s \quad \text { for all } \quad s \in \mathcal{L} .
$$

We conclude that $\nabla f(\bar{x})^{T} s=0$ for all $s \in \mathcal{L}$.
Conversely, if $f$ is a convex function and $\nabla f(\bar{x})^{T} s=0$ for all $s \in \mathcal{L}$ then, for any $x \in \mathcal{C}$,

$$
f(x)-f(\bar{x}) \geq \nabla f(\bar{x})^{T}(x-\bar{x})=0
$$

since $s=(x-\bar{x}) \in \mathcal{L}$, hence the theorem is proved.

A crucial assumption of the above lemma is that the set $\mathcal{C}$ is a relatively open set. In general this is not the case because the level sets of convex optimization problems are closed. However as we will see later the barrier function approach will result in such relatively open feasible sets. This is an important feature of interior point methods that will be discussed later on. If the set of feasible solutions is not relatively open, similar results by using similar techniques can be derived (see Theorem 2.14).

Exercise 2.4 We return to Tartaglia's problem (4) in Section 0.3.1.

1. Eliminate one of the variables and show that the resulting problem can be written as the problem of minimizing a univariate convex function on an open interval.
2. Show that the answer given by Tartaglia is indeed the optimal solution, by applying Theorem 2.9.

Now let us consider the general convex optimization problem, as given earlier in (1), but without equality constraints.

$$
\begin{array}{rll}
(C O) & \min & f(x) \\
& \text { s.t. } & g_{j}(x) \leq 0, \quad j=1, \cdots, m  \tag{2.3}\\
& x \in \mathcal{C},
\end{array}
$$

where $\mathcal{C} \subseteq \mathbb{R}^{n}$ is a convex set and $f, g_{1}, \cdots, g_{m}$ are convex functions on $\mathcal{C}$ (or on an open set that contains the set $\mathcal{C}$ ). Almost always we will assume that the functions are differentiable. The set of feasible solutions will be denoted by $\mathcal{F}$, hence

$$
\mathcal{F}=\left\{x \in \mathcal{C} \mid g_{j}(x) \leq 0, \quad j=1, \cdots, m\right\} .
$$

Definition 2.10 The vector $s \in \mathbb{R}^{n}$ is called a feasible direction at a point $x \in \mathcal{F}$ if there is a $\lambda_{0}>0$ such that $x+\lambda s \in \mathcal{F}$ for all $0 \leq \lambda \leq \lambda_{0}$. The set of feasible directions at the feasible point $x \in \mathcal{F}$ is denoted by $\mathcal{F D}(x)$.

Example 2.11 Assume that the feasible set $\mathcal{F} \subset \mathbb{R}^{2}$ is defined by the three constraints

$$
-x_{1}-x_{2}+1 \leq 0,1-x_{2} \leq 0, x_{1}-x_{2} \leq 0
$$

If $\bar{x}=(1,1)$, then the set of feasible directions at $\bar{x}$ is $\mathcal{F} \mathcal{D}(\bar{x})=\left\{s \in \mathbb{R}^{2} \mid s_{2} \geq s_{1}, s_{2} \geq 0\right\}$. Note that in this case $\mathcal{F D}(\bar{x})$ is a closed convex set.


Example 2.12 Assume that the feasible set $\mathcal{F} \subset \mathbb{R}^{2}$ is defined by the single constraint $x_{1}^{2}-x_{2} \leq 0$.
If $\bar{x}=(1,1)$, then the set of feasible directions at $\bar{x}$ is $\mathcal{F} \mathcal{D}(\bar{x})=\left\{s \in \mathbb{R}^{2} \mid s_{2}>2 s_{1}\right\}$. Observe that now $\mathcal{F} \mathcal{D}(\bar{x})$ is an open set.


Lemma 2.13 For any $x \in \mathcal{F}$ the set of feasible directions $\mathcal{F} \mathcal{D}(x)$ is a convex cone.

Proof: Let $\vartheta>0$. Obviously, $s \in \mathcal{F} \mathcal{D}(x)$ implies $(\vartheta s) \in \mathcal{F D}(x)$ since $x+\frac{\lambda}{\vartheta}(\vartheta s)=$ $x+\lambda s \in \mathcal{F}$, hence $\mathcal{F D}(x)$ is a cone. To prove the convexity of $\mathcal{F} \mathcal{D}(x)$ let us take $s, \bar{s} \in \mathcal{F} \mathcal{D}(x)$. Then by definition we have $x+\lambda s \in \mathcal{F}$ and $x+\lambda \bar{s} \in \mathcal{F}$ for some $\lambda>0$ (observe that a common $\lambda$ can be taken). Further, for $0 \leq \alpha \leq 1$ we write

$$
x+\lambda(\alpha s+(1-\alpha) \bar{s})=\alpha(x+\lambda s)+(1-\alpha)(x+\lambda \bar{s}) .
$$

Due to the convexity of $\mathcal{F}$ the right hand side of the above equation is in $\mathcal{F}$, hence the convexity of $\mathcal{F} \mathcal{D}(x)$ follows.

In view of the above lemma we may speak about the cone of feasible directions $\mathcal{F} \mathcal{D}(x)$ for any $x \in \mathcal{F}$. Note that the cone of feasible directions is not necessarily closed even if the set of feasible solutions $\mathcal{F}$ is closed. Figure 2.1 illustrates the cone $\mathcal{F} \mathcal{D}(x)$ for three different choices of $\mathcal{F}$ and $x$.

We will now formulate an optimality condition in terms of the cone of feasible directions. It states that a feasible solution is optimal if and only if the gradient of the objective in that point has an acute angle with all feasible directions at that point (no feasible descent direction exists).

Theorem 2.14 The feasible point $\bar{x} \in \mathcal{F}$ is an optimal solution of the convex optimization problem (CO) if and only if for all $s \in \mathcal{F} \mathcal{D}(\bar{x})$ one has $\delta f(\bar{x}, s) \geq 0$.

Proof: Observing that $s \in \mathcal{F} \mathcal{D}(\bar{x})$ if and only if $s=\lambda(x-\bar{x})$ for some $x \in \mathcal{F}$ and some $\lambda>0$, the result follows in the same way as in the proof of Theorem 2.9.

$x \in \mathcal{F}$




Figure 2.1: Convex feasible sets and cones of feasible directions.

### 2.2.1 A geometric interpretation

The purpose of this section is to give a geometric interpretation of the result of Theorem 2.14. In doing so, we will look at where we are going in the rest of this chapter, and what we would like to prove. The results in this section are not essential or even necessary in developing the theory further, but should provide some geometric insight and a taste of things to come.
Theorem 2.14 gives us necessary and sufficient optimality conditions, but is not a practical test because we do not have a description of the cone $\mathcal{F D}(\bar{x})$ and therefore cannot perform the test:

$$
\begin{equation*}
\text { 'is } \delta f(\bar{x}, s) \equiv \nabla f(\bar{x})^{T} s \geq 0 \text { for all } s \in \mathcal{F} \mathcal{D}(\bar{x}) ? \text { ' } \tag{2.4}
\end{equation*}
$$

It is easy to give a sufficient condition for (2.4) to hold, which we will now do. This condition will depend only on the constraint functions that are zero (active) at $\bar{x}$.

Definition 2.15 $A$ constraint $g_{i}(x) \leq 0$ is called active at $\bar{x} \in \mathcal{F}$ if $g_{i}(\bar{x})=0$.

Now let $I_{\bar{x}}$ denote the index set of the active constraints at $\bar{x}$, and assume that $\mathcal{C}=\mathbb{R}^{n}$.
We now give a sufficient condition for (2.4) to hold (i.e. for $\bar{x}$ to be an optimal solution).

Theorem 2.16 A point $\bar{x} \in \mathcal{F}$ is an optimal solution of problem (CO) (if $\mathcal{C}=\mathbb{R}^{n}$ ) if

$$
\begin{equation*}
\nabla f(\bar{x})=-\sum_{i \in I_{\bar{x}}} \bar{y}_{i} \nabla g_{i}(\bar{x}) \tag{2.5}
\end{equation*}
$$

for some nonnegative vector $\bar{y}$, where $I_{\bar{x}}$ denotes the index set of the active constraints at $\bar{x}$, as before.

The condition (2.5) is called the Karush-Kuhn-Tucker (KKT) optimality condition,. One can check whether it holds for a given $\bar{x} \in \mathcal{F}$ by using techniques from linear optimization.

The proof that (2.5) is indeed a sufficient condition for optimality follows from the next two exercises.

Exercise 2.5 Let $s \in \mathcal{F D}(\bar{x})$ be a given feasible direction at $\bar{x} \in \mathcal{F}$ for (CO) and let $\mathcal{C}=\mathbb{R}^{n}$. One has

$$
\nabla g_{i}(\bar{x})^{T} s \leq 0 \text { for all } i \in I_{\bar{x}}
$$

(Hint: Use Lemma 1.49.)
Exercise 2.6 Let $\bar{x} \in \mathcal{F}$ be a feasible solution of (CO) where $\mathcal{C}=\mathbb{R}^{n}$. Use the previous exercise and Theorem 2.14 to show that, if there exists a $\bar{y} \geq 0$ such that

$$
\nabla f(\bar{x})=-\sum_{i \in I_{\bar{x}}} \bar{y}_{i} \nabla g_{i}(\bar{x}),
$$

then $\bar{x}$ is an optimal solution of (CO).
Exercise 2.7 We wish to design a cylindrical can with height $h$ and radius $r$ such that the volume is at least $V$ units and the total surface area is minimal.

We can formulate this as the following optimization problem:

$$
p^{*}:=\min 2 \pi r^{2}+2 \pi r h
$$

subject to

$$
\pi r^{2} h \geq V, r>0, h>0
$$

1. Show that we can rewrite the above problem as the following optimization problem:

$$
p^{*}=\min 2 \pi\left(e^{2 x_{1}}+e^{x_{1}+x_{2}}\right),
$$

subject to

$$
\ln \left(\frac{V}{\pi}\right)-2 x_{1}-x_{2} \leq 0, x_{1} \in \mathbb{R}, x_{2} \in \mathbb{R}
$$

2. Prove that the new problem is a convex optimization problem (CO).
3. Prove that the optimal design is where $r=\frac{1}{2} h=\left(\frac{V}{2 \pi}\right)^{\frac{1}{3}}$ by using the result of Exercise 2.6.

The KKT condition (2.5) is sufficient for optimality, but is not a necessary condition for optimality in general, as the next example shows.

Example 2.17 Consider the problem of the form (CO):

$$
\min x \text { subject to } x^{2} \leq 0, x \in \mathbb{R} \text {. }
$$

Obviously, the unique optimal solution is $\bar{x}=0$, and the constraint $g(x):=x^{2} \leq 0$ is active at $\bar{x}$.
If we write out condition (2.5), we get

$$
1 \equiv \nabla f(\bar{x})=-\bar{y} \nabla g(\bar{x}) \equiv-\bar{y}(2(0))=0,
$$

which is obviously not satisfied for any choice of $\bar{y} \geq 0$. In other words, we cannot prove that $\bar{x}=0$ is an optimal solution by using the KKT condition.

In the rest of the chapter we will show that the KKT conditions are also necessary optimality conditions for (CO), if the feasible set $\mathcal{F}$ satisfies an additional assumption called the Slater condition.

### 2.2.2 The Slater condition

We still consider the convex optimization (CO) problem in the form:

$$
\begin{array}{rll}
(C O) & \min & f(x) \\
& \text { s.t. } & g_{j}(x) \leq 0, \quad j=1, \cdots, m \\
& x \in \mathcal{C},
\end{array}
$$

where $\mathcal{C} \subseteq \mathbb{R}^{n}$ is a convex set and $f, g_{1}, \cdots, g_{m}$ are convex functions on $\mathcal{C}$ (or on an open set that contains the set $\mathcal{C}$ ). Almost always we will assume that the functions $f$ and $g_{j}$ are differentiable. The set of indices $\{1, \cdots, m\}$ is denoted by $J$, and the set of feasible solutions by $\mathcal{F}$, hence

$$
\mathcal{F}=\left\{x \in \mathcal{C} \mid g_{j}(x) \leq 0, \quad j \in J\right\} .
$$

We now introduce the assumption on $\mathcal{F}$ that we referred to in the previous section, namely the Slater condition.

Definition 2.18 A vector (point) $x^{0} \in \mathcal{C}^{0}$ is called a Slater point of (CO) if

$$
\begin{array}{ll}
g_{j}\left(x^{0}\right)<0, & \text { for all } j \text { where } \\
g_{j} \text { is nonlinear, } \\
g_{j}\left(x^{0}\right) \leq 0, & \text { for all } j \text { where } \quad g_{j} \text { is linear. }
\end{array}
$$

If a Slater point exists we say that (CO) is Slater regular or (CO) satisfies the Slater condition, or (CO) satisfies the Slater constraint qualification.

## Example 2.19

1. Let us consider the optimization problem

$$
\begin{array}{cc}
\min & f(x) \\
\text { s.t. } & x_{1}^{2}+x_{2}^{2} \leq 4 \\
& x_{1}-x_{2} \geq 2 \\
& x_{2} \geq 0 \\
& \mathcal{C}=\mathbb{R}^{2} .
\end{array}
$$

The feasible region $\mathcal{F}$ contains only one point, (2, 0), for which the non-linear constraint becomes an equality. Hence, the problem is not Slater regular.

2. Let us consider the optimization problem

$$
\begin{array}{cc}
\min & f(x) \\
\text { s.t. } & x_{1}^{2}+x_{2}^{2} \leq 4 \\
& x_{1}-x_{2} \geq 2 \\
& x_{2} \geq-1 \\
& \mathcal{C}=\left\{x \mid x_{1} \leq 1\right\} .
\end{array}
$$

Again the feasible region contains only one point, $(1,-1)$. For this point the non-linear constraint holds with strict inequality. However, this point does not lie in the relative interior of $\mathcal{C}$. Hence, the problem is not Slater regular.


Exercise 2.8 Assume that (CO) satisfies the Slater condition. Prove that any $x \in \mathcal{F}^{0}$ is a Slater point of (CO).

Exercise 2.9 By solving a so-called first-phase problem one can check whether a given problem of the form ( $C O$ ) satisfies the Slater condition. Let us assume that $\mathcal{C}=\mathbb{R}^{n}$ and consider the first-phase problem

$$
\begin{array}{ll}
\min & \tau \\
\text { s.t. } & g_{j}(x)-\tau \leq 0, \quad j=1, \cdots, m \\
& x \in \mathbb{R}^{n}, \tau \in \mathbb{R},
\end{array}
$$

where $\tau$ is an auxiliary variable.
(a) Show that the first-phase problem is Slater regular.
(b) What information can you gain about problem (CO) by looking at the optimal objective value $\tau^{*}$ of the first-phase problem? (Consider the cases: $\tau^{*}>0, \tau^{*}=0$ and $\tau^{*}<0$.) $\triangleleft$

We can further refine our definition. Some constraint functions $g_{j}(x)$ might take the value zero for all feasible points. Such constraints are called singular while the others are called regular. Hence the index set of singular constraints is defined as

$$
J_{s}=\left\{j \in J \mid g_{j}(x)=0 \text { for all } x \in \mathcal{F}\right\},
$$

while the index set of regular (qualified) constraints is defined as the complement of the singular set

$$
J_{r}=J \backslash J_{s}=\left\{j \in J \mid g_{j}(x)<0 \text { for some } x \in \mathcal{F}\right\} .
$$

Remark: Note, that if (CO) is Slater regular, then all singular functions must be linear.

Definition 2.20 A Slater point $x^{*} \in \mathcal{C}^{0}$ is called an Ideal Slater point of the convex optimization problem (CO) if

$$
\begin{array}{ll}
g_{j}\left(x^{*}\right)<0 & \text { for all } j \in J_{r}, \\
g_{j}\left(x^{*}\right)=0 & \text { for all } j \in J_{s} .
\end{array}
$$

First we show an elementary property.
Lemma 2.21 If the convex optimization problem (CO) is Slater regular then there exists an ideal Slater point $x^{*} \in \mathcal{F}$.

Proof: According to the assumption, there exists a Slater point $x^{0} \in \mathcal{C}^{0}$ and there exist points $x^{k} \in \mathcal{F}$ for all $k \in J_{r}$ such that $g_{k}\left(x^{k}\right)<0$. Let $\lambda_{0}>0, \lambda_{k}>0$ for all $k \in J_{r}$ such that $\lambda_{0}+\sum_{k \in J_{r}} \lambda_{k}=1$, then $x^{*}=\lambda_{0} x^{0}+\sum_{j \in J_{r}} \lambda_{k} x^{k}$ is an ideal Slater point. This last statement follows from the convexity of the functions $g_{j}$.

## Example 2.22

1. Let us consider the optimization problem

$$
\begin{array}{cc}
\text { min } & f(x) \\
\text { s.t. } & x_{1}^{2}+x_{2}^{2} \leq 4 \\
& x_{1}-x_{2} \geq 2 \\
& x_{2} \geq-1 \\
& \mathcal{C}=\left\{x \mid x_{1}=1\right\} .
\end{array}
$$

The feasible region contains only one point, $(1,-1)$, but now this point does lie in the relative interior of the convex set $\mathcal{C}$. Hence, this point is an ideal Slater point.

2. Let us consider the optimization problem

$$
\begin{array}{rc}
\min & f(x) \\
\text { s.t. } & x_{1}^{2}+x_{2}^{2} \leq 4 \\
& x_{1}-x_{2} \geq 2 \\
& x_{2} \geq-1 \\
& \mathcal{C}=\mathbb{R}^{2} .
\end{array}
$$

Now, the point $(1,-1)$ is again a Slater point, but not an ideal Slater point. The point $\left(\frac{3}{2},-\frac{3}{4}\right)$ is an ideal Slater point.


Exercise 2.10 Prove that any ideal Slater point of (CO) is in the relative interior of $\mathcal{F}$. $\triangleleft$

### 2.2.3 Convex Farkas lemma

The convex Farkas lemma is an example of a theorem of alternatives, which means that it is a statement of the type: for two specific systems of inequalities (I) and (II), (I) has a solution if and only if (II) has no solution. It will play an essential role in developing the KKT theory.
Before stating the convex Farkas lemma, we present a simple separation theorem. It essentially states that disjoint convex sets can be separated by a (hyper)plane (geometrically, in $\mathbb{R}^{2}$ or $\mathbb{R}^{3}$, the convex sets lie on different sides of the separating plane). Its proof can be found in most textbooks (see e.g. [2]).

Theorem 2.23 Let $\mathcal{U} \subseteq \mathbb{R}^{n}$ be a convex set and a point $w \in \mathbb{R}^{n}$ with $w \notin \mathcal{U}$ be given. Then there is a separating hyperplane $\left\{x \mid a^{T} x=\alpha\right\}$, with $a \in \mathbb{R}^{n}, \alpha \in \mathbb{R}$ such that

1. $a^{T} w \leq \alpha$;
2. $a^{T} u \geq \alpha$ for all $u \in \mathcal{U}$ but $\mathcal{U}$ is not a subset of the hyperplane.

Note that the last property says that there is a $\bar{u} \in \mathcal{U}$ such that $a^{T} \bar{u}>\alpha$.
Now we are ready to prove the convex Farkas Lemma. The proof here is a simplified version of the proofs in the books [38, 42].

Lemma 2.24 (Farkas) The convex optimization problem (CO) is given and we assume that the Slater regularity condition is satisfied. The inequality system

$$
\begin{align*}
& f(x)<0 \\
& g_{j}(x) \leq 0, \quad j=1, \cdots, m  \tag{2.6}\\
& x \in \mathcal{C}
\end{align*}
$$

has no solution if and only if there exists a vector $y=\left(y_{1}, \cdots, y_{m}\right) \geq 0$ such that

$$
\begin{equation*}
f(x)+\sum_{j=1}^{m} y_{j} g_{j}(x) \geq 0 \quad \text { for all } x \in \mathcal{C} \tag{2.7}
\end{equation*}
$$

Before proving this important result we make remark. The systems (2.6) and (2.7) are called alternative systems, i.e. exactly one of them has a solution.
Proof: If the system (2.6) has a solution then clearly (2.7) cannot be true for that solution. This is the trivial part of the lemma. Note that this part is true without any regularity condition.

To prove the other side let us assume that (2.6) has no solution. With $u=\left(u_{0}, \cdots, u_{m}\right)$, we define the set $\mathcal{U} \in \mathbb{R}^{m+1}$ as follows.

$$
\mathcal{U}=\left\{u \mid \exists x \in \mathcal{C} \text { with } u_{0}>f(x), u_{j} \geq g_{j}(x) \text { if } j \in J_{r}, u_{j}=g_{j}(x) \text { if } j \in J_{s}\right\} .
$$

Clearly the set $\mathcal{U}$ is convex (note that due to the Slater condition singular functions are linear) and due to the infeasibility of (2.6) it does not contain the origin. Hence according to Theorem 2.23 there exists a separating hyperplane defined by an appropriate vector ( $y_{0}, y_{1}, \cdots, y_{m}$ ) and $\alpha=0$ such that

$$
\begin{equation*}
\sum_{j=0}^{m} y_{j} u_{j} \geq 0 \quad \text { for all } u \in \mathcal{U} \tag{2.8}
\end{equation*}
$$

and for some $\bar{u} \in \mathcal{U}$ one has

$$
\begin{equation*}
\sum_{j=0}^{m} y_{j} \bar{u}_{j}>0 \tag{2.9}
\end{equation*}
$$

The rest of the proof is divided into four parts.
I. First we prove that $y_{0} \geq 0$ and $y_{j} \geq 0$ for all $j \in J_{r}$.
II. Secondly we establish that (2.8) holds for $u=\left(f(x), g_{1}(x), \cdots, g_{m}(x)\right)$ if $x \in \mathcal{C}$.
III. Then we prove that $y_{0}$ must be positive.
IV. Finally, it is shown by using induction that we can assume $y_{j}>0$ for all $j \in J_{s}$.
I. First we show that $y_{0} \geq 0$ and $y_{j} \geq 0$ for all $j \in J_{r}$. Let us assume that $y_{0}<0$. Let us take an arbitrary $\left(u_{0}, u_{1}, \cdots, u_{m}\right) \in \mathcal{U}$. By definition $\left(u_{0}+\lambda, u_{1}, \cdots, u_{m}\right) \in \mathcal{U}$ for all $\lambda \geq 0$. Hence by (2.8) one has

$$
\lambda y_{0}+\sum_{j=0}^{m} y_{j} u_{j} \geq 0 \quad \text { for all } \lambda \geq 0
$$

For sufficiently large $\lambda$ the left hand side is negative, which is a contradiction, i.e. $y_{0}$ must be nonnegative. The proof of the nonnegativity of all $y_{j}$ as $j \in J_{r}$ goes analogously.
II. Secondly we establish that

$$
\begin{equation*}
y_{0} f(x)+\sum_{j=1}^{m} y_{j} g_{j}(x) \geq 0 \quad \text { for all } x \in \mathcal{C} \tag{2.10}
\end{equation*}
$$

This follows from the observation that for all $x \in \mathcal{C}$ and for all $\lambda>0$ one has $u=$ $\left(f(x)+\lambda, g_{1}(x), \cdots, g_{m}(x)\right) \in \mathcal{U}$, thus

$$
y_{0}(f(x)+\lambda)+\sum_{j=1}^{m} y_{j} g_{j}(x) \geq 0 \quad \text { for all } x \in \mathcal{C}
$$

Taking the limit as $\lambda \longrightarrow 0$ the claim follows.
III. Thirdly we show that $y_{0}>0$. The proof is by contradiction. We already know that $y_{0} \geq 0$. Let us assume to the contrary that $y_{0}=0$. Hence from (2.10) we have

$$
\sum_{j \in J_{r}} y_{j} g_{j}(x)+\sum_{j \in J_{s}} y_{j} g_{j}(x)=\sum_{j=1}^{m} y_{j} g_{j}(x) \geq 0 \quad \text { for all } x \in \mathcal{C}
$$

Taking an ideal Slater point $x^{*} \in \mathcal{C}^{0}$ one has

$$
g_{j}\left(x^{*}\right)=0 \text { if } j \in J_{s},
$$

whence

$$
\sum_{j \in J_{r}} y_{j} g_{j}\left(x^{*}\right) \geq 0
$$

Since $y_{j} \geq 0$ and $g_{j}\left(x^{*}\right)<0$ for all $j \in J_{r}$, this implies $y_{j}=0$ for all $j \in J_{r}$. This results in

$$
\begin{equation*}
\sum_{j \in J_{s}} y_{j} g_{j}(x) \geq 0 \quad \text { for all } x \in \mathcal{C} \tag{2.11}
\end{equation*}
$$

Now, from (2.9), with $\bar{x} \in \mathcal{C}$ such that $\bar{u}_{j}=g_{j}(\bar{x})$ if $i \in J_{s}$ we have

$$
\begin{equation*}
\sum_{j \in J_{s}} y_{j} g_{j}(\bar{x})>0 \tag{2.12}
\end{equation*}
$$

Because the ideal Slater point $x^{*}$ is in the relative interior of $\mathcal{C}$ there exist a vector $\tilde{x} \in \mathcal{C}$ and $0<\lambda<1$ such that $x^{*}=\lambda \bar{x}+(1-\lambda) \tilde{x}$. Using that $g_{j}\left(x^{*}\right)=0$ for $j \in J_{s}$ and that the singular functions are linear one gets

$$
\begin{aligned}
0 & =\sum_{j \in J_{s}} y_{j} g_{j}\left(x^{*}\right) \\
& =\sum_{j \in J_{s}} y_{j} g_{j}(\lambda \bar{x}+(1-\lambda) \tilde{x}) \\
& =\lambda \sum_{j \in J_{s}} y_{j} g_{j}(\bar{x})+(1-\lambda) \sum_{j \in J_{s}} y_{j} g_{j}(\tilde{x}) \\
& >(1-\lambda) \sum_{j \in J_{s}} y_{j} g_{j}(\tilde{x}) .
\end{aligned}
$$

Here the last inequality follows from (2.12). The inequality

$$
(1-\lambda) \sum_{j \in J_{s}} y_{j} g_{j}(\tilde{x})<0
$$

contradicts (2.11). Hence we have proved that $y_{0}>0$.
At this point we have (2.10) with $y_{0}>0$ and $y_{j} \geq 0$ for all $j \in J_{r}$. Dividing by $y_{0}>0$ in (2.10) and by defining $y_{j}:=\frac{y_{j}}{y_{0}}$ for all $j \in J$ we obtain

$$
\begin{equation*}
f(x)+\sum_{j=1}^{m} y_{j} g_{j}(x) \geq 0 \quad \text { for all } x \in \mathcal{C} \tag{2.13}
\end{equation*}
$$

We finally show that $y$ may be taken such that $y_{j}>0$ for all $j \in J_{s}$.
IV. To complete the proof we show by induction on the cardinality of $J_{s}$ that one can make $y_{j}$ positive for all $j \in J_{s}$. Observe that if $J_{s}=\emptyset$ then we are done. If $\left|J_{s}\right|=1$ then we apply the results proved till this point to the inequality system

$$
\begin{align*}
& g_{s}(x)<0, \\
& g_{j}(x) \leq 0, \quad j \in J_{r},  \tag{2.14}\\
& x \in \mathcal{C}
\end{align*}
$$

where $\{s\}=J_{s}$. The system (2.14) has no solution, it satisfies the Slater condition, and therefore there exists a $\hat{y} \in \mathbb{R}^{m-1}$ such that

$$
\begin{equation*}
g_{s}(x)+\sum_{j \in J_{r}} \hat{y}_{j} g_{j}(x) \geq 0 \quad \text { for all } x \in \mathcal{C}, \tag{2.15}
\end{equation*}
$$

where $\hat{y}_{j} \geq 0$ for all $j \in J_{r}$. Adding a sufficiently large positive multiple of (2.15) to (2.13) one obtains a positive coefficient $\hat{y}_{s}>0$ for $g_{s}(x)$.

The general inductive step goes analogously. Assuming that the result is proved if $\left|J_{s}\right|=k$ then the result is proved for the case $\left|J_{s}\right|=k+1$. Let $s \in J_{s}$ then $\left|J_{s} \backslash\{s\}\right|=k$,
and hence the inductive assumption applies to the system

$$
\begin{align*}
& g_{s}(x)<0 \\
& g_{j}(x) \leq 0, \quad j \in J_{s} \backslash\{s\},  \tag{2.16}\\
& g_{j}(x) \leq 0, \quad j \in J_{r}, \\
& x \in \mathcal{C} .
\end{align*}
$$

By construction the system (2.16) has no solution, it satisfies the Slater condition, and by the inductive assumption we have a $\hat{y} \in \mathbb{R}^{m-1}$ such that

$$
\begin{equation*}
g_{s}(x)+\sum_{j \in J_{r} \cup J_{s} \backslash\{s\}} \hat{y}_{j} g_{j}(x) \geq 0 \quad \text { for all } x \in \mathcal{C} . \tag{2.17}
\end{equation*}
$$

where $\hat{y}_{j}>0$ for all $j \in J_{s} \backslash\{s\}$ and $\hat{y}_{j} \geq 0$ for all $j \in J_{r}$. Adding a sufficiently large multiple of (2.17) to (2.13), one obtains the desired nonnegative multipliers.

Remark: Note, that finally we proved slightly more than was stated. We have proved that the multipliers of all the singular constraints can be made strictly positive.

Example 2.25 [Farkas Lemma]

1. Let us consider the convex optimization problem

$$
\begin{array}{lll}
\text { (CO) } & \min & x \\
& \text { s.t. } & x^{2} \leq 0 \\
& x \in \mathbb{R} .
\end{array}
$$

Then (CO) is not Slater regular.
The system

$$
\begin{aligned}
x & <0 \\
x^{2} & \leq 0
\end{aligned}
$$

has no solution, but for every $y>0$ the quadratic function $f(x)=x+y x^{2}$ has two zeroes.


So, there is no $y \geq 0$ such that

$$
x+y x^{2} \geq 0 \text { for all } x \in \mathbb{R} .
$$

Hence, the Farkas Lemma does not hold for (CO).
2. Let us consider the convex optimization problem

$$
\begin{array}{lll}
(\mathrm{CO}) & \min & 1+x \\
& \text { s.t. } & x^{2}-1 \leq 0 \\
& x \in \mathbb{R} .
\end{array}
$$

Then (CO) is Slater regular (0 is an ideal Slater point). The system

$$
\begin{array}{r}
1+x<0 \\
x^{2}-1 \leq 0
\end{array}
$$

has no solution. If we let $y=\frac{1}{2}$ the quadratic function

$$
g(x)=x+1+y\left(x^{2}-1\right)=\frac{1}{2} x^{2}+x+\frac{1}{2}
$$


has only one zero, thus one has

$$
\frac{1}{2} x^{2}+x+\frac{1}{2} \geq 0 \text { for all } x \in \mathbb{R}
$$

Exercise 2.11 Let the matrices $A: m \times n$ and the vector $b \in \mathbb{R}^{m}$ be given. Apply the convex Farkas Lemma 2.24 to prove that exactly one of the following alternative systems (I) or (II) is solvable:

$$
(I) \quad A x \leq b, \quad x \geq 0
$$

or

$$
(I I) \quad A^{T} y \geq 0, \quad y \geq 0, \quad b^{T} y<0
$$

Exercise 2.12 Let the matrices $A: m \times n, B: k \times n$ and the vectors $a \in \mathbb{R}^{m}, b \in \mathbb{R}^{k}$ be given. With a proper reformulation, apply the convex Farkas Lemma 2.24 to the inequality system

$$
A x \leq a, \quad B x<b, \quad x \geq 0
$$

Exercise 2.13 Let the matrix $A: m \times n$ and the vectors $c \in \mathbb{R}^{n}$ and $b \in R^{m}$ be given. Apply the convex Farkas Lemma 2.24 to prove the so-called Goldman-Tucker theorem for the LO problem:

$$
\min \left\{c^{T} x: A x=b, \quad x \geq 0\right\}
$$

when it admits an optimal solution. In other words, prove that there exists an optimal solution $x^{*}$ and an optimal solution $\left(y^{*}, s^{*}\right)$ of the dual LO problem

$$
\max \left\{b^{T} y: A^{T} y+s=c, \quad s \geq 0\right\}
$$

such that

$$
x^{*}+s^{*}>0
$$

### 2.2.4 Karush-Kuhn-Tucker theory

For the convex optimization problem (CO) one defines the Lagrange function (or $L a$ grangian)

$$
\begin{equation*}
L(x, y):=f(x)+\sum_{j=1}^{m} y_{j} g_{j}(x) \tag{2.18}
\end{equation*}
$$

where $x \in \mathcal{C}$ and $y \geq 0$. Note that for fixed $y$ the Lagrange function is convex in $x$.

Definition 2.26 A vector pair $(\bar{x}, \bar{y}) \in \mathbb{R}^{n+m}, \bar{x} \in \mathcal{C}$ and $\bar{y} \geq 0$ is called a saddle point of the Lagrange function $L$ if

$$
\begin{equation*}
L(\bar{x}, y) \leq L(\bar{x}, \bar{y}) \leq L(x, \bar{y}) \tag{2.19}
\end{equation*}
$$

for all $x \in \mathcal{C}$ and $y \geq 0$.

One easily sees that (2.19) is equivalent with

$$
L(\bar{x}, y) \leq L(x, \bar{y}) \quad \text { for all } x \in \mathcal{C}, \quad y \geq 0
$$

We will see (in the proof of Theorem 2.30) that the $\bar{x}$ part of a saddle point is always an optimal solution of $(\mathrm{CO})$.

Example 2.27 [Saddle point] Let us consider the convex optimization problem

$$
\begin{array}{lll}
(\mathrm{CO}) & \min & -x+2 \\
& \text { s.t. } & e^{x}-4 \leq 0 \\
& x \in \mathbb{R}
\end{array}
$$

Then the Lagrange function of $(\mathrm{CO})$ is given by

$$
L(x, y)=-x+2+y\left(e^{x}-4\right)
$$

where the Lagrange multiplier $y$ is non-negative. For fixed $y>0$ we have

$$
\frac{\partial}{\partial x} L(x, y)=-1+y e^{x}=0
$$

for $x=-\log y$, thus $L(-\log y, y)=\log y-4 y+3$ is a minimum.
On the other hand, for feasible $x$, i.e. if $x \leq \log 4$, we have

$$
\sup _{y \geq 0} y\left(e^{x}-4\right)=0
$$

Hence, defining $\psi(y)=\inf _{x \in \mathbb{R}} L(x, y)$ and $\phi(x)=\sup _{y \geq 0} L(x, y)$ we have

$$
\begin{aligned}
& \psi(y)= \begin{cases}\log y-4 y+3 & \text { for } y>0 \\
-\infty & \text { for } y=0\end{cases} \\
& \phi(x)= \begin{cases}-x+2 & \text { for } x \leq \log 4 \\
\infty & \text { for } x>\log 4\end{cases}
\end{aligned}
$$

Now, we have

$$
\frac{d}{d y} \psi(y)=\frac{1}{y}-4=0
$$

for $y=\frac{1}{4}$, i.e. this value gives the maximum of $\psi(y)$. Hence, $\sup _{y \geq 0} \psi(y)=-\log 4+2$. The function $\phi(x)$ is minimal for $x=\log 4$, thus $\inf _{x \in \mathbb{R}} \phi(x)=-\log 4+2$ and we conclude that $\left(\log 4, \frac{1}{4}\right)$ is a saddle point of the Lagrange function $L(x, y)$. Note that $x=\log 4$ is the optimal solution of (CO).

Lemma 2.28 $A$ saddle point $(\bar{x}, \bar{y}) \in \mathbb{R}^{n+m}, \bar{x} \in \mathcal{C}$ and $\bar{y} \geq 0$ of $L(x, y)$ satisfies the relation

$$
\begin{equation*}
\inf _{x \in \mathcal{C}} \sup _{y \geq 0} L(x, y)=L(\bar{x}, \bar{y})=\sup _{y \geq 0} \inf _{x \in \mathcal{C}} L(x, y) \tag{2.20}
\end{equation*}
$$

Proof: For any $(\hat{x}, \hat{y})$ one has

$$
\inf _{x \in \mathcal{C}} L(x, \hat{y}) \leq L(\hat{x}, \hat{y}) \leq \sup _{y \geq 0} L(\hat{x}, y)
$$

hence one can take the supremum of the left hand side and the infimum of the right hand side resulting in

$$
\begin{equation*}
\sup _{y \geq 0} \inf _{x \in \mathcal{C}} L(x, y) \leq \inf _{x \in \mathcal{C}} \sup _{y \geq 0} L(x, y) \tag{2.21}
\end{equation*}
$$

Further using the saddle point inequality (2.19) one obtains

$$
\begin{equation*}
\inf _{x \in \mathcal{C}} \sup _{y \geq 0} L(x, y) \leq \sup _{y \geq 0} L(\bar{x}, y) \leq L(\bar{x}, \bar{y}) \leq \inf _{x \in \mathcal{C}} L(x, \bar{y}) \leq \sup _{y \geq 0} \inf _{x \in \mathcal{C}} L(x, y) \tag{2.22}
\end{equation*}
$$

Combining (2.22) and (2.21) the equality (2.20) follows.

Condition (2.20) is a property of saddle points. If some $\bar{x} \in \mathcal{C}$ and $\bar{y} \geq 0$ satisfy (2.20), it does not imply that $(\bar{x}, \bar{y})$ is a saddle point though, as the following example shows.

Example 2.29 Let (CO) be given by

$$
\min e^{x} \text { subject to } x \leq 0
$$

Here

$$
L(x, y)=e^{x}+y x
$$

It is easy to verify that $\bar{x}=-1$ and $\bar{y}=e^{-1}$ satisfy (2.20). Indeed, $L(\bar{x}, \bar{y})=0$ and

$$
\inf _{x \in \mathcal{C}} \sup _{y \geq 0} L(x, y)=0
$$

by letting $x$ tend to $-\infty$. Likewise

$$
\sup _{y \geq 0} \inf _{x \in \mathcal{C}} L(x, y)=0
$$

However, $(\bar{x}, \bar{y})$ is not a saddle point of $L$. (This example does not have an optimal solution, and, as we have mentioned, the $\bar{x}$ part of a saddle point is always an optimal solution of (CO).)

We still do not know if a saddle point exists or not. Assuming Slater regularity, the next result states that $L(x, y)$ has a saddle point if and only if (CO) has an optimal solution.

Theorem 2.30 (Karush-Kuhn-Tucker) The convex optimization problem (CO) is given. Assume that the Slater regularity condition is satisfied. The vector $\bar{x}$ is an optimal solution of (CO) if and only if there is a vector $\bar{y}$ such that $(\bar{x}, \bar{y})$ is a saddle point of the Lagrange function $L$.

Proof: The easy part of the theorem is to prove that if $(\bar{x}, \bar{y})$ is a saddle point of $L(x, y)$ then $\bar{x}$ is optimal for (CO). The proof of this part does not need any regularity condition. From the saddle point inequality (2.19) one has

$$
f(\bar{x})+\sum_{j=1}^{m} y_{j} g_{j}(\bar{x}) \leq f(\bar{x})+\sum_{j=1}^{m} \bar{y}_{j} g_{j}(\bar{x}) \leq f(x)+\sum_{j=1}^{m} \bar{y}_{j} g_{j}(x)
$$

for all $y \geq 0$ and for all $x \in \mathcal{C}$. From the first inequality one easily derives $g_{j}(\bar{x}) \leq 0$ for all $j=1, \cdots, m$ hence $\bar{x} \in \mathcal{F}$ is a feasible solution of (CO). Taking the two extreme sides of the above inequality and substituting $y=0$ we have

$$
f(\bar{x}) \leq f(x)+\sum_{j=1}^{m} \bar{y}_{j} g_{j}(x) \leq f(x)
$$

for all $x \in \mathcal{F}$, i.e. $\bar{x}$ is optimal.
To prove the other direction we need Slater regularity and the Convex Farkas Lemma 2.24. Let us take an optimal solution $\bar{x}$ of the convex optimization problem (CO). Then the inequality system

$$
\begin{aligned}
& f(x)-f(\bar{x})<0 \\
& g_{j}(x) \leq 0, \quad j=1, \cdots, m \\
& x \in \mathcal{C}
\end{aligned}
$$

is infeasible. By the Convex Farkas Lemma 2.24 there exists a $\bar{y} \geq 0$ such that

$$
f(x)-f(\bar{x})+\sum_{j=1}^{m} \bar{y}_{j} g_{j}(x) \geq 0
$$

for all $x \in \mathcal{C}$. Using that $\bar{x}$ is feasible one easily derive the saddle point inequality

$$
f(\bar{x})+\sum_{j=1}^{m} y_{j} g_{j}(\bar{x}) \leq f(\bar{x}) \leq f(x)+\sum_{j=1}^{m} \bar{y}_{j} g_{j}(x)
$$

for all $y \geq 0$ and $x \in \mathcal{C}$, which completes the proof.

The following corollaries lead us to the Karush-Kuhn-Tucker (KKT) optimality conditions.

Corollary 2.31 Under the assumptions of Theorem 2.30 the vector $\bar{x} \in \mathcal{C}$ is an optimal solution of (CO) if and only if there exists $a \bar{y} \geq 0$ such that

$$
\begin{aligned}
& \text { (i) } f(\bar{x})=\min _{x \in \mathcal{C}}\left\{f(x)+\sum_{j=1}^{m} \bar{y}_{j} g_{j}(x)\right\} \quad \text { and } \\
& \text { (ii) } \sum_{j=1}^{m} \bar{y}_{j} g_{j}(\bar{x})=\max _{y \geq 0}\left\{\sum_{j=1}^{m} y_{j} g_{j}(\bar{x})\right\} .
\end{aligned}
$$

Proof: Easily follows from the theorem.

Corollary 2.32 Under the assumptions of Theorem 2.30 the vector $\bar{x} \in \mathcal{F}$ is an optimal solution of (CO) if and only if there exists $a \bar{y} \geq 0$ such that
(i) $f(\bar{x})=\min _{x \in \mathcal{C}}\left\{f(x)+\sum_{j=1}^{m} \bar{y}_{j} g_{j}(x)\right\} \quad$ and
(ii) $\sum_{j=1}^{m} \bar{y}_{j} g_{j}(\bar{x})=0$.

Proof: Easily follows from the Corollary 2.31.

Corollary 2.33 Let us assume that $\mathcal{C}=\mathbb{R}^{n}$ and the functions $f, g_{1}, \cdots, g_{m}$ are continuously differentiable functions. Under the assumptions of Theorem 2.30 the vector $\bar{x} \in \mathcal{F}$ is an optimal solution of (CO) if and only if there exists a $\bar{y} \geq 0$ such that
(i) $0=\nabla f(\bar{x})+\sum_{j=1}^{m} \bar{y}_{j} \nabla g_{j}(\bar{x}) \quad$ and
(ii) $\sum_{j=1}^{m} \bar{y}_{j} g_{j}(\bar{x})=0$.

Proof: Follows directly from Corollary 2.32 and the convexity of the function $f(x)+$ $\sum_{j=1}^{m} \bar{y}_{j} g_{j}(x), x \in \mathcal{C}$.

Exercise 2.14 Prove the above three Corollaries.
$\triangleleft$
Note that the last corollary stays valid if $\mathcal{C}$ is a full dimensional open subset of $\mathbb{R}^{n}$. If the set $\mathcal{C}$ is not full dimensional, then the right hand side vector, the $x$-gradient of the Lagrange function has to be orthogonal to any direction in the affine hull of $\mathcal{C}$ ( $c f$. Theorem 2.9). To check the validity of these statements is left to the reader.

Now the notion of Karush-Kuhn-Tucker (KKT) point is defined.
Definition 2.34 (KKT point) Let us assume that $\mathcal{C}=\mathbb{R}^{n}$ and the functions $f, g_{1}, \cdots, g_{m}$ are continuously differentiable functions. The vector $(\bar{x}, \bar{y}) \in \mathbb{R}^{n+m}$ is called a Karush-Kuhn-Tucker (KKT) point of (CO) if

$$
\begin{aligned}
& \text { (i) } g_{j}(\bar{x}) \leq 0, \text { for all } j \in J, \\
& \text { (ii) } 0=\nabla f(\bar{x})+\sum_{j=1}^{m} \bar{y}_{j} \nabla g_{j}(\bar{x}) \\
& \text { (iii) } \sum_{j=1}^{m} \bar{y}_{j} g_{j}(\bar{x})=0, \\
& \text { (iv) } \bar{y} \geq 0 .
\end{aligned}
$$

It is important to understand that - under the assumptions of Corollary 2.33 - $(\bar{x}, \bar{y})$ is a saddle point of the Lagrangian of (CO) if and only if it is a KKT point of (CO). The proof is left as an exercise.

Exercise 2.15 Let us assume that $\mathcal{C}=\mathbb{R}^{n}$ and the functions $f, g_{1}, \cdots, g_{m}$ are continuously differentiable convex functions and the assumptions of Theorem 2.30 hold. Show that $(\bar{x}, \bar{y})$ is a saddle point of the Lagrangian of (CO) if and only if it is a KKT point of (CO).

The so-called Karush-Kuhn-Tucker sufficient optimality conditions now follow from Corollary 2.33.

Corollary 2.35 Let us assume that $\mathcal{C}=\mathbb{R}^{n}$ and the functions $f, g_{1}, \cdots, g_{m}$ are continuously differentiable convex functions and the assumptions of Theorem 2.30 hold. Let the vector $(\bar{x}, \bar{y})$ be a KKT point, then $\bar{x}$ is an optimal solution of (CO).

Thus we have derived necessary and sufficient optimality conditions for the convex optimization problem (CO) under the Slater regularity assumption. Note that if an optimization problem is not convex, or does not satisfy any regularity condition, then only weaker results can be proven.

## Chapter 3

## Duality in convex optimization

Every optimization problem has an associated dual optimization problem. Under some assumptions, a convex optimization problem (CO) and its dual have the same optimal objective values. We can therefore use the dual problem to show that a certain solution of (CO) is in fact optimal. Moreover, some optimization algorithms solve (CO) and its dual problem at the same time, and when the objective values are the same then optimality has been proved. One can easily derive dual problems and duality results from the KKT theory or from the Convex Farkas Lemma. First we define the more general Lagrange dual and then we specialize it to get the so-called Wolfe dual for convex problems.

### 3.1 Lagrange dual

Definition 3.1 Denote $\psi(y)=\inf _{x \in \mathcal{C}}\left\{f(x)+\sum_{j=1}^{m} y_{j} g_{j}(x)\right\}$. The problem

$$
\begin{array}{r}
(L D) \quad \sup \psi(y) \\
y \geq 0
\end{array}
$$

is called the Lagrange dual of the convex optimization problem (CO).

Lemma 3.2 The Lagrange Dual (LD) of (CO) is a convex optimization problem, even if the functions $f, g_{1}, \cdots, g_{m}$ are not convex.

Proof: Because the maximization of $\psi(y)$ is equivalent to the minimization of $-\psi(y)$, we have only to prove that the function $-\psi(y)$ is convex, i.e. $\psi(y)$ is concave. Let $\bar{y}, \hat{y} \geq 0$ and $0 \leq \lambda \leq 1$. Using that the infimum of the sum of two functions is larger
than the sum of the two separate infimums one has:

$$
\begin{aligned}
\psi(\lambda \bar{y}+(1-\lambda) \hat{y}) & =\inf _{x \in \mathcal{C}}\left\{f(x)+\sum_{j=1}^{m}\left(\lambda \bar{y}_{j}+(1-\lambda) \hat{y}_{j}\right) g_{j}(x)\right\} \\
& =\inf _{x \in \mathcal{C}}\left\{\lambda\left[f(x)+\sum_{j=1}^{m} \bar{y}_{j} g_{j}(x)\right]+(1-\lambda)\left[f(x)+\sum_{j=1}^{m} \hat{y}_{j} g_{j}(x)\right]\right\} \\
& \geq \inf _{x \in \mathcal{C}}\left\{\lambda\left[f(x)+\sum_{j=1}^{m} \bar{y}_{j} g_{j}(x)\right]\right\}+\inf _{x \in \mathcal{C}}\left\{(1-\lambda)\left[f(x)+\sum_{j=1}^{m} \hat{y}_{j} g_{j}(x)\right]\right\} \\
& =\lambda \psi(\bar{y})+(1-\lambda) \psi(\hat{y}) .
\end{aligned}
$$

Definition 3.3 If $\bar{x}$ is a feasible solution of (CO) and $\bar{y} \geq 0$ then we call the quantity

$$
f(\bar{x})-\psi(\bar{y})
$$

the duality gap at $\bar{x}$ and $\bar{y}$.

It is easy to prove the so-called weak duality theorem, which states that the duality gap is always nonnegative.

Theorem 3.4 If $\bar{x}$ is a feasible solution of (CO) and $\bar{y} \geq 0$ then

$$
\psi(\bar{y}) \leq f(\bar{x})
$$

and equality holds if and only if $\inf _{x \in \mathcal{C}}\left\{f(x)+\sum_{j=1}^{m} \bar{y}_{j} g_{j}(x)\right\}=f(\bar{x})$.
Proof: By straightforward calculations one has

$$
\psi(\bar{y})=\inf _{x \in \mathcal{C}}\left\{f(x)+\sum_{j=1}^{m} \bar{y}_{j} g_{j}(x)\right\} \leq f(\bar{x})+\sum_{j=1}^{m} \bar{y}_{j} g_{j}(\bar{x}) \leq f(\bar{x}) .
$$

Equality holds if and only if $\inf _{x \in \mathcal{C}}\left\{f(x)+\sum_{j=1}^{m} \bar{y}_{j} g_{j}(x)\right\}=f(\bar{x})$ and hence $\bar{y}_{j} g_{j}(x)=0$ for all $j \in J$.

One easily derives the following corollary.
Corollary 3.5 If $\bar{x}$ is a feasible solution of (CO), $\bar{y} \geq 0$ and $\psi(\bar{y})=f(\bar{x})$ then the vector $\bar{x}$ is an optimal solution of (CO) and $\bar{y}$ is optimal for (LD). Further if the functions $f, g_{1}, \cdots, g_{m}$ are continuously differentiable then $(\bar{x}, \bar{y})$ is a KKT-point.

To prove the so-called strong duality theorem one needs a regularity condition.

Theorem 3.6 Let us assume that (CO) satisfies the Slater regularity condition. Let $\bar{x}$ be a feasible solution of (CO). The vector $\bar{x}$ is an optimal solution of (CO) if and only if there exists a $\bar{y} \geq 0$ such that $\bar{y}$ is an optimal solution of (LD) and

$$
\psi(\bar{y})=f(\bar{x}) .
$$

Proof: Directly follows from Corollary 2.31.

Exercise 3.1 Prove Theorem 3.6.

Remark: If the convex optimization problem does not satisfy a regularity condition, then it is not true in general that the duality gap is zero. It is also not always true (even not under regularity condition) that the convex optimization problem has an optimal solution. Frequently only the supremum or the infimum of the objective function exists.

Example 3.7 [Lagrange dual] Let us consider again the problem (see Example 2.25)

$$
\begin{array}{lll}
\text { (CO) } & \min & x \\
& \text { s.t. } & x^{2} \leq 0 \\
& x \in \mathbb{R} .
\end{array}
$$

As we have seen this (CO) problem is not Slater regular and the Convex Farkas Lemma 2.24 does not apply to the system

$$
\begin{aligned}
x & <0 \\
x^{2} & \leq 0
\end{aligned}
$$

On the other hand, we have

$$
\psi(y)=\inf _{x \in \mathbb{R}}\left(x+y x^{2}\right)=\left\{\begin{aligned}
-\frac{1}{4 y} & \text { for } y>0 \\
-\infty & \text { for } y=0
\end{aligned}\right.
$$

The Lagrange dual reads

$$
\sup _{y \geq 0} \psi(y) .
$$

The optimal value of the Lagrange dual is zero, i.e. in spite of the lack of Slater regularity there is no duality gap.

### 3.2 Wolfe dual

Observing the similarity of the formulation of the Lagrange dual (LD) and the conditions occurring in the corollaries of the KKT-Theorem 2.30 the so-called Wolfe dual is obtained.

Definition 3.8 Assume that $\mathcal{C}=\mathbb{R}^{n}$ and the functions $f, g_{1}, \cdots, g_{m}$ are continuously differentiable and convex. The problem

$$
\begin{gathered}
(W D) \sup _{x, y}\left\{f(x)+\sum_{j=1}^{m} y_{j} g_{j}(x)\right\} \\
\nabla f(x)+\sum_{j=1}^{m} y_{j} \nabla g_{j}(x)=0, \\
y \geq 0, x \in \mathbb{R}^{n}
\end{gathered}
$$

is called the Wolfe Dual of the convex optimization problem (CO).
Note that the variables in $(W D)$ are both $y \geq 0$ and $x \in \mathbb{R}^{n}$, and that the Lagrangian $L(x, y)$ is the objective function of $(W D)$. For this reason, the Wolfe dual does not have a concave objective function in general, but it is still very useful tool, as we will see. In particular, if the Lagrange function has a saddle point, $\mathcal{C}=\mathbb{R}^{n}$ and the functions $f, g_{1}, \cdots, g_{m}$ are continuously differentiable and convex, then the two dual problems are equivalent. Using the results of the previous section one easily proves weak and strong duality results, as we will now show. A more detailed discussion of duality theory can be found in $[2,28]$.

Theorem 3.9 (Weak duality for the Wolfe dual) Assume that $\mathcal{C}=\mathbb{R}^{n}$ and the functions $f, g_{1}, \cdots, g_{m}$ are continuously differentiable and convex. If $\hat{x}$ is a feasible solution of (CO) and $(\bar{x}, \bar{y})$ is a feasible solution for (WD) then

$$
L(\bar{x}, \bar{y}) \leq f(\hat{x})
$$

In other words, weak duality holds for (CO) and (WD).

Proof: Let $(\bar{x}, \bar{y})$ be a feasible solution for (WD). Since the functions $f$ and $g_{1}, \ldots, g_{m}$ are convex and continuously differentiable, and $\bar{y} \geq 0$, the function

$$
h(x):=f(x)+\sum_{j=1}^{m} \bar{y}_{j} g_{j}(x)
$$

must also be convex and continuously differentiable (see Lemma 1.40). Since $(\bar{x}, \bar{y})$ is feasible for $(W D)$, one has

$$
\nabla h(\bar{x})=\nabla f(\bar{x})+\sum_{j=1}^{m} \bar{y}_{j} \nabla g_{j}(\bar{x})=0 .
$$

This means that $\bar{x}$ is a minimizer of the function $h$, by Lemma 2.6. In other words

$$
\begin{equation*}
f(\bar{x})+\sum_{j=1}^{m} \bar{y}_{j} g_{j}(\bar{x}) \leq f(x)+\sum_{j=1}^{m} \bar{y}_{j} g_{j}(x) \quad \forall x \in \mathbb{R}^{n} \tag{3.1}
\end{equation*}
$$

Let $\hat{x}$ be an arbitrary feasible solution of (CO). Setting $x=\hat{x}$ in (3.1) one gets

$$
f(\bar{x})+\sum_{j=1}^{m} \bar{y}_{j} g_{j}(\bar{x}) \leq f(\hat{x})+\sum_{j=1}^{m} \bar{y}_{j} g_{j}(\hat{x}) \leq f(\hat{x}),
$$

where the last inequality follows from $\bar{y} \geq 0$ and $g_{j}(\hat{x}) \leq 0(j=1, \ldots, m)$. This completes the proof.

Theorem 3.10 (Strong duality for the Wolfe dual) Assume that $\mathcal{C}=\mathbb{R}^{n}$ and the functions $f, g_{1}, \cdots, g_{m}$ are continuously differentiable and convex. Also assume that (CO) satisfies the Slater regularity condition. Let $\bar{x}$ be a feasible solution of (CO). Then $\bar{x}$ is an optimal solution of (CO) if and only if there exists $a \bar{y} \geq 0$ such that $(\bar{x}, \bar{y})$ is an optimal solution of (WD).

Proof: Follows directly from Corollary 2.33.

Warning! Remember, we are only allowed to form the Wolfe dual of a nonlinear optimization problem if it is a convex optimization problem. We may replace the infimum in the definition of $\psi(y)$ by the condition that the $x$-gradient is zero only if all the functions $f$ and $g_{j}, \forall j$ are convex and if we know that the infimum is attained. Else, the condition

$$
\nabla f(x)+\sum_{j=1}^{m} y_{j} \nabla g_{j}(x)=0
$$

allows solutions which are possibly maxima, saddle points or inflection points, or it may not have any solution. In such cases no duality relation holds in general. For nonconvex problems one has to work with the Lagrange dual.

Example 3.11 [Wolfe dual] Let us consider the convex optimization problem

$$
\text { (CO) } \begin{aligned}
\min x_{1}+e^{x_{2}} & \\
\text { s.t. } 3 x_{1}-2 e^{x_{2}} & \geq 10 \\
x_{2} & \geq 0 \\
x & \in \mathbb{R}^{2} .
\end{aligned}
$$

Then the optimal value is 5 with $x=(4,0)$. Note that the Slater condition holds for this example.
Wolfe dual The Wolfe dual of (CO) is given by

$$
\begin{array}{cc}
\sup & x_{1}+e^{x_{2}}+y_{1}\left(10-3 x_{1}+2 e^{x 2}\right)-y_{2} x_{2}  \tag{WD}\\
\text { s.t. } & 1-3 y_{1}=0 \\
& e^{x_{2}}+2 e^{x_{2}} y_{1}-y_{2}=0 \\
x \in \mathbb{R}^{2}, y \geq 0,
\end{array}
$$

which is a non-convex problem. The first constraint gives $y_{1}=\frac{1}{3}$, and thus the second constraint becomes

$$
\frac{5}{3} e^{x_{2}}-y_{2}=0
$$

Now we can eliminate $y_{1}$ and $y_{2}$ from the object function. We get the function

$$
f\left(x_{2}\right)=\frac{5}{3} e^{x_{2}}-\frac{5}{3} x_{2} e^{x_{2}}+\frac{10}{3}
$$

This function has a maximum when

$$
f^{\prime}\left(x_{2}\right)=-\frac{5}{3} x_{2} e^{x_{2}}=0
$$

which is only true when $x_{2}=0$ and $f(0)=5$. Hence the optimal value of (WD) is 5 and then $(x, y)=\left(4,0, \frac{1}{3}, \frac{5}{3}\right)$.
Lagrange dual We can double check this answer by using the Lagrange dual. Let

$$
\begin{aligned}
\psi(y) & =\inf _{x \in \mathbb{R}^{2}}\left\{x_{1}+e^{x_{2}}+y_{1}\left(10-3 x_{1}+2 e^{x 2}\right)-y_{2} x_{2}\right\} \\
& =\inf _{x_{1} \in \mathbb{R}}\left\{x_{1}-3 y_{1} x_{1}\right\}+\inf _{x_{2} \in \mathbb{R}}\left\{\left(1+2 y_{1}\right) e^{x_{2}}-y_{2} x_{2}\right\}+10 y_{1}
\end{aligned}
$$

We have

$$
\inf _{x_{1} \in \mathbb{R}}\left\{x_{1}-3 y_{1} x_{1}\right\}= \begin{cases}0 & \text { for } y_{1}=\frac{1}{3} \\ -\infty & \text { otherwise }\end{cases}
$$

Now, for fixed $y_{1}, y_{2}$, with $y_{2}>0$ let

$$
g\left(x_{2}\right)=\left(1+2 y_{1}\right) e^{x_{2}}-y_{2} x_{2}
$$

Then $g$ has a minimum when

$$
g^{\prime}\left(x_{2}\right)=\left(1+2 y_{1}\right) e^{x_{2}}-y_{2}=0
$$

i.e., when $x_{2}=\log \left(\frac{y_{2}}{1+2 y_{1}}\right)$. Further, $g\left(\log \left(\frac{y_{2}}{1+2 y_{1}}\right)\right)=y_{2}-y_{2} \log \left(\frac{y_{2}}{1+2 y_{1}}\right)$. Hence, we have

$$
\inf _{x_{2} \in \mathbb{R}}\left\{\left(1+2 y_{1}\right) e^{x_{2}}-y_{2} x_{2}\right\}= \begin{cases}y_{2}-y_{2} \log \left(\frac{y_{2}}{1+2 y_{1}}\right) & \text { for } y_{2}>0 \\ 0 & \text { for } y_{2}=0\end{cases}
$$

Thus the Lagrange dual becomes

$$
\begin{aligned}
(\mathrm{LD}) \quad \sup \psi(y) & =10 y_{1}+y_{2}-y_{2} \log \left(\frac{y_{2}}{1+2 y_{1}}\right) \\
\text { s.t. } \quad y_{1} & =\frac{1}{3} \\
y_{2} & \geq 0
\end{aligned}
$$

Now we have

$$
\frac{d}{d y_{2}} \psi\left(\frac{1}{3}, y_{2}\right)=\log \left(\frac{3 y_{2}}{5}\right)=0
$$

when $y_{2}=\frac{5}{3}$, and $\psi\left(\frac{1}{3}, \frac{5}{3}\right)=5$.

Exercise 3.2 Prove that - under the assumptions of Theorem 3.10 - the Lagrange and Wolfe duals of the optimization problem (CO) are equivalent.

Exercise 3.3 We wish to design a rectangular prism (box) with length l, width $b$, and height $h$ such that the volume of the box is at least $V$ units, and the total surface area is minimal. This problem has the following (nonconvex) formulation:

$$
\begin{equation*}
\min _{l, b, h} 2(l b+b h+l h), \quad l b h \geq V, \quad l, b, h>0 \tag{3.2}
\end{equation*}
$$

i) Transform the problem (3.2) by introducing new variables to obtain:

$$
\begin{equation*}
\min _{x_{1}, x_{2}, x_{3}} 2\left(e^{x_{1}+x_{2}}+e^{x_{2}+x_{3}}+e^{x_{1}+x_{3}}\right), \quad x_{1}+x_{2}+x_{3} \geq \ln (V), \quad x_{1}, x_{2}, x_{3} \in \mathbb{R} \tag{3.3}
\end{equation*}
$$

ii) Show that the transformed problem is convex and satisfies Slater's regularity condition.
iii) Show that the Lagrange dual of problem (3.3) is:

$$
\begin{equation*}
\max _{\lambda \geq 0}\left(\frac{3}{2}+\ln (V)\right) \lambda-\frac{3}{2} \lambda \ln \left(\frac{\lambda}{4}\right) . \tag{3.4}
\end{equation*}
$$

iv) Show that the Wolfe dual of problem (3.3) is the same as the Lagrange dual.
v) Use the KKT conditions of problem (3.3) to show that the cube $\left(l=b=h=V^{1 / 3}\right)$ is the optimal solution of problem (3.2).
vi) Use the dual problem (3.4) to derive the same result as in part v).

### 3.3 Examples for dual problems

In this section we derive the Lagrange and/or the Wolfe dual of some specific convex optimization problems.

## Linear optimization

Let $A: m \times n$ be a matrix, $b \in \mathbb{R}^{m}$ and $c, x \in \mathbb{R}^{n}$. The primal Linear Optimization (LO) problem is given as

$$
(\mathrm{LO}) \quad \min \left\{c^{T} x \mid A x=b, \quad x \geq 0\right\}
$$

Here we can say that $\mathcal{C}=\mathbb{R}^{n}$. Obviously all the constraints are continuously differentiable. The inequality constraints can be given as $g_{j}(x)=\left(a^{j}\right)^{T} x-b_{j}$ if $j=1, \cdots, m$ and $g_{j}(x)=\left(-a^{j-m}\right)^{T} x+b_{j-m}$ if $j=m+1, \cdots, 2 m$ and finally $g_{j}(x)=-x_{j-2 m}$ if $j=2 m+1, \cdots, 2 m+n$. Here $a^{j}$ denotes the $j$ th row of matrix $A$. Denoting the Lagrange multipliers by $y^{-}, y^{+}$and $s$ respectively the Wolfe dual (WD) of (LO) has the following form:

$$
\begin{gathered}
\max c^{T} x+\left(y^{-}\right)^{T}(A x-b)+\left(y^{+}\right)^{T}(-A x+b)+s^{T}(-x) \\
c+A^{T} y^{-}-A^{T} y^{+}-s=0 \\
y^{-} \geq 0, \quad y^{+} \geq 0, \quad s \geq 0
\end{gathered}
$$

As we substitute $c=-A^{T} y^{-}+A^{T} y^{+}+s$ in the objective and introduce the notation $y=y^{+}-y^{-}$the standard dual linear optimization problem follows.

$$
\begin{gathered}
\max \\
b^{T} y \\
A^{T} y+s=c, \\
s \geq 0
\end{gathered}
$$

## Quadratic optimization

The quadratic optimization problem is considered in the symmetric form. Let $A: m \times n$ be a matrix, $Q: n \times n$ be a positive semi-definite symmetric matrix, $b \in \mathbb{R}^{m}$ and $c, x \in \mathbb{R}^{n}$. The primal Quadratic Optimization (QO) problem is given as

$$
\text { (QO) } \quad \min \left\{\left.c^{T} x+\frac{1}{2} x^{T} Q x \right\rvert\, A x \geq b, \quad x \geq 0\right\}
$$

Here we can say that $\mathcal{C}=\mathbb{R}^{n}$. Obviously all the constraints are continuously differentiable. The inequality constraints can be given as $g_{j}(x)=\left(-a^{j}\right)^{T} x+b_{j}$ if $j=1, \cdots, m$ and $g_{j}(x)=-x_{j-m}$ if $j=m+1, \cdots, m+n$. Denoting the Lagrange multipliers by $y$ and $s$ respectively the Wolfe dual (WD) of (QO) has the following form:

$$
\begin{gathered}
\max c^{T} x+\frac{1}{2} x^{T} Q x+y^{T}(-A x+b)+s^{T}(-x) \\
c+Q x-A^{T} y-s=0 \\
y \geq 0, \quad s \geq 0
\end{gathered}
$$

As we substitute $c=-Q x+A^{T} y+s$ in the objective the dual quadratic optimization problem follows.

$$
\begin{gathered}
\max \quad b^{T} y-\frac{1}{2} x^{T} Q x \\
-Q x+A^{T} y+s=c, \\
y \geq 0, \quad s \geq 0
\end{gathered}
$$

Observe that the vector $x$ occurring in this dual is not necessarily feasible for (QO)! To eliminate the $x$ variables another form of the dual can be presented.

Since $Q$ is a positive semidefinite symmetric matrix, it can be represented as the product of two matrices $Q=D^{T} D$ (use e.g. Cholesky factorization), one can introduce the vector $z=D x$. Hence the following (QD) dual problem is obtained:

$$
\begin{gathered}
\max \quad b^{T} y-\frac{1}{2} z^{T} z \\
-D^{T} z+A^{T} y+s=c \\
y \geq 0, \quad s \geq 0
\end{gathered}
$$

Note that the optimality conditions are $x^{T} s=0, y^{T}(A x-b)=0$ and $z=D x$.

## Constrained maximum likelihood estimation

Maximum Likelihood Estimation frequently occurs in statistics. This problem can also be used to illustrate duality in convex optimization. In this problem we are given a finite set of sample points $x_{i}, \quad(1 \leq i \leq n)$. The most probable density values at the sample points are to be determined that satisfy some linear (e.g. convexity) constraints. Formally, the problem is defined as one has to determine the maximum of the Likelihood function $\prod_{i=1}^{n} x_{i}$ under the conditions

$$
A x \geq 0, \quad d^{T} x=1, \quad x \geq 0
$$

Here $A x \geq 0$ represents the linear constraints, the density values $x_{i}$ are nonnegative and the condition $d^{T} x=1$ ensures that the (approximate) integral of the density function is one. Since the logarithm function is monotone the objective can equivalently replaced by

$$
\min \quad-\sum_{i=1}^{n} \ln x_{i} .
$$

It is easy to check that the so defined problem is a convex optimization problem. Again we can take $\mathcal{C}=\mathbb{R}^{n}$ and all the constraints are linear, hence continuously differentiable. Denoting the Lagrange multipliers by $y \in \mathbb{R}^{m}, t \in \mathbb{R}$ and $s \in \mathbb{R}^{n}$ respectively the Wolfe dual (WD) of this problem has the following form:

$$
\begin{gathered}
\max -\sum_{i=1}^{n} \ln x_{i}+y^{T}(-A x)+t\left(d^{T} x-1\right)+s^{T}(-x) \\
-X^{-1} e-A^{T} y+t d-s=0 \\
y \geq 0, \quad s \geq 0
\end{gathered}
$$

Here the notation $e=(1, \cdots, 1) \in \mathbb{R}^{n}$ and $X=\operatorname{diag}(x)$ is used. Also note that for simplicity we did not split the equality constraint into two inequalities but we used immediately that its multiplier is a free variable. Multiplying the first constraint by $x^{T}$ one has

$$
-x^{T} X^{-1} e-x^{T} A^{T} y+t x^{T} d-x^{T} s=0
$$

Using $d^{T} x=1, x^{T} X^{-1} e=n$ and the optimality conditions $y^{T} A x=0, x^{T} s=0$ we have

$$
t=n .
$$

Observe further that due to the logarithm in the primal objective, the primal optimal solution is necessarily strictly positive, hence the dual variable $s$ must be zero at the optimum. Combining these results the dual problem is

$$
\begin{gathered}
\text { max }-\sum_{i=1}^{n} \ln x_{i} \\
X^{-1} e+A^{T} y=n d, \\
y \geq 0 .
\end{gathered}
$$

Eliminating the variables $x_{i}>0$ from the constraints one has $x_{i}=\frac{1}{n d_{i}-a_{i}^{T} y}$ and $-\ln x_{i}=$ $\ln \left(n d_{i}-a_{i}^{T} y\right)$ for all $i=1, \cdots, n$. Now we have the final form of our dual problem:

$$
\begin{gathered}
\max \sum_{i=1}^{n} \ln \left(n d_{i}-a_{i}^{T} y\right) \\
A^{T} y \leq n d, \\
y \geq 0 .
\end{gathered}
$$

### 3.4 Some examples with positive duality gap

Example 3.12 This example is due to Duffin. It shows that positive duality gap might occur for convex problems when the problem does not satisfy the Slater regularity condition. Moreover, it makes clear that the Wolfe dual might be significantly weaker than the Lagrange dual.
Let us consider the convex optimization problem

$$
\begin{array}{rll}
(\mathrm{CO}) & \min & e^{-x_{2}} \\
& \text { s.t. } & \sqrt{x_{1}^{2}+x_{2}^{2}}-x_{1} \leq 0 \\
& x \in \mathbb{R}^{2} .
\end{array}
$$

The feasible region is $\mathcal{F}=\left\{x \in \mathbb{R}^{2} \mid x_{1} \geq 0, x_{2}=0\right\}$. The only constraint is non-linear and singular, thus (CO) is not Slater regular. The optimal value of the object function is 1 .

The Lagrange function is given by

$$
L(x, y)=e^{-x_{2}}+y\left(\sqrt{x_{1}^{2}+x_{2}^{2}}-x_{1}\right)
$$

Let us first consider the Wolfe dual (WD):

$$
\begin{aligned}
\sup e^{-x_{2}}+y\left(\sqrt{x_{1}^{2}+x_{2}^{2}}-x_{1}\right) & \\
-y+y \frac{x_{1}}{\sqrt{x_{1}^{2}+x_{2}^{2}}} & =0 \\
-e^{-x_{2}}+y \frac{x_{2}}{\sqrt{x_{1}^{2}+x_{2}^{2}}} & =0 \\
y & \geq 0
\end{aligned}
$$

The first constraint imply that $x_{2}=0$ and $x_{1} \geq 0$, but these values do not satisfy the second constraint. Thus the Wolfe dual is infeasible, yielding an infinitely large duality gap.

Let us see if we can do better by using the Lagrange dual. Now, let $\epsilon=\sqrt{x_{1}^{2}+x_{2}^{2}}-x_{1}$, then

$$
x_{2}^{2}-2 \epsilon x_{1}-\epsilon^{2}=0
$$

Hence, for any $\epsilon>0$ we can find $x_{1}>0$ such that $\epsilon=\sqrt{x_{1}^{2}+x_{2}^{2}}-x_{1}$ even if $x_{2}$ goes to infinity. However, when $x_{2}$ goes to infinity $e^{-x_{2}}$ goes to 0 . So,

$$
\psi(y)=\inf _{x \in \mathbb{R}^{2}} e^{-x_{2}}+y\left(\sqrt{x_{1}^{2}+x_{2}^{2}}-x_{1}\right)=0
$$

thus the optimal value of the Lagrange dual

$$
\begin{array}{lcl}
(\mathrm{LD}) & \max & \psi(y) \\
& \text { s.t. } & y \geq 0
\end{array}
$$

is 0 . This gives a nonzero duality gap that equals to 1 .
Observe that the Wolfe dual becomes infeasible because the infimum in the definition of $\psi(y)$ exists, but it is not attained.

Example 3.13 [Basic model with zero duality gap] Let us first consider the following simple convex optimization problem.

$$
\begin{align*}
\min & x_{1} \\
\text { s.t. } & x_{1}^{2}  \tag{3.5}\\
-x_{2} & \leq 0 \\
-1-x_{1} & \leq 0
\end{align*}
$$

Here the convex set $\mathcal{C}$ where the above functions are defined is $\mathbb{R}^{2}$. It is clear that the set of feasible solutions is given by

$$
\mathcal{F}=\left\{\left(x_{1}, x_{2}\right) \mid x_{1}=0, \quad x_{2} \geq 0\right\}
$$

thus any feasible vector $\left(x_{1}, x_{2}\right) \in \mathcal{F}$ is optimal and the optimal value of this problem is 0 . Because $x_{1}=0$ for all feasible solutions the Slater regularity condition does not hold for (3.5).

Let us make the Lagrange dual of (3.5). The Lagrange multipliers $\left(y_{1}, y_{2}, y_{3}\right)$ are nonnegative and the Lagrange function

$$
L(x, y)=x_{1}+y_{1} x_{1}^{2}-y_{2} x_{2}-y_{3}\left(1+x_{1}\right)
$$

is defined on $x \in \mathbb{R}^{2}$ and $y \in \mathbb{R}^{3}, y \geq 0$.
The Lagrange dual is defined as

$$
\begin{array}{cl}
\max & \psi(y)  \tag{3.6}\\
\text { s.t. } \quad y \geq 0
\end{array}
$$

where

$$
\begin{aligned}
\psi(y) & =\inf _{x \in \mathbb{R}^{2}}\left\{x_{1}+y_{1} x_{1}^{2}-y_{2} x_{2}-y_{3}\left(1+x_{1}\right)\right\} \\
& =\inf _{x \in \mathbb{R}^{2}}\left\{x_{1}\left(1-y_{3}\right)+y_{1} x_{1}^{2}-y_{2} x_{2}-y_{3}\right\} \\
& = \begin{cases}-\infty & \text { if } \quad y_{2} \neq 0 \text { or } y_{1}=0 \text { but } y_{3} \neq 1 \\
0 & \text { if } \quad y_{2}=0, \quad y_{1}=0 \text { and } y_{3}=1 \\
-y_{3}-\frac{\left(1-y_{3}\right)^{2}}{4 y_{1}} & y_{2}=0 \text { and } y_{1} \neq 0\end{cases}
\end{aligned}
$$

The last expression in the formula above is obtained by minimizing the convex quadratic function $x_{1}\left(1-y_{3}\right)+y_{1} x_{1}^{2}-y_{3}$ where $y_{1}$ and $y_{3}$ are fixed. Because this last expression is nonpositive, the maximum of $\psi(y)$ is zero. Thus for this problem both the primal and the dual problems have optimal solutions with equal (zero) optimal objective values.

Example 3.14 [A variant with positive duality gap] Let us consider the same problem as in the previous example (see problem (3.5)) with a different representation of the feasible set. As we will see the new formulation results in a quite different dual. The new dual has also an optimal solution but now the duality gap is positive.

$$
\begin{align*}
\min x_{1} & \\
\text { s.t. } x_{0}-s_{0} & =0 \\
x_{1}-s_{1} & =0 \\
x_{2}-s_{2} & =0  \tag{3.7}\\
1+x_{1}-s_{3} & =0 \\
x_{0} & =0 \\
x \in \mathbb{R}^{3}, s \in \mathcal{C} &
\end{align*}
$$

Note that (3.7) has the correct form: the constraints are linear, hence convex, and the vector $(x, s)$ of the variables belong to the convex set $\mathbb{R}^{3} \times \mathcal{C}$. Here the convex set $\mathcal{C}$ is defined as follows:

$$
\mathcal{C}=\left\{s=\left(s_{0}, s_{1}, s_{2}, s_{3}\right) \mid s_{0} \geq 0, s_{2} \geq 0, s_{3} \geq 0, s_{0} s_{2} \geq s_{1}^{2}\right\}
$$

It is clear that the set of feasible solutions is

$$
\mathcal{F}=\left\{(x, s) \mid x_{0}=0, x_{1}=0, x_{2} \geq 0, s_{0}=0, s_{1}=0, s_{2} \geq 0, s_{3}=1\right\}
$$

thus any feasible vector $(x, s) \in \mathcal{F}$ is optimal and the optimal value of this problem is 0 .

## Exercise 3.4 1. Prove that the function $s_{1}^{2}-s_{0} s_{2}$ is not convex.

2. Prove that the set $\mathcal{C}$ is convex.
3. Prove that problem (3.7) does not satisfy the Slater regularity condition.

Due to the equality constraints the Lagrange multipliers $\left(y_{0}, y_{1}, y_{2}, y_{3}, y_{4}\right)$ are free and the Lagrange function

$$
\left.L(x, s, y)=x_{1}+y_{0}\left(x_{0}-s_{0}\right)+y_{1}\left(x_{1}-s_{1}\right)+y_{2}\left(x_{2}-s_{2}\right)+y_{3}\left(1+x_{1}-s_{3}\right)++y_{4} x_{0}\right)
$$

is defined for $x \in \mathbb{R}^{3}, s \in \mathcal{C}$ and $y \in \mathbb{R}^{5}$.
The Lagrange dual is defined as

$$
\begin{array}{r}
\max \quad \psi(y)  \tag{3.8}\\
\text { s.t. } \quad y \in \mathbb{R}^{5}
\end{array}
$$

where

$$
\begin{aligned}
\psi(y) & =\inf _{x \in \mathbb{R}^{3}, s \in \mathcal{C}} L(x, s, y) \\
& =\inf _{x \in \mathbb{R}^{3}, s \in \mathcal{C}}\left\{x_{1}\left(1+y_{1}+y_{3}\right)+x_{0}\left(y_{4}+y_{0}\right)+x_{2} y_{2}-s_{0} y_{0}-s_{1} y_{1}-s_{2} y_{2}-s_{3} y_{3}+y_{3}\right\} \\
& =\left\{\begin{array}{lr}
y_{3} & \text { if } 1+y_{1}+y_{3}=0, y_{4}+y_{0}=0, y_{2}=0, y_{3} \leq 0, y_{0} \leq 0, y_{1}=0 ; \\
-\infty & \text { otherwise } .
\end{array}\right.
\end{aligned}
$$

The last equality requires some explanation.

- If $1+y_{1}+y_{3} \neq 0$ then $L(x, s, y)=x_{1}\left(1+y_{1}+y_{3}\right)+y_{3}$ for $x_{0}=x_{2}=0, s=0$. So $\inf L(x, s, y)=-\infty$ in this case.
- If $y_{4}+y_{0} \neq 0$ then $L(x, s, y)=x_{0}\left(y_{4}+y_{0}\right)+y_{3}$ for $x_{1}=x_{2}=0, s=0$. So inf $L(x, s, y)=-\infty$ in this case.
- If $y_{2} \neq 0$ then $L(x, s, y)=x_{2} y_{2}+y_{3}$ for $x_{0}=x_{1}=0, s=0$. So inf $L(x, s, y)=-\infty$ in this case.
- If $y_{0}>0$ then $L(x, s, y)=-s_{0} y_{0}+y_{3}$ for $x=0, s_{0} \geq 0, s_{1}=0, s_{2}=0, s_{3}=0$. So $\inf L(x, s, y)=-\infty$ in this case.
- If $y_{3}>0$ then $L(x, s, y)=-s_{3} y_{3}+y_{3}$ for $x=0, s_{0}=0, s_{1}=0, s_{2}=0, s_{3} \geq 0$. So $\inf L(x, s, y)=-\infty$ in this case.
- If $y_{2}=0$ but $y_{1} \neq 0$ then $L(x, s, y)=-\frac{1}{\tau} y_{0}-\frac{y_{1}}{\left|y_{1}\right|} \tau y_{1}+y_{3}$ for $x=0, s_{3}=0$ and $\left(s_{0}, s_{1}, s_{2}, s_{3}\right)=$ $\left(\frac{1}{\tau}, \frac{y_{1}}{\left|y_{1}\right|} \tau, \tau^{2}, 0\right) \in \mathcal{C}$. So $\inf L(x, s, y)=-\infty$ (one obtains this a let $\left.\tau \rightarrow \infty\right)$ in this case.

Summarizing the above results we conclude that the Lagrange dual reduces to

$$
\begin{gathered}
\max y_{3} \\
y_{0} \leq 0, y_{1}=0, y_{2}=0, y_{3}=-1, y_{4}=-y_{0}
\end{gathered}
$$

Here for any feasible solution $y_{3}=-1$, thus the optimal value of the Lagrange dual is -1 , i.e. both the primal problem (3.7) and its dual (3.8) have optimal solutions, but their optimal values are not equal.

Exercise 3.5 Modify the above problem so that for a given $\gamma>0$ the nonzero duality gap at optimum will be equal to $\gamma$.

Example 3.15 [Duality for non convex problems 1] Let us consider the non-convex optimization problem

$$
\begin{array}{lll}
(\mathrm{NCO}) & \min & x_{1}^{2}-2 x_{2} \\
& \text { s.t. } & x_{1}^{2}+x_{2}^{2}=4 \\
& x \in \mathbb{R}^{2}
\end{array}
$$

Then the optimal value is -4 , with $x=(0,2)$.

Lagrange dual The Lagrange function of (NCO) is given by

$$
L(x, y)=x_{1}^{2}-2 x_{2}+y\left(x_{1}^{2}+x_{2}^{2}-4\right), \quad \text { where } y \in \mathbb{R}
$$

and then

$$
\begin{aligned}
\psi(y) & =\inf _{x \in \mathbb{R}^{2}} L(x, y) \\
& =\inf _{x_{1}}\left\{(1+y) x_{1}^{2}\right\}+\inf _{x_{2}}\left\{y x_{2}^{2}-2 x_{2}\right\}-4 y
\end{aligned}
$$

We have

$$
\begin{aligned}
& \inf _{x_{1}}\left\{(1+y) x_{1}^{2}\right\}= \begin{cases}0 & \text { for } y \geq-1 \\
-\infty & \text { for } y<-1\end{cases} \\
& \inf _{x_{2}}\left\{y x_{2}^{2}-2 x_{2}\right\}= \begin{cases}-\frac{1}{y} & \text { for } y>0 \\
-\infty & \text { for } y \leq 0\end{cases}
\end{aligned}
$$

Hence, the Lagrange dual is

$$
\begin{gathered}
(\mathrm{LD}) \sup \quad-\frac{1}{y}-4 y \\
y>0
\end{gathered}
$$

which is a convex problem, and the optimal value is -4 , with $y=\frac{1}{2}$. Note that although the problem is not convex, and does not satisfy the Slater regularity condition, the duality gap is zero.

Example 3.16 [Duality for non convex problems 2] Let us consider the non-convex optimization problem

$$
\begin{array}{lll}
(\mathrm{CLO}) & \min & x_{1}^{2}-x_{2}^{2} \\
& \text { s.t. } & x_{1}+x_{2} \leq 2 \\
& x \in \mathcal{C}=\left\{x \in \mathbb{R}^{2} \mid-2 \leq x_{1}, x_{2} \leq 4\right\}
\end{array}
$$

Then we have the optimal value -12 with $x=(-2,4)$. The Lagrange function of (CLO) is given by

$$
L(x, y)=x_{1}^{2}-x_{2}^{2}+y\left(x_{1}+x_{2}-2\right), \quad \text { where } \quad y \geq 0
$$

Thus for $y \geq 0$ we have

$$
\begin{aligned}
\psi(y) & =\inf _{x \in \mathcal{C}} L(x, y) \\
& =\inf _{-2 \leq x_{1} \leq 4}\left\{x_{1}^{2}+y x_{1}\right\}+\inf _{-2 \leq x_{2} \leq 4}\left\{-x_{2}^{2}+y x_{2}\right\}-2 y
\end{aligned}
$$

Now, $x_{1}^{2}+y x_{1}$ is a parabola which has its minimum at $x_{1}=-\frac{y}{2}$. So, this minimum lies within $\mathcal{C}$ when $y \leq 4$. When $y \geq 4$ the minimum is reached at the boundary of $\mathcal{C}$. The minimum of the parabola $-x_{2}^{2}+y x_{2}$ is always reached at the boundaries of $\mathcal{C}$, at $x_{2}=-2$ when $y \geq 2$, and at $x_{2}=4$ when $y \leq 2$. Hence, we have

$$
\psi(y)= \begin{cases}-\frac{y^{2}}{4}+2 y-16 & \text { for } y \leq 2 \\ -\frac{y^{2}}{4}-4 y-4 & \text { for } 2 \leq y \leq 4 \\ -6 y & \text { for } y \geq 4\end{cases}
$$

Maximizing $\psi(y)$ for $y \geq 0$ gives

$$
\begin{aligned}
\sup _{0 \leq y \leq 2} \psi(y) & =-13 \\
\sup _{2 \leq y \leq 4} \psi(y) & =-13 \\
\sup _{y \geq 4} \psi(y) & =-24
\end{aligned}
$$

Hence, the optimal value of the Lagrange dual is -13 , and we have a nonzero duality gap that equals to 1 .

### 3.5 Semidefinite optimization

## The Primal and the Dual Problem

Let $A_{0}, A_{1}, \cdots, A_{n} \in \mathbb{R}^{m \times m}$ be symmetric matrices. Further let $c \in \mathbb{R}^{n}$ be a given vector and $x \in \mathbb{R}^{n}$ be the vector of unknowns in which the optimization is done. The primal semidefinite optimization problem is defined as

$$
\begin{array}{lll}
(P S O) & \min & c^{T} x  \tag{3.9}\\
& \text { s.t. } & -A_{0}+\sum_{k=1}^{n} A_{k} x_{k} \succeq 0,
\end{array}
$$

where $\succeq 0$ indicates that the left hand side matrix has to be positive semidefinite. It is clear that the primal problem ( $P S O$ ) is a convex optimization problem since the
convex combination of positive semidefinite matrices is also positive semidefinite. For convenience the notation

$$
F(x)=-A_{0}+\sum_{k=1}^{n} A_{k} x_{k}
$$

will be used.
The dual problem of the semidefinite optimization problem, as given e.g. in [44], is as follows:

$$
\begin{align*}
(D S P) \quad \max & \operatorname{Tr}\left(A_{0} Z\right)  \tag{3.10}\\
\text { s.t. } & \operatorname{Tr}\left(A_{k} Z\right)=c_{k}, \quad \text { for all } \quad k=1, \cdots, n, \\
& Z \succeq 0,
\end{align*}
$$

where $Z \in \mathbb{R}^{m \times m}$ is the matrix of variables. Again, the dual of the semidefinite optimization problem is convex. The trace of a matrix is a linear function of the matrix and the convex combination of positive semidefinite matrices is also positive semidefinite.

Theorem 3.17 (Weak duality) If $x \in \mathbb{R}^{n}$ is primal feasible and $Z \in \mathbb{R}^{m \times m}$ is dual feasible, then

$$
c^{T} x \geq \operatorname{Tr}\left(A_{0} Z\right)
$$

with equality if and only if

$$
F(x) Z=0
$$

Proof: Using the dual constraints and some elementary properties of the trace of matrices one may write

$$
c^{T} x-\operatorname{Tr}\left(A_{0} Z\right)=\sum_{k=1}^{n} \operatorname{Tr}\left(A_{k} Z\right) x_{k}-\operatorname{Tr}\left(A_{0} Z\right)=\operatorname{Tr}\left(\left(\sum_{k=1}^{n} A_{k} x_{k}-A_{0}\right) Z\right)=\operatorname{Tr}(F(x) Z) \geq 0 .
$$

Here the last inequality holds because both matrices $F(x)$ and $Z$ are positive semidefinite. Equality holds if and only if $F(x) Z=0$, which completes the proof.

## The Dual as Lagrange-Wolfe Dual

First we give another equivalent form of the ( $P S O$ ) problem in order to be able to derive the dual problem more easily. Clearly problem (PSO) can equivalently be given as

$$
\begin{array}{ll}
\left(P S O^{\prime}\right) \quad & \min ^{T} x  \tag{3.11}\\
\text { s.t. } & -F(x)+S=0 \\
& S \succeq 0,
\end{array}
$$

where $S \in \mathbb{R}^{m \times m}$ is a symmetric matrix. It plays the role of the usual "slack variables". The Lagrange function $L(x, S, Z)$ of problem $\left(P S O^{\prime}\right)$ is defined on the set $\left\{(x, S, Z) \mid x \in \mathbb{R}^{n}, S \in \mathbb{R}^{m \times m}, S \succeq 0, Z \in \mathbb{R}^{m \times m},\right\}$ and is given by

$$
L(x, S, Z)=c^{T} x-e^{T}(F(x) \circ Z) e+e^{T}(S \circ Z) e
$$

where $e^{T}=(1, \cdots, 1) \in \mathbb{R}^{n}$ and $X \circ Z$ denotes the Minkowski (coordinatewise) product of matrices. Before going on we observe that $e^{T}(S \circ Z) e=\operatorname{Tr}(S Z)$, hence the Lagrange function can be reformulated as

$$
\begin{equation*}
L(x, S, Z)=c^{T} x-\sum_{k=1}^{n} x_{k} \operatorname{Tr}\left(A_{k} Z\right)+\operatorname{Tr}\left(A_{0} Z\right)+\operatorname{Tr}(S Z) \tag{3.12}
\end{equation*}
$$

Before formulating the Lagrange dual of $\left(P S O^{\prime}\right)$ note that we can assume that the matrix $Z$ is symmetric, since $F(x)$ is symmetric. The Lagrange dual of problem $\left(P S O^{\prime}\right)$ is

$$
\begin{equation*}
(D S D L) \quad \max \left\{\psi(Z) \mid Z \in \mathbb{R}^{m \times m}\right\} \tag{3.13}
\end{equation*}
$$

where

$$
\begin{equation*}
\psi(Z)=\min \left\{L(x, S, Z) \mid x \in \mathbb{R}^{n}, S \in \mathbb{R}^{m \times m}, S \succeq 0\right\} \tag{3.14}
\end{equation*}
$$

As we did in deriving the Wolfe dual, one easily derives optimality conditions to calculate $\psi(Z)$. Since the minimization in (3.14) is done in the free variable $x$, the positive semidefinite matrix of variables $S$ and, further the function $L(x, S, Z)$ is separable w.r.t. $x$ and $S$ we can take these minimums separately.

If we minimize in $S$ all the terms in (3.12) but $\operatorname{Tr}(S Z)$ are constant. The matrix $S$ is positive semidefinite, hence

$$
\min _{S} \operatorname{Tr}(S Z)= \begin{cases}0 & \text { if } Z \succeq 0  \tag{3.15}\\ -\infty & \text { otherwise }\end{cases}
$$

If we minimize (3.14) in $x$, we need to equate the $x$-gradient of $L(x, S, Z)$ to zero (remember to the Wolfe dual). This requirement leads to

$$
\begin{equation*}
c_{k}-\operatorname{Tr}\left(A_{k} Z\right)=0 \quad \text { for all } \quad k=1, \cdots, n \tag{3.16}
\end{equation*}
$$

Multiplying the equations of (3.16) by $x_{k}$ and summing up one obtains

$$
c^{T} x-\sum_{k=1}^{n} x_{k} \operatorname{Tr}\left(A_{k} Z\right)=0
$$

By combining the last formula and the results presented in (3.15) and in (3.16) the simplified form of the Lagrange dual (3.13), the Lagrange-Wolfe dual

$$
\begin{aligned}
(D S O) \quad \max & \operatorname{Tr}\left(A_{0} Z\right) \\
\text { s.t. } & \operatorname{Tr}\left(A_{k} Z\right)=c_{k}, \quad \text { for all } \quad k=1, \cdots, n, \\
& Z \succeq 0,
\end{aligned}
$$

follows. The reader readily verifies that this is identical to (3.10).

Exercise 3.6 Consider the problem given in Example 3.13.

1. Prove that problem (3.7) is a semidefinite optimization problem.
2. Give the dual semidefinite problem.
3. Prove that there is a positive duality gap for this primal-dual semidefinite optimization pair.

### 3.6 Duality in cone-linear optimization

In this section we deal with cone-linear optimization problems. A cone-linear optimization problem is a natural generalization of the well known standard linear optimization problem

$$
\min \left\{c^{T} x \mid A x \geq b, x \geq 0\right\}
$$

where $A \in \mathbb{R}^{m \times n}, b \in \mathbb{R}^{m}$ and $c \in \mathbb{R}^{n}$. The inequality conditions can be reformulated by observing that the conditions $A x \geq b$ and $x \geq 0$ mean that the vector $A x-b$ and $x$ should be in the positive orthant

$$
\mathbb{R}_{+}^{m}:=\left\{x \in \mathbb{R}^{m} \mid x \geq 0\right\}
$$

and $\mathbb{R}_{+}^{n}$, respectively. One observes, that the positive orthants $\mathbb{R}_{+}^{m}$ and $\mathbb{R}_{+}^{n}$ are convex cones, i.e. the linear optimization problem can be restated as the following cone-linear optimization problem

$$
\begin{aligned}
\min c^{T} x & \\
A x-b & \in \mathbb{R}_{+}^{m} \\
x & \in \mathbb{R}_{+}^{n}
\end{aligned}
$$

The dual problem

$$
\max \left\{b^{T} y \mid A^{T} y \leq c, y \geq 0\right\}
$$

can similarly be reformulated in the conic form:

$$
\begin{aligned}
\max b^{T} y & \\
c-A^{T} y & \in \mathbb{R}_{+}^{n} \\
y & \in \mathbb{R}_{+}^{m} .
\end{aligned}
$$

The natural question arises: how one can derive dual problems for general cone-linear optimization problems where, in the above given formulation the simple polyhedral convex cones $\mathbb{R}_{+}^{m}$ and $\mathbb{R}_{+}^{n}$ are replaced by arbitrary convex cones $\mathcal{C}_{1} \subseteq \mathbb{R}^{m}$ and $\mathcal{C}_{2} \subseteq \mathbb{R}^{n}$. The cone-linear optimization problem is defined as follows:

$$
\begin{align*}
\min c^{T} x & \\
A x-b & \in \mathcal{C}_{1}  \tag{3.17}\\
x & \in \mathcal{C}_{2} .
\end{align*}
$$

## The Dual of a Cone-linear Problem

First, by introducing slack variables $s$, we give another equivalent form of the conelinear problem (3.17)

$$
\begin{aligned}
\min c^{T} x & \\
s-A x+b & =0 \\
s & \in \mathcal{C}_{1} \\
x & \in \mathcal{C}_{2} .
\end{aligned}
$$

In this optimization problem we have linear equality constraints $s-A x+b=0$ and the vector $(s, x)$ must be in the convex cone

$$
\mathcal{C}_{1} \times \mathcal{C}_{2}:=\left\{(s, x) \mid s \in \mathcal{C}_{1}, x \in \mathcal{C}_{2}\right\}
$$

The Lagrange function $L(s, x, y)$ of the above problem is defined on the set

$$
\left\{(s, x, y) \mid s \in \mathcal{C}_{1}, x \in \mathcal{C}_{2}, y \in \mathbb{R}^{m}\right\}
$$

and is given by

$$
\begin{equation*}
L(s, x, y)=c^{T} x+y^{T}(s-A x+b)=b^{T} y+s^{T} y+x^{T}\left(c-A^{T} y\right) \tag{3.18}
\end{equation*}
$$

Hence, the Lagrange dual of the cone-linear problem is given by

$$
\max _{y \in \mathbb{R}^{m}} \psi(y)
$$

where

$$
\begin{equation*}
\psi(y)=\min \left\{L(s, x, y) \mid s \in \mathcal{C}_{1}, x \in \mathcal{C}_{2}\right\} \tag{3.19}
\end{equation*}
$$

As we did in deriving the Wolfe dual, one easily derives optimality conditions to calculate $\psi(y)$. Since the minimization in (3.19) is done in the variables $s \in \mathcal{C}_{1}$ and $x \in \mathcal{C}_{2}$, and the function $L(s, x, y)$ is separable w.r.t. $x$ and $s$, we can take these minimums separately.

If we minimize in $s$ all the terms in (3.18) but $s^{T} y$ are constant. The vector $s$ is in the cone $\mathcal{C}_{1}$, hence

$$
\min _{s \in \mathcal{C}_{1}} s^{T} y= \begin{cases}0 & \text { if } y \in \mathcal{C}_{1}^{*}  \tag{3.20}\\ -\infty & \text { otherwise }\end{cases}
$$

If we minimize (3.19) in $x$ then all the terms in (3.18) but $x^{T}\left(c-A^{T} y\right)$ are constant. The vector $x$ is in the cone $\mathcal{C}_{2}$, hence

$$
\min _{x \in \mathcal{C}_{2}} x^{T}\left(c-A^{T} y\right)= \begin{cases}0 & \text { if } c-A^{T} y \in \mathcal{C}_{2}^{*}  \tag{3.21}\\ -\infty & \text { otherwise }\end{cases}
$$

By combining (3.20) and (3.21) we have

$$
\psi(y)= \begin{cases}b^{T} y & \text { if } y \in \mathcal{C}_{1}^{*} \text { and } c-A^{T} y \in \mathcal{C}_{2}^{*}  \tag{3.22}\\ -\infty & \text { otherwise }\end{cases}
$$

Thus the dual of the cone-linear optimization problem (3.17) is the following cone-linear problem:

$$
\begin{align*}
& \max b^{T} y \\
& c-A^{T} y \in \mathcal{C}_{2}^{*}  \tag{3.23}\\
& y \in \mathcal{C}_{1}^{*} .
\end{align*}
$$

Exercise 3.7 Derive the dual semidefinite optimization problem (DSO) by using the general cone-dual problem (3.23).

To illustrate the duality relation between (3.17) and (3.23) we prove the following weak duality theorem.

Theorem 3.18 (Weak duality) If $x \in \mathbb{R}^{n}$ is a feasible solution of the primal problem (3.17) and $y \in \mathbb{R}^{m}$ is a feasible solution of the dual problem (3.23) then

$$
c^{T} x \geq b^{T} y
$$

with equality if and only if

$$
x^{T}\left(c-A^{T} y\right)=0 \quad \text { and } \quad y^{T}(A x-b)=0 .
$$

Proof: Using the definition of the dual cone one may write

$$
c^{T} x-b^{T} y=x^{T}\left(c-A^{T} y\right)+y^{T}(A x-b) \geq 0 .
$$

Due to the nonnegativity of the vectors $x, c-A^{T} y, y$ and $A x-b$, equality holds if and only if $x^{T}\left(c-A^{T} y\right)=0$ and $y^{T}(A x-b)=0$, which completes the proof.

## Chapter 4

## Algorithms for unconstrained optimization

### 4.1 A generic algorithm

The problem considered in this section is

$$
\begin{array}{ll}
\min & f(x)  \tag{4.1}\\
\text { s.t. } & x \in \mathcal{C},
\end{array}
$$

where $\mathcal{C}$ is a relatively open convex set. For typical unconstrained optimization problems one has $\mathcal{C}=\mathbb{R}^{n}$, the trivial full dimensional open set, but for other applications (like in interior point methods) one frequently has lower dimensional relatively open convex sets. A generic algorithm for minimizing the function $f(x)$ can be presented as follows.

## Generic Algorithm

## Input:

$x^{0}$ is a given (relative interior) feasible point;
For $k=0,1, \ldots$ do
Step 1: Find a search direction $s^{k}$ with $\delta f\left(x^{k}, s^{k}\right)<0$;
(This should be a descending feasible direction in the constrained case.)
Step 1a: If no such direction exists STOP, optimum found.
Step 2: Line search : find $\lambda_{k}=\arg \min _{\lambda} f\left(x^{k}+\lambda s^{k}\right)$;
Step 3: $x^{k+1}=x^{k}+\lambda_{k} s^{k}, \quad k=k+1 ;$
Step 4: If stopping criteria are satisfied STOP.
The crucial elements of all algorithms, besides the selection of a starting point are printed boldface in the scheme, given above.

To generate a search direction is the crucial element of all minimization algorithms. Once a search direction is obtained, then one performs the line search procedure. Before we discuss these aspects in detail we turn to the question of the convergence rate of an algorithm.

### 4.2 Rate of convergence

Assume that an algorithm generates an $n$ dimensional convergent sequence of iterates $x^{1}, x^{2}, \ldots, x^{k}, \ldots \rightarrow \bar{x}$, as a minimizing sequence of the continuous function $f(x)$ : $\mathbb{R}^{n} \rightarrow \mathbb{R}$.

One can define a scalar sequence $\alpha_{k}=\left\|x^{k}-\bar{x}\right\|$ with limit $\alpha=0$, or a sequence $\alpha_{k}=f\left(x^{k}\right)$ with limit $\alpha=f(\bar{x})$. The rate of convergence of these sequences gives an indication of 'how fast' the iterates converge. In order to quantify the concept of rate (or order) of convergence, we need the following definition.

Definition 4.1 Let $\alpha_{1}, \alpha_{2}, \ldots, \alpha_{k}, \ldots \rightarrow \alpha$ be a convergent sequence with $\alpha_{k} \neq \alpha$ for all $k$. We say that the order of convergence of this sequence is $p^{*}$ if

$$
p^{*}=\sup \left\{p: \limsup _{k \rightarrow \infty} \frac{\left|\alpha_{k+1}-\alpha\right|}{\left|\alpha_{k}-\alpha\right|^{p}}<\infty\right\}
$$

The larger $p^{*}$ is, the faster the convergence. Let $\beta=\limsup _{k \rightarrow \infty} \frac{\left|\alpha_{k+1}-\alpha\right|}{\left|\alpha_{k}-\alpha\right|^{p^{*}}}$. If $p^{*}=1$ and $0<\beta<1$ we are speaking about linear (or geometric rate of) convergence. If $p^{*}=1$ and $\beta=0$ the convergence rate is super-linear, while if $\beta=1$ the convergence rate is sub-linear. If $p^{*}=2$ then the convergence is quadratic.

Exercise 4.1 Show that the sequence $\alpha_{k}=a^{k}$, where $0<a<1$ converges linearly to zero while $\beta=a$.

Exercise 4.2 Show that the sequence $\alpha_{k}=a^{\left(2^{k}\right)}$, where $0<a<1$, converges quadratically to zero.

Exercise 4.3 Show that the sequence $\alpha_{k}=\frac{1}{k}$ converges sub-linearly to zero.
Exercise 4.4 Show that the sequence $\alpha_{k}=\left(\frac{1}{k}\right)^{k}$ converges super-linearly to zero. $\triangleleft$
Exercise 4.5 Construct a sequence that converges to zero with the order of four.

### 4.3 Line search

Line search in fact means one dimensional optimization, since the function $f\left(x^{k}+\lambda s^{k}\right)$ is the function of the single variable $\lambda$. Hence our problem in this part is to find the minimum of a one dimensional function $\phi(\lambda):=f\left(x^{k}+\lambda s^{k}\right)$, or if it is differentiable one has to find a zero of its derivative $\phi^{\prime}(\lambda)$.

Exercise 4.6 Assume that $f$ is continuously differentiable, $x^{k}$ and $s^{k}$ are given, and $\lambda_{k}$ is obtained via exact line search:

$$
\lambda_{k}=\arg \min _{\lambda} f\left(x^{k}+\lambda s^{k}\right)
$$

Show that $\nabla f\left(x^{k}+\lambda_{k} s^{k}\right)^{T} s^{k}=0$.
Below we present four line search methods, that require different levels of information about $\phi(\lambda)$ :

- The Dichotomous search and Golden section methods, that use only function evaluations of $\phi$;
- bisection, that evaluates $\phi^{\prime}(\lambda)$ ( $\phi$ has to be continuously differentiable);
- Newton's method, that evaluates both $\phi^{\prime}(\lambda)$ and $\phi^{\prime \prime}(\lambda)$.


### 4.3.1 Dichotomous and Golden section search

Assume that $\phi$ is convex and has a minimizer, and that we know an interval $[a, b]$ that contains this minimizer. We wish to reduce the size of this 'interval of uncertainty' by evaluating $\phi$ at points in $[a, b]$.
Say we evaluate $\phi(\lambda)$ at two points $\bar{a} \in(a, b)$ and $\bar{b} \in(a, b)$, where $\bar{a}<\bar{b}$.
Lemma 4.2 If $\phi(\bar{a})<\phi(\bar{b})$ then the minimum of $\phi$ is contained in the interval $[a, \bar{b}]$. If $\phi(\bar{a}) \geq \phi(\bar{b})$ then the minimum of $\phi$ is contained in the interval $[\bar{a}, b]$.

Exercise 4.7 Prove Lemma 4.2.
The lemma suggest a simple algorithm to reduce the interval of uncertainty.

## Input:

$\epsilon>0$ is the accuracy parameter;
$a_{0}, b_{0}$ are given such that $\left[a_{0}, b_{0}\right]$ contains the minimizer of $\phi(\lambda)$;
For $k=0,1, \ldots$, do:
Step 1: If $\left|a_{k}-b_{k}\right|<\epsilon$ STOP.
Step 2: Choose $\bar{a}_{k} \in\left(a_{k}, b_{k}\right)$ and $\bar{b}_{k} \in\left(a_{k}, b_{k}\right)$, such that $\bar{a}_{k}<\bar{b}_{k}$;
Step 3a: If $\phi\left(\bar{a}_{k}\right)<\phi\left(\bar{b}_{k}\right)$ then the minimum of $\phi$ is contained in the interval $\left[a_{k}, \bar{b}_{k}\right]$; set $b_{k+1}=\bar{b}_{k}$ and $a_{k+1}=a_{k}$;
Step 3b: If $\phi\left(\bar{a}_{k}\right) \geq \phi\left(\bar{b}_{k}\right)$ then the minimum of $\phi$ is contained in the interval $\left[\bar{a}_{k}, b_{k}\right]$; set $a_{k+1}=\bar{a}_{k}$ and $b_{k+1}=b_{k}$;
We have not specified yet how we should choose the values $\bar{a}_{k}$ and $\bar{b}_{k}$ in iteration $k$ (Step 2 of the algorithm). There are many ways to do this. One is to choose $\bar{a}_{k}=$ $\frac{1}{2}\left(a_{k}+b_{k}\right)-\delta$ and $\bar{b}_{k}=\frac{1}{2}\left(a_{k}+b_{k}\right)+\delta$ where $\delta>0$ is a (very) small fixed constant. This is called Dichotomous Search.

Exercise 4.8 Prove that - when using Dichotomous Search - the interval of uncertainty is reduced by a factor $\left(\frac{1}{2}+\delta\right)^{t / 2}$ after $t$ function evaluations.

There is a more clever way to choose $\bar{a}_{k}$ and $\bar{b}_{k}$, which reduces the number of function evaluations per iteration from two to one, while still shrinking the interval of uncertainty by a constant factor. It is based on a geometric concept called the Golden section.
The golden section of a line segment is its division into two unequal segments, such that the ratio of the longer of the two segments to the whole segment is equal to the ratio of the shorter segment to the longer segment.


Figure 4.1: The golden section: $\alpha \approx 0.618$.

With reference to Figure 4.1, we require that the value $\alpha$ is chosen such that the following ratios are equal:

$$
\frac{1-\alpha}{\alpha}=\frac{\alpha}{1} .
$$

This is the same as $\alpha^{2}+\alpha-1=0$ which has only one root in the interval [ 0,1$]$, namely $\alpha \approx 0.618$.

Returning to the line search procedure, we simply choose $\bar{a}_{k}$ and $\bar{b}_{k}$ as the points that correspond to the golden section (see Figure 4.2).


Figure 4.2: Choosing $\bar{a}_{k}$ and $\bar{b}_{k}$ via the Golden section rule.
The reasoning behind this is as follows. Assume that we know the values $\phi\left(\bar{a}_{k}\right)$ and $\phi\left(\bar{b}_{k}\right)$ during iteration $k$. Assume that $\phi\left(\bar{a}_{k}\right)<\phi\left(\bar{b}_{k}\right)$, so that we set $b_{k+1}=\bar{b}_{k}$ and $a_{k+1}=a_{k}$. Now, by the definition of the golden section, $\bar{b}_{k+1}$ is equal to $\bar{a}_{k}$ (see Figure 4.3).

In other words, we do not have to evaluate $\phi$ at $\bar{b}_{k+1}$, because we already know this value. In iteration $k+1$ we therefore only have to evaluate $\phi\left(\bar{a}_{k+1}\right)$ in this case. The analysis for the case where $\phi\left(\bar{a}_{k}\right) \geq \phi\left(\bar{b}_{k}\right)$ is perfectly analogous.

Iteration $k+1$

Iteration $k$


Figure 4.3: Illustration of consecutive iterations of the Golden section rule when $\phi\left(\bar{a}_{k}\right)<$ $\phi\left(\bar{b}_{k}\right)$.

Exercise 4.9 Prove that - when using Golden section search - the interval of uncertainty is reduced by a factor $0.618^{t-1}$ after $t$ function evaluations.

The Golden section search requires fewer function evaluations than the Dichotomous search method to reduce the length interval of uncertainty to a given $\epsilon>0$; see Exercise 4.10. If one assumes that the time it takes to evaluate $\phi$ dominates the work per iteration, then it is more important to count the total number of function evaluations than the number of iterations.

Exercise 4.10 Show that the Dichotomous search algorithm terminates after at most

$$
2\left(\frac{\log \left(\frac{b_{0}-a_{0}}{\epsilon}\right)}{\log \left(\frac{2}{1+2 \delta}\right)}\right)
$$

function evaluations, and that the Golden section search terminates after at most

$$
1+\left(\frac{\log \left(\frac{b_{0}-a_{0}}{\epsilon}\right)}{\log \left(\frac{1}{0.618}\right)}\right)
$$

function evaluations. Which of the two bounds is better?

### 4.3.2 Bisection

The Bisection method (also called Bolzano's method) is used to find a root of $\phi^{\prime}(\lambda)$ (here we assume $\phi$ to be continuously differentiable). Recall that such a root corresponds to a minimum of $\phi$ if $\phi$ is convex.
The algorithm is similar to the Dichotomous and Golden section search ones, in the sense that it too uses an interval of uncertainty that is reduced at each iteration. In the case of the bisection method the interval of uncertainty contains a root of $\phi^{\prime}(\lambda)$.

The algorithm proceeds as follows.

## Input:

$\epsilon>0$ is the accuracy parameter;
$a_{0}, b_{0}$ are given such that $\phi^{\prime}\left(a_{0}\right)<0$ and $\phi^{\prime}\left(b_{0}\right)>0$;

For $k=0,1, \ldots$, do:
Step 1: If $\left|b_{k}-a_{k}\right|<\epsilon$ STOP.
Step 2: Let $\lambda=\frac{1}{2}\left(a_{k}+b_{k}\right)$;
Step 3: If $\phi^{\prime}(\lambda)<0$ then $a_{k+1}:=\lambda$ and $b_{k+1}=b_{k}$;
Step 4: If $\phi^{\prime}(\lambda)>0$ then $b_{k+1}:=\lambda$ and $a_{k+1}=a_{k}$.

Exercise 4.11 Prove that the bisection algorithm uses at most $\log _{2} \frac{\left|b_{0}-a_{0}\right|}{\epsilon}$ function evaluations before terminating.

Note that the function $\phi^{\prime}(\lambda)$ does not have to be differentiable in order to perform the bisection procedure.

### 4.3.3 Newton's method

Newton's method is another algorithm for finding a root of $\phi^{\prime}$. Once again, such a root corresponds to a minimum of $\phi$ if $\phi$ is convex. Newton's method requires that $\phi$ be twice continuously differentiable and strictly convex, and works as follows: we construct the linear Taylor approximation to $\phi^{\prime}$ at the current iterate $\lambda_{k}$, namely

$$
l(\lambda):=\phi^{\prime}\left(\lambda_{k}\right)+\phi^{\prime \prime}\left(\lambda_{k}\right)\left(\lambda-\lambda_{k}\right) .
$$

Next we find the root of $l(\lambda)$ and set $\lambda_{k+1}$ to be equal to this root. This means that $\lambda_{k+1}$ is given by

$$
\lambda_{k+1}=\lambda_{k}-\frac{\phi^{\prime}\left(\lambda_{k}\right)}{\phi^{\prime \prime}\left(\lambda_{k}\right)} .
$$

Now we repeat the process with $\lambda_{k+1}$ as the current iterate.
There is an equivalent interpretation of this procedure: take the quadratic Taylor approximation of $\phi$ at the current iterate $\lambda_{k}$, namely

$$
q(\lambda)=\phi\left(\lambda_{k}\right)+\phi^{\prime}\left(\lambda_{k}\right)\left(\lambda-\lambda_{k}\right)+\frac{1}{2} \phi^{\prime \prime}\left(\lambda_{k}\right)\left(\lambda-\lambda_{k}\right)^{2}
$$

and set $\lambda_{k+1}$ to be the minimum of $q$. The minimum of $q$ is attained at

$$
\lambda_{k+1}=\lambda_{k}-\frac{\phi^{\prime}\left(\lambda_{k}\right)}{\phi^{\prime \prime}\left(\lambda_{k}\right)},
$$

and $\lambda_{k+1}$ becomes the new iterate (new approximation to the minimum). Note that the two interpretations are indeed equivalent.
Newton's algorithm can be summarized as follows.

## Input:

$\epsilon>0$ is the accuracy parameter;
$\lambda_{0}$ is the given initial point; $k=0$;

For $k=0,1, \ldots$, do:
Step 1: Let $\lambda_{k+1}=\lambda_{k}-\frac{\phi^{\prime}\left(\lambda_{k}\right)}{\phi^{\prime \prime}\left(\lambda_{k}\right)}$;
Step 2: If $\left|\lambda_{k+1}-\lambda_{k}\right|<\epsilon$ STOP.
Newton's method as presented above may not converge to the global minimum of $\phi$. On the other hand, Newton's method has some spectacular properties. It converges quadratically if the following conditions are met:

1. the starting point is sufficiently close to the minimum point;
2. in addition to being convex, the function $\phi$ has a property called self-concordance, which we will discuss later.

The next two examples illustrate the possible scenarios.

Example 4.3 Let us apply Newton's method to $\phi(\lambda)=\lambda-\log (1+\lambda)$. Note that the domain of $\phi$ is $(-1, \infty)$. The first and second derivatives of $\phi$ are given by

$$
\phi^{\prime}(\lambda)=\frac{\lambda}{1+\lambda}, \phi^{\prime \prime}(\lambda)=\frac{1}{(1+\lambda)^{2}}
$$

and it is therefore clear that $\phi$ is strictly convex on its domain, and that $\lambda=0$ is the minimizer of $\phi$.
The iterates from Newton's method satisfy the recursive relation

$$
\lambda_{k+1}=\lambda_{k}-\left[\phi^{\prime \prime}\left(\lambda_{k}\right)\right]^{-1} \phi^{\prime}\left(\lambda_{k}\right)=\lambda_{k}-\lambda_{k}\left(1+\lambda_{k}\right)=-\lambda_{k}^{2}
$$

This implies quadratic convergence if $\left|\lambda_{0}\right|<1$ (see Exercise 4.12).
On the other hand, note that Newton's method fails if $\lambda_{0} \geq 1$. For example, if $\lambda_{0}=1$ then $\lambda_{1}=-1$, which is not in the domain of $\phi$ !

We mention that the convex function $\phi$ has the self-concordance property mentioned above. This will be shown in Exercise 6.16.

Exercise 4.12 This exercise refers to Example 4.3. Prove that, if the sequence $\left\{\lambda_{k}\right\}$ satisfies

$$
\lambda_{k+1}=-\left(\lambda_{k}\right)^{2}
$$

then $\lambda_{k} \rightarrow 0$ with a quadratic rate of convergence if $\left|\lambda_{0}\right|<1$.
In the following example, Newton's method converges to the minimum, but the rate of convergence is only linear.

Example 4.4 Let $m \geq 2$ be even and define

$$
\phi(\lambda)=\lambda^{m} .
$$

Clearly, $\phi$ has a unique minimizer, namely $\lambda=0$. Suppose we start Newton's method at some nonzero $\lambda_{0} \in \mathbb{R}$.

The derivatives of $\phi$ are

$$
\begin{aligned}
\phi^{\prime}(\lambda) & =m \lambda^{m-1} \\
\phi^{\prime \prime}(\lambda) & =m(m-1) \lambda^{m-2}
\end{aligned}
$$

Hence, the iterates from Newton's method satisfy the recursive relation

$$
\lambda_{k+1}=\lambda_{k}-\left(\phi^{\prime \prime}\left(\lambda_{k}\right)\right)^{-1} \phi^{\prime}\left(\lambda_{k}\right)=\lambda_{k}+\frac{-1}{m-1} \lambda_{k}=\frac{m-2}{m-1} \lambda_{k}
$$

This shows that Newton's method is exact if $\phi$ is quadratic (if $m=2$ ), whereas for $m>2$ the Newton process converges to 0 with a linear convergence rate (see Exercise 4.13).

Exercise 4.13 This exercise refers to Example 4.4. Prove that, if the sequence $\left\{\lambda_{k}\right\}$ satisfies

$$
\lambda_{k+1}=\frac{m-2}{m-1} \lambda_{k},
$$

where $m>2$ is even, then $\lambda_{k} \rightarrow 0$ with a linear rate of convergence, if $\lambda_{0} \neq 0$.

### 4.4 Gradient method

We now return to the generic algorithm on page 77, and look at some different choices for the search direction. The gradient method uses the negative gradient $\left(-\nabla f\left(x^{k}\right)\right)$ of the function $f$ as the search direction. ${ }^{1}$ This direction is frequently referred to as the steepest descent direction. This name is justified by observing that the normalized directional derivative is minimized by the negative gradient

$$
\delta f(x,-\nabla f(x))=-\nabla f(x)^{T} \nabla f(x)=\min _{\|s\|=\|\nabla f(x)\|}\left\{\nabla f(x)^{T} s\right\} .
$$

Exercise 4.14 Let $f: \mathbb{R}^{n} \mapsto \mathbb{R}$ be continuously differentiable and let $\bar{x} \in \mathbb{R}^{n}$ be given. Assume that the level set $\left\{x \in \mathbb{R}^{n} \mid f(x)=f(\bar{x})\right\}$, is in fact a curve (contour). Show that $\nabla f(\bar{x})$ is orthogonal to the tangent line to the curve at $\bar{x}$.

To calculate the gradient is relatively cheap which indicates that the gradient method can be quite efficient. Although it works fine in many applications, several theoretical and practical disadvantages can be mentioned. First, the minimization of a convex quadratic function by the gradient method is not a finite process in general. Slow convergence, due to a sort of "zigg-zagging" sometimes takes place. Secondly, the order of convergence is no better than linear in general.
Figure 4.4 illustrates the zig-zag behavior that may occur when using the gradient method.

Exercise 4.15 Calculate the steepest descent direction for the quadratic function

$$
f(x)=\frac{1}{2} x^{T} Q x+q^{T} x-\beta,
$$

where the matrix $Q$ is positive definite. Calculate the exact step length in the line search as well.

[^2]

Figure 4.4: Iterates of the gradient method for the function $f(x)=9 x_{1}^{2}+2 x_{1} x_{2}+x_{2}^{2}$.

Exercise 4.16 Prove that subsequent search directions of the gradient method are always orthogonal (i.e. $s^{k} \perp s^{k+1} ; k=0,1,2, \ldots$ ) if exact line search is used.

The following theorem gives a convergence result for the gradient method.
Theorem 4.5 Let $f$ be continuously differentiable. Starting from the initial point $x^{0}$ using exact line search the gradient method produces a decreasing sequence $x^{0}, x^{1}, x^{2}, \ldots$ such that $f\left(x^{k}\right)>f\left(x^{k+1}\right)$ for $k=0,1,2, \cdots$. Assume that the level set $D=\{x$ : $\left.f(x) \leq f\left(x^{0}\right)\right\}$ is compact, then any accumulation point $\bar{x}$ of the generated sequence $x^{0}, x^{1}, x^{2}, \cdots, x^{k}, \cdots$ is a stationary point (i.e. $\nabla f(\bar{x})=0$ ) of $f$. Further if the function $f$ is a convex function, then $\bar{x}$ is a global minimizer of $f$.

Proof: Since $D$ is compact and $f$ is continuous we have that $f$ is bounded on $D$, hence we have a convergent subsequence $x^{k_{j}} \rightarrow \bar{x}$ with $f\left(x^{k_{j}}\right) \rightarrow f^{*}$ as $k_{j} \rightarrow \infty$. By continuity of $f$ we have $f(\bar{x})=f^{*}$. Since the search direction is the gradient of $f$ we have

$$
\bar{s}=\lim _{k_{j} \rightarrow \infty} s^{k_{j}}=-\lim _{k_{j} \rightarrow \infty} \nabla f\left(x^{k_{j}}\right)=-\nabla f(\bar{x}) .
$$

Multiplying by $\nabla f(\bar{x})$ we have

$$
\begin{equation*}
\bar{s}^{T} \nabla f(\bar{x})=-\nabla f(\bar{x})^{T} \nabla f(\bar{x}) \leq 0 . \tag{4.2}
\end{equation*}
$$

On the other hand using the construction of the iteration sequence and the convergent subsequence we write

$$
f\left(x^{k_{j+1}}\right) \leq f\left(x^{k_{j}+1}\right) \leq f\left(x^{k_{j}}+\lambda s^{k_{j}}\right)
$$

Taking the limit in the last inequality we have

$$
f(\bar{x}) \leq f(\bar{x}+\lambda \bar{s})
$$

which leads to $\delta f(\bar{x}, \bar{s})=\bar{s}^{T} \nabla f(\bar{x}) \geq 0$. Combining this result with (4.2) we have $\nabla f(\bar{x})=0$, and the theorem is proved.

### 4.5 Newton's method

We now extend Newton's method to the multivariate case. To apply Newton's method we have to assume that the function $f$ is a twice continuously differentiable function with positive definite Hessian on its domain. Newton's search direction in multidimensional optimization is again based on minimizing the second order approximation of the function $f$. The quadratic Taylor approximation at the current iterate $x^{k}$ is given by:

$$
q(x):=f\left(x^{k}\right)+\nabla f\left(x^{k}\right)^{T}\left(x-x^{k}\right)+\frac{1}{2}\left(x-x^{k}\right)^{T} \nabla^{2} f\left(x^{k}\right)\left(x-x^{k}\right)
$$

Since the Hessian $\nabla^{2} f\left(x^{k}\right)$ is positive definite, the function $q(x)$ is strictly convex (see Exercise 1.20). Hence the minimum of $q(x)$ is attained when its gradient

$$
\nabla q(x)=\nabla f\left(x^{k}\right)+\nabla^{2} f\left(x^{k}\right)\left(x-x^{k}\right)
$$

equals to the zero vector, i.e. at the point

$$
x^{k+1}=x^{k}-\left(\nabla^{2} f\left(x^{k}\right)\right)^{-1} \nabla f\left(x^{k}\right)
$$

The classical Newton method does not apply line search, one takes the full Newton step. If line search is applied then typically we are far from the solution, the step length is usually less than one. We refer to this as the damped Newton method.

In addition we have to mention that to compute and invert the Hesse matrix is more expensive than to compute only the gradient. Several methods are developed to reduce this cost while preserving the advantages of Newton's method. These are the so-called quasi-Newton methods of which the most popular are the methods which use conjugate directions, to be discussed later.

Anyway, the compensation for the extra cost in Newton's method is a better search direction. Just note that the minimization of a convex quadratic function happens in one step.

Exercise 4.17 Let $f(x)=\frac{1}{2} x^{T} A x-b^{T} x$ where $A$ is positive definite and $b \in \mathbb{R}^{n}$. Assume that we apply Newton's method to minimize $f$. Show that $x^{1}=A^{-1} b$, i.e. $x^{1}$ is the minimum of $f$, regardless of the starting point $x^{0}$.

If the Hessian $\nabla^{2} f(x)$ is not positive definite, or is ill-conditioned (the ratio of the largest and smallest eigenvalue is large) then it is not (or hardly) invertible. In this case additional techniques are needed to circumvent these difficulties. In the trust region method, $\nabla^{2} f(x)$ is replaced by $\left(\nabla^{2} f(x)+\alpha I\right)$ where $I$ is the identity matrix and $\alpha$ is changed dynamically. Observe that if $\alpha=0$ then we have the Hessian, hence we have the Newton step, while as $\alpha \rightarrow \infty$ this matrix approaches a multiple of the identity matrix and so the search direction is asymptotically getting parallel to the negative gradient.
The interested reader can consult the following books for more details on trust region methods $[2,3,16,9]$.

Exercise 4.18 Let $x \in \mathbb{R}^{n}$ and $f$ be twice continuously differentiable. Show that $s=$ $-H \nabla f(x)$ is a descent direction of $f$ at $x$ for any positive definite matrix $H$, if $\nabla f(x) \neq 0$. Which choice of $H$ gives:

- the steepest descent direction?
- Newton's direction (for convex f)?

Exercise 4.19 Consider the unconstrained optimization problem:

$$
\min \left(x_{1}-2\right)^{4}+\left(x_{1}-2 x_{2}\right)^{2}
$$



Figure 4.5: Contours of the function. Note that the minimum is at $[2,1]^{T}$.

1. Perform two iterations of the gradient method, starting from $x^{0}=[0,3]^{T}$.
2. Perform four iterations of Newton's method (without line search), with the same starting point $x^{0}$.

## Relation with Newton's method for solving nonlinear equations

The reader may be familiar with Newton's method to solve nonlinear systems of equations. Here we show that Newton's optimization method is obtained by setting the gradient of $f$ to zero and using Newton's method for nonlinear equations to solve the resulting equations.

Assume we have a nonlinear system of equations

$$
F(x)=0
$$

to solve, where $F(x)$ is a differentiable mapping from $\mathbb{R}^{n} \rightarrow \mathbb{R}^{m}$. Given any point $x^{k} \in \mathbb{R}^{n}$, Newton's method proceeds as follows. Let us first linearize the nonlinear equation at $x^{k}$ by approximating $F(x)$ by $F\left(x^{k}\right)+J F\left(x^{k}\right)\left(x-x^{k}\right)$ where $J F(x)$ denotes the Jacobian of $F$, an $m \times n$ matrix defined as

$$
J F(x)_{i j}=\frac{\partial F_{i}(x)}{\partial x_{j}} \quad \text { where } \quad i=1, \cdots, m ; \quad j=1, \cdots, n
$$

Now we take a step so that the iterate after the step satisfies the linearized equation

$$
\begin{equation*}
J F\left(x^{k}\right)\left(x^{k+1}-x^{k}\right)=-F\left(x^{k}\right) \tag{4.3}
\end{equation*}
$$

This is a linear system of equations, hence a solution (if it exists) can be found by standard linear algebra.

Observe, that if we want to minimize a strictly convex function $f(x)$ one can interpret this problem as solving the nonlinear system of equations $\nabla f(x)=0$. The solution of this system by Newton's method, as we have a point $x^{k}$, leads to (apply (4.3))

$$
\nabla^{2} f\left(x^{k}\right)\left(x^{k+1}-x^{k}\right)=-\nabla f\left(x^{k}\right)
$$

The Jacobian of the gradient is exactly the Hessian of the function $f(x)$ hence it is positive definite and we have

$$
x^{k+1}=x^{k}-\left(\nabla^{2} f\left(x^{k}\right)\right)^{-1} \nabla f\left(x^{k}\right)
$$

as we have seen above.

### 4.6 Methods of Conjugate directions

Let $A$ be an $n \times n$ symmetric positive definite matrix and $b \in \mathbb{R}^{n}$. We consider the problem of minimizing the strictly convex quadratic function

$$
q(x)=\frac{1}{2} x^{T} A x-b^{T} x
$$

We will study a class of algorithms that use so-called conjugate search directions to minimize $q$.

Definition 4.6 The directions (vectors) $s^{1}, \cdots, s^{k} \in \mathbb{R}^{n}$ are called conjugate (or $A$-conjugate) directions if $\left(s^{i}\right)^{T} A s^{j}=0$ for all $1 \leq i \neq j \leq k$.

Note that conjugate directions are mutually orthogonal if $A=I$.

Exercise 4.20 Let $A$ be $n \times n$ symmetric positive definite and $s^{1}, \ldots, s^{k}(k \leq n)$ be $A$ conjugate. Prove that $s^{1}, \ldots, s^{k}$ are linearly independent.

If one uses $A$-conjugate directions in the generic algorithm to minimize $q$, then the minimum is found in at most $n$ iterations. The next theorem establishes this important fact.

Theorem 4.7 Let $s^{0}, \cdots, s^{k} \in \mathbb{R}^{n}$ be conjugate directions with respect to $A$. Let $x^{0}$ be given and let

$$
x^{i+1}:=\arg \min \quad q\left(x^{i}+\lambda s^{i}\right) \quad i=0, \cdots, k .
$$

Then $x^{k+1}$ minimizes $q(x)$ on the affine space $H=x^{0}+\operatorname{span}\left(s^{0}, \cdots, s^{k}\right)$.

Proof: One has to show (see Theorem 2.9) that $\nabla q\left(x^{k+1}\right) \perp s^{1}, \cdots, s^{k}$. Recall that

$$
x^{i+1}:=x^{i}+\lambda^{i} s^{i} \quad i=0, \cdots, k
$$

where $\lambda^{i}$ indicates the line-search minimum, thus

$$
x^{k+1}:=x^{1}+\lambda^{0} s^{0}+\cdots+\lambda^{k} s^{k}=x^{i}+\lambda^{i} s^{i}+\cdots+\lambda^{k} s^{k},
$$

for any fixed $i \leq k$. Due to exact line-search we have $\nabla q\left(x^{i+1}\right)^{T} s^{i}=0$ (see Exercise 4.6). Using $\nabla q(x)=A x-b$, we get

$$
\nabla q\left(x^{k+1}\right):=\nabla q\left(x^{i}+\lambda^{i} s^{i}\right)+\sum_{j=i+1}^{k} \lambda^{j} A s^{j} .
$$

Taking the inner product on both sides with $s^{i}$ yields

$$
\left(s^{i}\right)^{T} \nabla q\left(x^{k+1}\right):=\left(s^{i}\right)^{T} \nabla q\left(x^{i+1}\right)+\sum_{j=i+1}^{k} \lambda^{j}\left(s^{i}\right)^{T} A s^{j} .
$$

Hence $\left(s^{i}\right)^{T} \nabla q\left(x^{k+1}\right)=0$.

Corollary 4.8 Let $x^{k}$ be defined as in Theorem 4.7. Then $x^{n}=A^{-1} b$, i.e. $x^{n}$ is the minimizer of $q(x)=\frac{1}{2} x^{T} A x-b^{T} x$.

Exercise 4.21 Show that the result in Corollary 4.8 follows from Theorem 4.7.

### 4.6.1 The method of Powell

To formulate algorithms that use conjugate directions, we need tools to construct conjugate directions. The next theorem may seem a bit technical, but it gives us such a tool.

Theorem 4.9 Let $\mathcal{L}$ be a linear subspace, $H_{1}:=x^{1}+\mathcal{L}$ and $H_{2}:=x^{2}+\mathcal{L}$ be two parallel affine spaces where $x^{1}$ and $x^{2}$ are the minimizers of $q(x)$ over $H_{1}$ and $H_{2}$, respectively.

Then for every $s \in \mathcal{L},\left(x^{2}-x^{1}\right)$ and $s$ are conjugate with respect to $A$.

Proof: Assume $x^{1}$ minimizes $q(x)$ over $H_{1}=x^{1}+\mathcal{L}$ and $x^{2}$ minimizes $q(x)$ over $H_{2}=x^{2}+\mathcal{L}$. Let $s \in \mathcal{L}$. Now

$$
\begin{aligned}
& q\left(x^{1}+\lambda s\right) \geq q\left(x^{1}\right) \quad \Rightarrow \quad s^{T} \nabla q\left(x^{1}\right)=0 \\
& q\left(x^{2}+\lambda s\right) \geq q\left(x^{2}\right) \Rightarrow s^{T} \nabla q\left(x^{2}\right)=0
\end{aligned}
$$

This implies that

$$
s^{T}\left(\nabla q\left(x^{2}\right)-\nabla q\left(x^{1}\right)\right)=s^{T} A\left(x^{2}-x^{1}\right)=0
$$

In other words, for any $s \in \mathcal{L}, s$ and $x^{2}-x^{1}$ are $A$-conjugate directions.

The basic ingredients of the method of Powell are as follows:

- The algorithm constructs conjugate directions $t^{1}, \ldots, t^{n}$ by using the result of Theorem 4.9. The method requires one cycle of $n+1$ line searches to construct each conjugate direction $t^{i}$. Thus the first conjugate direction $t^{1}$ is constructed at the end of cycle 1, etc.
- It starts with a fixed set of linearly independent directions $s^{1}, \ldots, s^{n}$ to achieve this. (Usually the standard unit vectors.)
- In the first cycle, the method performs successive exact line searches using the directions $s^{1}, \ldots, s^{n}$ (in that order). In the second cycle the directions $s^{2}, \ldots, s^{n}, t^{1}$ are used (in that order). In the third cycle the directions $s^{3}, \ldots, s^{n}, t^{1}, t^{2}$ are used, etc.
- The method terminates after $n$ cycles due to the result in Theorem 4.7.

We will now state the algorithm, but first a word about notation. As mentioned before, the second cycle uses the directions $s^{2}, \ldots, s^{n}, t^{1}$. In order to state the algorithm in a compact way, the search directions used during cycle $k$ are called $s^{(k, 1)}, \ldots, s(k, n)$.

The iterates generated during cycle $k$ via successive line searches will be called $z^{(k, 1)}, \ldots, z^{(k, n)}$, and $x^{k}$ will denote the iterate at the end of cycle $k$.

## Powell's algorithm

Input A starting point $x^{0}$, a set of linearly independent vectors $s^{1}, \ldots, s^{n}$.
Initialization Set $s^{(1, i)}=s^{i}, \quad i=1, \cdots, n$.

For $k=1,2, \ldots, n$ do:
(Cycle $k$ :)
Let $z^{(k, 1)}=x^{k-1}$ and $z^{(k, i+1)}:=\arg \min q\left(z^{(k, i)}+\lambda s^{(k, i)}\right) \quad i=1, \cdots, n$.
Let $x^{k}:=\operatorname{argmin} q\left(z^{(k, n+1)}+\lambda t^{k}\right)$ where $t^{k}:=z^{(k, n+1)}-x^{k-1}$.
Let $s^{(k+1, i)}=s^{(k, i+1)}, \quad i=1, \cdots, n-1$ and $s^{(k+1, n)}:=t^{k}$.

It may not be clear to the reader why the directions $t^{1}, t^{2}, \ldots$ are indeed conjugate directions. As mentioned before, we will invoke Theorem 4.9 to prove this.

Lemma 4.10 The vectors $t^{1}, \ldots, t^{n}$ generated by Powell's algorithm are $A$-conjugate.

Proof: The proof is by induction. Assume that $t^{1}, \ldots, t^{k}$ are conjugate at the end of cycle $k$ of the algorithm. By the definition of $x^{k}$ in the statement of the algorithm, and by Theorem 4.7, $x^{k}$ minimizes $q$ over the affine space $x^{k}+\operatorname{span}\left\{t^{1}, \ldots, t^{k}\right\}$.
In cycle $k+1, z^{(k+1, n+1)}$ is obtained after successive line searches along the directions $\left\{s^{1}, \ldots, s^{n-k}, t^{1}, \ldots, t^{k}\right\}$. By Theorem 4.7, $z^{(k+1, n+1)}$ minimizes $q$ over the affine space $z^{(k+1, n)}+\operatorname{span}\left\{t^{1}, \ldots, t^{k}\right\}$.
Now define $t^{k+1}=z^{(k, n+1)}-x^{k-1}$. By Theorem 4.9, $t^{k+1}$ is $A$-conjugate to every vector in $\operatorname{span}\left\{t^{1}, \ldots, t^{k}\right\}$, and in particular to $\left\{t^{1}, \ldots, t^{k}\right\}$.

Example 4.11 We consider the problem

$$
\min f(x)=5 x_{1}^{2}+2 x_{1} x_{2}+x_{2}^{2}+7 .
$$

The minimum is attained at $x_{1}=x_{2}=0$.
We choose $s^{1}$ and $s^{2}$ as the standard unit vectors in $\mathbb{R}^{2}$, and the starting point is: $x^{0}=[1,2]^{T}$. The progress of Powell's method for this example is illustrated in Figure 4.6. We will describe the progress with giving the actual numerical values, in order to keep things simple.
Note that, at the start of cycle 1, successive line searches are done using $s^{1}=\left[\begin{array}{ll}01\end{array}\right]^{T}$ and $s^{2}=[10]^{T}$. Then the first conjugate direction $t^{1}$ is generated by connecting $x^{0}$ with the last point obtained, and a line search is performed along $t^{1}$ to obtain the point $x^{1}$.
In cycle 2 , successive line searches are done using $s^{2}$ and $t^{1}$. Then the second conjugate direction $t^{2}$ is generated by connecting $x^{1}$ with the last point obtained, and a line search is performed along $t^{2}$ to obtain the point $x^{2}$.
Note that $x^{2}$ is optimal, as it should be.


Figure 4.6: Iterates generated by Powell's algorithm for the function $f(x)=5 x_{1}^{2}+$ $2 x_{1} x_{2}+x_{2}^{2}+7$, starting from $x^{0}=[1,2]^{T}$.

## Discussion of Powell's method

- We may apply Powell's algorithm to any function (not necessarily quadratic); The only change to the Algorithm on page 91 is that the quadratic function $q(x)$ is replaced by a general $f(x)$. Of course, in this case it does not make sense to speak of conjugate directions, and there is no guarantee that $x^{n}$ will be the optimal solution. For this reason it is customary to restart the algorithm from $x^{n}$.
- Powell's algorithm uses only line searches, and finds the exact minimum of a strictly convex quadratic function after at most $n(n+1)$ line-searches. For a general (convex) function $f$, Powell's method can be combined with the Golden section line search procedure to obtain an algorithm for minimizing $f$ that does not require gradient information.
- Storage requirements: The algorithm stores $n n$-vectors (the current set of search directions) at any given time.

Let us compare Powell's method to Newton's method and the gradient method. Newton's method requires only one step to minimize a strictly convex quadratic function, but requires both gradient and Hessian information for general functions. The gradient method requires only gradient information, but does not always converge in a finite number of steps (not even for strictly convex quadratic functions).
In conclusion, Powell method is an attractive algorithm for minimizing 'black box' functions where gradient and Hessian information is not available (or too expensive to compute).

### 4.6.2 The Fletcher-Reeves method

The method of Fletcher and Reeves is also a conjugate gradient method to

$$
\operatorname{minimize} q(x)=\frac{1}{2} x^{T} A x-b^{T} x
$$

but is simpler to state than the method of Powell.
Before giving a formal statement of the algorithm, we list the key ingredients:

- The first search direction is the steepest descent direction: $s^{0}=-\nabla q\left(x^{0}\right)$.
- The search direction at iteration $k$, namely $s^{k}$, is constructed so that it is conjugate with respect to the preceding directions $s^{0}, \ldots, s^{k-1}$, as well as a linear combination of $-\nabla q\left(x^{k}\right)$ and $s^{0}, \ldots, s^{k-1}$.
- We will show that these requirements imply that

$$
s^{k}=-\nabla q\left(x^{k}\right)+\left(\frac{\left\|\nabla q\left(x^{k}\right)\right\|^{2}}{\left\|\nabla q\left(x^{k-1}\right)\right\|^{2}}\right) s^{k-1} .
$$

- Note that, unlike Powell's method, this method requires gradient information. The advantage over Powell's method is that we only have to store two $n$-vectors and do $n+1$ line searches.
- We may again use the method to minimize non-quadratic functions, but then convergence is not assured.

Let us consider the situation during iteration $k$, i.e. assume that $x^{k}, \nabla q\left(x^{k}\right)$ and $s^{1}, \cdots, s^{k-1}$ conjugate directions be given.

We want to find values $\beta_{k, 0} \ldots \beta_{k, k-1}$ such that

$$
s^{k}:=-\nabla q\left(x^{k}\right)+\beta_{k, 0} s^{0}+\cdots+\beta_{k, k-1} s^{k-1}
$$

and $s^{k}$ is conjugate with respect to $s^{0}, \cdots, s^{k-1}$.
We require $A$-conjugacy, i.e. $s_{i}^{T} A s_{k}=0$, which implies:

$$
\beta_{k, i}=\frac{\nabla q\left(x^{k}\right)^{T} A s^{i}}{\left(s^{i}\right)^{T} A s^{i}} \quad(i=0, \ldots, k-1)
$$

We will now show that $\beta_{k, i}=0$ if $i<k-1$. To this end, note that

$$
\nabla q\left(x^{i+1}\right)-\nabla q\left(x^{i}\right)=A\left(x_{i+1}-x_{i}\right)=\lambda_{i} A s^{i} .
$$

Therefore

$$
\beta_{k, i}=\frac{\nabla q\left(x^{k}\right)^{T}\left(\nabla q\left(x^{i+1}\right)-\nabla q\left(x^{i}\right)\right)}{\left(s^{i}\right)^{T}\left(\nabla q\left(x^{i+1}\right)-\nabla q\left(x^{i}\right)\right)} \quad(i<k) .
$$

For any $i<k$ we have

$$
s^{i}=-\nabla q\left(x^{i}\right)+\beta_{i, 1} s^{1}+\cdots+\beta_{i, i-1} s^{i-1} .
$$

By Theorem 4.7 we have

$$
\nabla q\left(x^{k}\right) \perp s^{i} \quad(i=0, \ldots, k-1) .
$$

Therefore

$$
\nabla q\left(x^{i}\right)^{T} \nabla q\left(x^{k}\right)=0 \quad(i<k),
$$

and

$$
\nabla q\left(x^{i}\right)^{T} s^{i}=-\left\|\nabla q\left(x^{i}\right)\right\|^{2} \quad(i<k) .
$$

Therefore $\beta_{k, i}=0$ if $i<k-1$. Also, due to exact line-search, we have $\left(s^{i}\right)^{T}\left(\nabla q\left(x^{i+1}\right)\right)=$ 0 (see Exercise 4.6). Therefore

$$
\beta_{k, k-1}=\frac{\left\|\nabla q\left(x^{k}\right)\right\|^{2}}{\left\|\nabla q\left(x^{k-1}\right)\right\|^{2}} .
$$

## Fletcher-Reeves algorithm

Let $x^{0}$ be an initial point.
Step 0. Let $s^{0}=-\nabla q\left(x^{0}\right)$ and $x^{1}:=\arg \min q\left(x^{0}+\lambda s^{0}\right)$.

Step $k$. Let $x^{k}, \nabla q\left(x^{k}\right)$ and $s^{0}, \cdots, s^{k-1}$ conjugate directions be given. Set

$$
s^{k}=-\nabla q\left(x^{k}\right)+\left(\frac{\left\|\nabla q\left(x^{k}\right)\right\|^{2}}{\left\|\nabla q\left(x^{k-1}\right)\right\|^{2}}\right) s^{k-1} .
$$

Set $x^{k+1}:=\operatorname{argmin} q\left(x^{k}+\lambda s^{k}\right)$.

## Exercise 4.22

$$
\min x_{1}^{2}+2 x_{2}^{2}+2 x_{3}^{2}+2 x_{1} x_{2}+2 x_{2} x_{3} .
$$

1. Solve this problem using the conjugate gradient method of Powell. Use exact line search and the starting point $[2,4,10]^{T}$. Use the standard unit vectors as $s^{1}, s^{2}$ and $s^{3}$.
2. Solve this problem using the Fletcher-Reeves conjugate gradient method. Use exact line search and the starting point $[2,4,10]^{T}$.

### 4.7 Quasi-Newton methods

Recall that the Newton direction at iteration $k$ is given by:

$$
s^{k}=-\left[\nabla^{2} f\left(x^{k}\right)\right]^{-1} \nabla f\left(x^{k}\right) .
$$

Quasi-Newton methods use a positive definite approximation $H_{k}$ to $\left[\nabla^{2} f\left(x^{k}\right)\right]^{-1}$. The approximation $H_{k}$ is updated at each iteration, say

$$
H_{k+1}=H_{k}+D_{k},
$$

where $D_{k}$ denotes the update.
Let $A$ be an $n \times n$ symmetric PD matrix, and consider once more the strictly convex quadratic function

$$
q(x)=\frac{1}{2} x^{T} A x-b^{T} x
$$

The Newton direction for $q$ at $x^{k}$ is:

$$
s^{k}=-\left[\nabla^{2} q\left(x^{k}\right)\right]^{-1} \nabla q\left(x^{k}\right)=-A^{-1} \nabla q\left(x^{k}\right)
$$

Note that

$$
\nabla q\left(x^{k+1}\right)-\nabla q\left(x^{k}\right)=A\left(x^{k+1}-x^{k}\right)
$$

We introduce the notation

$$
y^{k}:=\nabla q\left(x^{k+1}\right)-\nabla q\left(x^{k}\right), \quad \sigma^{k}=x^{k+1}-x^{k} .
$$

Notice that $y^{k}=A \sigma^{k}$ i.e. $\sigma^{k}=A^{-1} y^{k}$.

## The secant condition

We will use a search direction of the form

$$
s^{k}=-H_{k} \nabla q\left(x^{k}\right)
$$

where $H_{k}$ is an approximation of $\left[\nabla^{2} q\left(x^{k}\right)\right]^{-1}=A^{-1}$, and subsequently perform the usual line search:

$$
x^{k+1}=\arg \min q\left(x^{k}+\lambda s^{k}\right)
$$

Since

$$
y^{k}:=\nabla q\left(x^{k+1}\right)-\nabla q\left(x^{k}\right), \quad \sigma^{k}=x^{k+1}-x^{k}
$$

and $\sigma^{k}=A^{-1} y^{k}$, we require that $\sigma^{k}=H_{k+1} y^{k}$. This is called the secant condition (quasi-Newton property).

## The hereditary property

Since

$$
y^{k}:=\nabla q\left(x^{k+1}\right)-\nabla q\left(x^{k}\right), \quad \sigma^{k}=x^{k+1}-x^{k}
$$

and $\nabla q(x)=A x-b$, it holds that

$$
\sigma^{i}=A^{-1} y^{i} \quad(i=0, \ldots, k-1)
$$

We therefore require that our approximation $H_{k}$ also satisfies

$$
\sigma^{i}=H_{k} y^{i} \quad(i=0, \ldots, k-1)
$$

This is called the hereditary property.
Since

$$
\sigma^{i}=A^{-1} y^{i} \text { and } \sigma^{i}=H_{n} y^{i} \quad(i=0, \ldots, n-1)
$$

it follows that $H_{n} A \sigma^{i}=\sigma^{i}(i=0, \ldots, n-1)$. If the $\sigma^{i}(i=0, \ldots, n-1)$ are linearly independent, this implies $H_{n}=A^{-1}$.

## Discussion

We showed that - if the $\sigma^{i}(i=0, \ldots, n-1)$ are linearly independent - we have $H_{n}=A^{-1}=\left[\nabla^{2} q\left(x^{n}\right)\right]^{-1}$ (the approximation has become exact!) In iteration $n$, we therefore use the search direction

$$
s^{n}=-H_{n} \nabla q\left(x^{n}\right)=-A^{-1} \nabla q\left(x^{n}\right)
$$

But this is simply the Newton direction at $x^{n}$ ! In other words, we find the minimum of $q$ no later than in iteration $n$.

## Generic Quasi-Newton algorithm

Step 0: Let $x^{0}$ be given and set $H_{0}=I$.
Step $k$ : Calculate the search direction $s^{k}=-H_{k} \nabla q\left(x^{k}\right)$ and perform the usual line search $x^{k+1}=\arg \min q\left(x^{k}+\lambda s^{k}\right)$.

We choose $D_{k}$ in such a way that:
i $H_{k+1}=H_{k}+D_{k}$ is symmetric positive definite;
ii $\sigma^{k}=H_{k+1} y^{k}$ (secant condition);
iii $\sigma^{i}=H_{k+1} y^{i}(i=0, \ldots, k-1)$ (hereditary property).

### 4.7.1 The DFP update

The Davidon-Fletcher-Powell (DFP) rank-2 update is defined by

$$
D_{k}=\frac{\sigma^{k} \sigma^{k^{T}}}{\sigma^{k^{T}} y^{k}}-\frac{H_{k} y^{k} y^{k T} H_{k}}{y^{k^{T}} H_{k} y^{k}} .
$$

We will show that:
i If $y_{k}^{T} \sigma_{k}>0$, then $H_{k+1}$ is positive definite.
ii $H_{k+1}=H_{k}+D_{k}$ satisfies the secant condition: $\sigma^{k}=H_{k+1} y^{k}$.
iii The hereditary property holds: $\sigma^{i}=H_{k+1} y^{i}(i=0, \ldots, k-1)$.
Exercise 4.23 Show that, if $H_{k}$ is positive definite, then

$$
H_{k+1}=H_{k}+D_{k}=H_{k}+\frac{\sigma^{k} \sigma^{k^{T}}}{\sigma^{k^{T}} y^{k}}-\frac{H_{k} y^{k} y^{k T} H_{k}}{y^{k^{T}} H_{k} y^{k}},
$$

is also positive definite if $\left(\sigma^{k}\right)^{T} y^{k}>0$.
Hint 1: For ease of notation, show that

$$
H+\frac{\sigma \sigma^{T}}{\sigma^{T} y}-\frac{H y y^{T} H}{y^{T} H y}
$$

is positive definite if the matrix $H$ is P.D. and the vectors $y, \sigma$ satisfy $y^{T} \sigma>0$.
Hint 2: Set $H=L L^{T}$ and show that

$$
v^{T}\left(H+\frac{\sigma \sigma^{T}}{\sigma^{T} y}-\frac{H y y^{T} H}{y^{T} H y}\right) v>0 \quad \forall v \in \mathbb{R}^{n} \backslash\{0\}
$$

Hint 3: Use the Cauchy-Schwartz inequality

$$
\left(a^{T} a\right)\left(b^{T} b\right)-\left(a^{T} b\right)^{2}>0 \text { if } a \neq k b \text { for all } k \in \mathbb{R},
$$

to obtain the required inequality.

Exercise 4.24 Prove that $H_{k+1}=H_{k}+D_{k}$ satisfies the secant condition: $\sigma^{k}=H_{k+1} y^{k}$. $\triangleleft$

We now prove that the DFP update satisfies the hereditary property. At the same time, we will show that the search directions of the DFP method are conjugate.

Lemma 4.12 Let $H_{0}=I$. One has

$$
\begin{equation*}
\sigma^{i}=H_{k+1} y^{i} \quad(i=0, \ldots, k), k \geq 0 \tag{4.4}
\end{equation*}
$$

and $\sigma^{0}, \ldots, \sigma^{k}$ are mutually conjugate.

Proof: We will use induction on $k$. The reader may verify that (4.4) holds for $k=0$. Induction assumption:

$$
\sigma^{i}=H_{k} y^{i} \quad(i=0, \ldots, k-1)
$$

and $\sigma^{0}, \ldots, \sigma^{k-1}$ are mutually conjugate.
We now use

$$
\sigma^{k}=\lambda_{k} s^{k}=-\lambda_{k} H_{k} \nabla q\left(x^{k}\right),
$$

to get

$$
\begin{aligned}
\left(\sigma^{k}\right)^{T} A \sigma^{i} & =-\lambda_{k}\left(H_{k} \nabla q\left(x^{k}\right)\right)^{T} A \sigma^{i} \\
& =-\lambda_{k} \nabla q\left(x^{k}\right)^{T} H_{k} A \sigma^{i}
\end{aligned}
$$

Now use the induction assumption $\sigma^{i}=H_{k} y^{i} \equiv H_{k} A \sigma^{i}(i=0, \ldots, k-1)$, to get:

$$
\left(\sigma^{k}\right)^{T} A \sigma^{i}=\nabla q\left(x^{k}\right)^{T} \sigma^{i} \quad(i=0, \ldots, k-1) .
$$

Since $\sigma^{0}, \ldots, \sigma^{k-1}$ mutually conjugate, Theorem 4.7 implies that:

$$
\nabla q\left(x^{k}\right)^{T} \sigma^{i}=0 \quad(i=0, \ldots, k-1)
$$

Substituting we get

$$
\left(\sigma^{k}\right)^{T} A \sigma^{i}=0 \quad(i=0, \ldots, k-1)
$$

i.e. $\sigma^{0}, \ldots, \sigma^{k}$ are mutually conjugate. We use this to prove the hereditary property. Note that

$$
H_{k+1} y^{i}=H_{k} y^{i}+\frac{\sigma^{k} \sigma^{k} y^{i}}{\sigma^{k^{T}} y^{k}}-\frac{H_{k} y^{k} y^{k^{T}} H_{k} y^{i}}{y^{k^{T}} H_{k} y^{k}} .
$$

We can simplify this, using:

$$
\sigma^{k^{T}} y^{i}=\sigma^{k^{T}} A \sigma^{i}=0 \quad(i=0, \ldots, k-1)
$$

We get

$$
\begin{equation*}
H_{k+1} y^{i}=H_{k} y^{i}-\frac{H_{k} y^{k} y^{k T} H_{k} y^{i}}{y^{k^{T}} H_{k} y^{k}} \tag{4.5}
\end{equation*}
$$

By the induction assumption $\sigma^{i}=H_{k} y^{i}(i=0, \ldots, k-1)$, and therefore

$$
y^{k^{T}} H_{k} y^{i}=y^{k^{T}} \sigma^{i}=\sigma^{k^{T}} A \sigma^{i}=0 \quad(i=0, \ldots, k-1)
$$

Substituting in (4.5) we get the required

$$
H_{k+1} y^{i}=H_{k} y^{i}=\sigma^{i} \quad(i=0, \ldots, k-1)
$$

## DFP updates: discussion

- We have shown that the DFP updates preserve the required properties: positive definiteness, the secant condition, and the hereditary property.
- We have also shown that the DFP directions are mutually conjugate for quadratic functions.
- The DFP method can be applied to non-quadratic functions, but then the convergence of the DFP method is an open problem, even if the function is convex.
- In practice DFP performs quite well, but the method of choice today is the socalled BFGS update.


### 4.7.2 The BFGS update

The Broyden-Fletcher-Goldfarb-Shanno (BFGS) update is defined via

$$
D_{k}=\frac{\tau_{k} \sigma^{k} \sigma^{k} T}{}-\sigma^{k} y^{k^{T}} H_{k}-H_{k} y^{k} \sigma^{k^{T}},
$$

where

$$
\tau_{k}=1+\frac{y^{k T} H_{k} y^{k}}{\sigma^{k} T y^{k}}
$$

i If $y_{k}^{T} \sigma_{k}>0$, then $H_{k+1}=H_{k}+D_{k}$ is positive definite.
ii $H_{k+1}$ satisfies the secant and hereditary conditions.

Exercise 4.25 Consider the BFGS update:

$$
D_{k}=\frac{\tau_{k} \sigma^{k} \sigma^{k^{T}}-\sigma^{k} y^{k^{T}} H_{k}-H_{k} y^{k} \sigma^{k^{T}}}{\sigma^{k^{T}} y^{k}}
$$

where $y^{k}:=\nabla q\left(x^{k+1}\right)-\nabla q\left(x^{k}\right), \sigma^{k}:=x^{k+1}-x^{k}$, and

$$
\tau_{k}=1+\frac{y^{k^{T}} H_{k} y^{k}}{\sigma^{k^{T}} y^{k}} .
$$

(a) Show that if $y_{k}^{T} \sigma_{k}>0$, and $H_{k}$ is positive definite, then $H_{k+1}=H_{k}+D_{k}$ is positive definite.
(b) Show that the BFGS update satisfies the secant condition: $\sigma^{k}=H_{k+1} y^{k}$.

How do we guarantee $\sigma^{k T} y^{k}>0$ ? Note that $\sigma^{k}=\lambda_{k} s^{k}$ and $y^{k}=\nabla f\left(x^{k+1}\right)-\nabla f\left(x^{k}\right)$. Thus we need to maintain

$$
\nabla f\left(x^{k+1}\right)^{T} s^{k}>\nabla f\left(x^{k}\right)^{T} s^{k}
$$

This can be guaranteed by using a special line-search.
The convergence of the BFGS method for convex functions was proved in 1976 by Powell. In practice, BFGS outperforms DFP and is currently the Quasi-Newton method of choice.

Exercise 4.26 Consider the unconstrained optimization problem:

$$
\min 5 x_{1}^{2}+2 x_{1} x_{2}+x_{2}^{2}+7
$$

See Figure 4.7 for a contour plot.

1. Perform two iterations using the DFP Quasi-Newton method. Use exact line search and the starting point $[1,2]^{T}$. Plot the iterates.
2. Perform two iterations using the BFGS Quasi-Newton method. Use exact line search and the starting point $[1,2]^{T}$. Plot the iterates.


Figure 4.7: Contours of the objective function. Note that the minimum is at $[0,0]^{T}$.
$\triangleleft$

### 4.8 Stopping criteria

The stopping criteria is a relatively simple but essential part of the algorithms. If the algorithm generates both primal and dual solutions then the algorithm stops to iterate as the (relative) duality gap is less than a predefined threshold value $\epsilon>0$. The duality gap is defined as primal obj. value - dual obj. value
while the relative duality gap is usually defined as

$$
\frac{\text { primal obj. value }- \text { dual obj. value }}{1+\text { primal obj. value } \mid} .
$$

In unconstrained optimization it happens often that one uses a primal algorithm and then there is no such absolute measure to show how close we are to the optimum. Usually the algorithm is then stopped as there is no sufficient improvement in the objective, or if the iterates are too close to each other or if the length of the gradient or the length of the Newton step in an appropriate norm is small. All these criteria can be scaled (relative to) some characteristic number describing the dimensions of the problem. We give just two examples here. The relative improvement of the objective is not sufficient and the algorithm is stopped if at two subsequent iterate $x^{k}, x^{k+1}$

$$
\frac{\left|f\left(x^{k}\right)-f\left(x^{k+1}\right)\right|}{1+\left|f\left(x^{k}\right)\right|} \leq \epsilon
$$

In Newton's method we conclude that we are close to the minimum of the function if the length of the full Newton step in the norm induced by the Hessian is small, i.e.

$$
\begin{aligned}
\left\|\left(\nabla^{2} f\left(x^{k}\right)\right)^{-1} \nabla f\left(x^{k}\right)\right\|_{\nabla^{2} f\left(x^{k}\right)} & =\left(\nabla f\left(x^{k}\right)\right)^{T}\left(\nabla^{2} f\left(x^{k}\right)\right)^{-1} \nabla^{2} f\left(x^{k}\right)\left(\nabla^{2} f\left(x^{k}\right)\right)^{-1} \nabla f\left(x^{k}\right) \\
& =\left(\nabla f\left(x^{k}\right)\right)^{T}\left(\nabla^{2} f\left(x^{k}\right)\right)^{-1} \nabla f\left(x^{k}\right) \\
& \leq \epsilon .
\end{aligned}
$$

This criteria can also be interpreted as the length of the gradient measured in the norm induced by the inverse Hessian. This last measure is used in interior point methods to control the Newton process efficiently.

## Chapter 5

## Algorithms for constrained optimization

### 5.1 The reduced gradient method

The reduced gradient method can be viewed as the logical extension of the gradient method to constrained optimization problems. We start with linearly constrained optimization problems.
To this end, consider the following linearly constrained convex problem

$$
\begin{array}{rll}
(L C) & \min & f(x) \\
\text { s.t. } & A x=b,  \tag{5.1}\\
& x \geq 0 .
\end{array}
$$

## Assumptions:

- $f$ is continuously differentiable;
- Every subset of $m$ columns of the $m \times n$ matrix $A$ is linearly independent;
- each extreme point of the feasible set has at least $m$ positive components (nondegeneracy assumption).

Exercise 5.1 Prove that under the non-degeneracy assumption, every $x \in \mathcal{F}$ has at least $m$ positive components.

If $x \in \mathcal{F}$, we call a set of $m$ columns $B$ of $A$ a basis if $x_{i}>0$ when column $i$ is a column of $B$. We partition $x$ into basic $x_{B}$ and non-basic variables $x_{N}$ such that the basic variables $x_{B}>0$ correspond to the columns of $B$. Note that $x_{N}$ does not have to be zero.

For simplicity of notation we assume that we can partition the matrix $A$ as $A=[B, N]$. We partition $x$ accordingly: $x^{T}=\left[x_{B}, x_{N}\right]^{T}$. Thus we can rewrite $A x=b$ as

$$
B x_{B}+N x_{N}=b,
$$

such that

$$
x_{B}=B^{-1} b-B^{-1} N x_{N} .
$$

(Recall that $B^{-1}$ exists by assumption.)
Given $x \in \mathcal{F}$, we will choose $B$ as the columns corresponding to the $m$ largest components of $x$.

The basic variables $x_{B}$ can now be eliminated from problem (5.1) to obtain the reduced problem

$$
\begin{array}{cl}
\min & f_{N}\left(x_{N}\right) \\
\text { s.t. } & B^{-1} b-B^{-1} N x_{N} \geq 0, \\
& x_{N} \geq 0
\end{array}
$$

where $f_{N}\left(x_{N}\right)=f(x)=f\left(B^{-1} b-B^{-1} N x_{N}, x_{N}\right)$.
Note that any feasible direction $s$ for problem (LC) in (5.1) must satisfy $A s=0$. If we write $s^{T}=\left[s_{B}^{T}, s_{N}^{T}\right]$ for a given basis $B$, the condition $A s=0$ can be rewritten as

$$
B s_{B}+N s_{N}=0
$$

We can solve this equation to obtain:

$$
\begin{equation*}
s_{B}=-(B)^{-1} N s_{N} . \tag{5.2}
\end{equation*}
$$

## The choice of search direction

Recall that $s$ is a descent direction of $f$ at $x \in \mathcal{F}$ if and only if $\nabla f(x)^{T} s<0$, which translates to

$$
\nabla_{B} f(x)^{T} s_{B}+\nabla_{N} f(x)^{T} s_{N}<0
$$

Here $\nabla_{B} f(x)$ is the gradient with respect to the basic variables, etc.
Substitute $s_{B}$ from (5.4) to get:

$$
\nabla f(x)^{T} s=\left(-\nabla_{B} f(x)^{T}(B)^{-1} N+\nabla_{N} f(x)^{T}\right) s_{N}
$$

Definition 5.1 We call

$$
\begin{equation*}
r:=\left(-\nabla_{B} f(x)^{T}(B)^{-1} N+\nabla_{N} f(x)^{T}\right)^{T} \tag{5.3}
\end{equation*}
$$

the reduced gradient of $f$ at $x$ for the given basis $B$.

Note that

$$
\nabla f(x)^{T} s=r^{T} s_{N}
$$

In other words, the reduced gradient $r$ plays the same role in the reduced problem as the gradient $\nabla f$ did in the original problem (LC). In fact, the reduced gradient is exactly the gradient of the function $f_{N}$ with respect to $x_{N}$ in the reduced problem.

Exercise 5.2 Prove that $r=\nabla_{N} f_{N}\left(x_{N}\right)$, where $f_{N}\left(x_{N}\right)=f(x)=f\left(B^{-1} b-B^{-1} N x_{N}, x_{N}\right)$. $\triangleleft$

Recall that the gradient method uses the search direction $s=-\nabla f(x)$. Analogously, the basic idea for the reduced gradient method is to use the negative reduced gradient $s_{N}=-r$ as search direction for the variables $x_{N}$, and then calculating the search direction for the variables $x_{B}$ from

$$
\begin{equation*}
s_{B}=-(B)^{-1} N s_{N} \tag{5.4}
\end{equation*}
$$

At iteration $k$ of the algorithm we then perform a line search: find $0 \leq \lambda \leq \lambda_{\max }$ such that

$$
x^{k+1}=x^{k}+\lambda s^{k} \geq 0
$$

where $\lambda_{\max }$ is an upper bound on the maximal feasible step length and is given by

$$
\lambda_{\max }= \begin{cases}\min _{j:\left(s^{k}\right)_{j}<0} \frac{-\left(x^{k}\right)_{j}}{\left(s^{k}\right)_{j}} & \text { if } s^{k} \nsupseteq 0  \tag{5.5}\\ \infty & \text { if } s^{k} \geq 0\end{cases}
$$

This choice for $\lambda_{\max }$ guarantees that $x^{k+1} \geq 0$ and $A x^{k+1}=b$.

## Necessary modifications to the search direction

If we choose $s_{N}=-r$, then it may happen that

$$
\left(s_{N}\right)_{i}<0 \text { and }\left(x_{N}\right)_{i}=0
$$

for some index $i$.
In this case $\lambda_{\max }=0$ and we cannot make a step. One possible solution is as follows: for the nonbasic components set

$$
\left(s_{N}\right)_{i}= \begin{cases}-\left(x_{N}\right)_{i} r_{i} & \text { if } r_{i}>0  \tag{5.6}\\ -r_{i} & \text { if } r_{i} \leq 0\end{cases}
$$

Note that this prevents zero and 'very small' steps.

## Convergence results

Since the reduced gradient method may be viewed as an extension of the gradient method, it may come as no surprise that analogous converge results hold for the reduced gradient method as for the gradient method. In this section we state some convergence results and emphasize the analogy with the results we have already derived for the gradient method (see Theorem 4.5).
Assume that the reduced gradient method generates iterates $\left\{x^{k}\right\}, k=0,1,2, \ldots$

Theorem 5.2 The search direction $s^{k}$ at $x^{k}$ is always a feasible descent direction unless $s^{k}=0$. If $s^{k}=0$, then $x^{k}$ is a KKT point of problem (LC).

Compare this to the gradient method where, by definition, $s^{k}=0$ if and only if $x^{k}$ is a stationary point $\left(\nabla f\left(x^{k}\right)=0\right)$.

Exercise 5.3 Prove Theorem 5.2.

Theorem 5.3 Any accumulation point of $\left\{x^{k}\right\}$ is a KKT point.

Compare this to the gradient method where any accumulation point of $\left\{x^{k}\right\}$ is a stationary point under some assumptions (see Theorem 4.5).
The proof of Theorem 5.3 is beyond the scope of this course. A detailed proof is given in [2], Theorem 10.6.3.

## The reduced gradient algorithm: a summary

To summarize, we give a statement of the complete algorithm.

## 1. Initialization

Choose a starting point $x^{0} \geq 0$ such that $A x=b$. Let $k=0$.

## 2. Main step

[1.1] Form $B$ from those columns of $A$ that correspond to the $m$ largest components of $x^{k}$. Define $N$ as the remaining columns of $A$. Define $x_{B}$ as the elements of $x^{k}$ that correspond to $B$, and define $x_{N}$ similarly.
[1.2] Compute the reduced gradient $r$ from (5.3).
[1.3] Compute $s_{N}$ from (5.6) and $s_{B}$ from (5.4). Form $s^{k}$ from $s_{B}$ and $s_{N}$.
[1.4] If $s^{k}=0, \operatorname{STOP}\left(x^{k}\right.$ is a KKT point $)$.

## 3. Line search

[2.1] Compute $\lambda_{\max }$ from (5.5).
[2.2] Perform the line search

$$
\lambda_{k}:=\arg \min _{0 \leq \lambda \leq \lambda \max } f\left(x^{k}+\lambda s^{k}\right)
$$

[2.3] Set $x^{k+1}=x^{k}+\lambda_{k} s^{k}$ and replace $k$ by $k+1$.
[2.4] Repeat the main step.

## Remarks:

- During the algorithm the solution $x^{k}$ is not necessarily a basic solution, hence positive coordinates in $x_{N}^{k}$ may appear. These variables are usually referred to as superbasic variables.
- Recall that we have made a non-degeneracy assumption that is difficult to check in practice. If degeneracy occurs in practice, similar techniques as in the linear optimization case are applied to resolve degeneracy and prevent cycling.
- The convex simplex method is obtained as the specialization of the above reduced gradient scheme if the definition of the search direction $s_{N}$ is modified. We allow only one coordinate $j$ of $s_{N}$ to be nonzero and defined as $s_{j}=-\frac{\partial f_{N}\left(x_{N}^{k}\right)}{\partial x_{j}}>0$. The rest of the $s_{N}$ coordinates is defined to be zero and $s_{B}=-B^{-1} N s_{N}=-B^{-1} a_{j} s_{j}$, where $a_{j}$ is the $j$-th column of the matrix $A$.
- The simplex method of LO is obtained as a further specialization of the convex simplex method. One assumes that the objective function is linear and the initial solution is a basic solution.

Example 5.4 [Reduced gradient method 1] Consider the following problem:

$$
\begin{aligned}
\min & x^{2} \\
\text { s.t. } & x \geq 2 \\
& x \geq 0
\end{aligned}
$$

We solve this problem by using the Reduced Gradient Method starting from the starting point $x^{0}=5$ with objective value 25 . We start with converting the constraint in an equality-constraint:

$$
\begin{aligned}
\min x^{2} & \\
\text { s.t. } x-y & =2 \\
x, y & \geq 0
\end{aligned}
$$

The value of the slack variable $y^{0}$ is 3 . We therefore choose variable $x$ as the basic variable. This results in $B=1$ and $N=-1$. We eliminate the basic variable:

$$
f_{N}\left(x_{N}\right)=f\left(B^{-1} b-B^{-1} N x_{N}, x_{N}\right)=f(2+y, y)
$$

This gives us the following problem:

$$
\begin{aligned}
\min \quad(2+y)^{2} & \\
\text { s.t. } 2+y & \geq 0 \\
y & \geq 0 .
\end{aligned}
$$

## Iteration 1

The search directions are:

$$
\begin{aligned}
& s_{N}^{0}=s_{y}^{0}=-\frac{\delta f_{N}\left(y^{0}\right)}{\delta y}=-\left(2\left(2+y^{0}\right)\right)=-10 \\
& s_{B}^{0}=s_{x}^{0}=-B^{-1} N s_{N}^{0}=(-1) \cdot(-1) \cdot(-10)=-10
\end{aligned}
$$

The new values of the variables are, depending on the step-length $\lambda$ :

$$
\begin{aligned}
x^{1}=x^{0}+\lambda s_{x}^{0} & =5-10 \lambda \\
y^{1}=y^{0}+\lambda s_{y}^{0} & =3-10 \lambda
\end{aligned}
$$

which stay non-negative if $\lambda \leq \bar{\lambda}=\frac{3}{10}$.
We now have to solve the one-dimensional problem:

$$
\min (5-10 \lambda)^{2}
$$

The minimum is attained when

$$
-20(5-10 \lambda)=0
$$

i.e. when $\lambda=\frac{1}{2}$. Since the $\lambda=\frac{1}{2}$ is larger than $\bar{\lambda}=\frac{3}{10}$ that preserves the non-negativity of the variables, we have to take $\lambda=\frac{3}{10}$ as the step-length. This results in $x^{1}=2$ and $y^{1}=0$ with 4 as the objective value.

## Iteration 2

Since $x>y$, we use the variable $x$ as basic variable again. First we compute the search direction of $y$. Because $y^{0}=0$ the search direction has to be non-negative else it will get the value 0 :

$$
s_{N}^{3}=s_{y}^{3}=-2\left(2+y^{3}\right)=-4
$$

This means:

$$
\begin{aligned}
s_{N}^{3} & =s_{y}^{3}=0 \\
s_{B}^{3} & =s_{x}^{3}=0
\end{aligned}
$$

Thus the optimum point is $x^{o p t}=2$.

Example 5.5 [Reduced gradient method 2] Consider the following problem:

$$
\min \begin{aligned}
x_{1}^{2}+x_{2}^{2}+x_{3}^{2}+x_{4}^{2}-2 x_{1}-3 x_{4} & \\
\text { s.t. } 2 x_{1}+x_{2}+x_{3}+4 x_{4} & =7 \\
x_{1}+x_{2}+2 x_{3}+x_{4} & =6 \\
x_{1}, x_{2}, x_{3}, x_{4} & \geq 0
\end{aligned}
$$

We perform one iteration of the Reduced Gradient Method starting from the point $x^{0}=\left(x_{1}^{0}, x_{2}^{0}, x_{3}^{0}, x_{4}^{0}\right)^{T}=$ $(2,2,1,0)^{T}$ with an objective value 5 . At this point $x^{0}$ we consider $x_{1}$ and $x_{2}$ as basic variables. This results in

$$
B=\left(\begin{array}{cc}
2 & 1 \\
1 & 1
\end{array}\right) \text { and } N=\left(\begin{array}{cc}
1 & 4 \\
2 & 1
\end{array}\right)
$$

We eliminate the basic variables to obtain the reduced problem:

$$
\begin{aligned}
f_{N}\left(x_{N}\right) & =f\left(B^{-1} b-B^{-1} N x_{N}, x_{N}\right) \\
& =f\left(\left(\begin{array}{cc}
1 & -1 \\
-1 & 2
\end{array}\right)\binom{7}{6}-\left(\begin{array}{cc}
1 & -1 \\
-1 & 2
\end{array}\right)\left(\begin{array}{cc}
1 & 4 \\
2 & 1
\end{array}\right)\binom{x_{3}}{x_{4}}, x_{3}, x_{4}\right) \\
& =f\left(1+x_{3}-3 x_{4}, 5-3 x_{3}+2 x_{4}, x_{3}, x_{4}\right) .
\end{aligned}
$$

This results in the following problem:

$$
\begin{aligned}
\min \left(1+x_{3}-3 x_{4}\right)^{2}+\left(5-3 x_{3}+2 x_{4}\right)^{2} & +x_{3}^{2}+x_{4}^{2}-2\left(1+x_{3}-3 x_{4}\right)-3 x_{4} \\
1+x_{3}-3 x_{4} & \geq 0 \\
5-3 x_{3}+2 x_{4} & \geq 0 \\
x_{3}, x_{4} & \geq 0
\end{aligned}
$$

## Iteration 1

The search directions are:

$$
\begin{aligned}
s_{N}^{0}=\binom{s_{3}^{0}}{s_{4}^{0}} & =\binom{-\frac{\delta f_{N}\left(x_{3}^{0}\right)}{\delta x_{3}}}{-\frac{\delta f_{N}\left(x_{4}^{0}\right)}{\delta x_{4}}} \\
& =\binom{-\left(2\left(1+x_{3}^{0}-3 x_{4}^{0}\right)-6\left(5-3 x_{3}^{0}+2 x_{4}^{0}\right)+2 x_{3}^{0}-2\right)}{-\left(-6\left(1+x_{3}^{0}-3 x_{4}^{0}\right)+4\left(5-3 x_{3}^{0}+2 x_{4}^{0}\right)+2 x_{4}^{0}+3\right)} \\
& =\binom{8}{1} .
\end{aligned}
$$

Because $x_{4}^{0}=0$ the search direction $s_{4}^{0}$ has to be non-negative. We see that this is true.

$$
s_{B}^{0}=\binom{s_{1}^{0}}{s_{2}^{0}}=-B^{-1} N s_{N}^{0}=-\left(\begin{array}{cc}
1 & -1 \\
-1 & 2
\end{array}\right)\left(\begin{array}{cc}
1 & 4 \\
2 & 1
\end{array}\right)\binom{8}{1}=\binom{5}{-22}
$$

We now have to make a line search to obtain the new variables. These new variables as a function of the step length $\lambda$ are:

$$
\begin{aligned}
x_{1}^{1}=x_{1}^{0}+\lambda s_{1}^{0} & =2+5 \lambda \\
x_{2}^{1}=x_{2}^{0}+\lambda s_{2}^{0} & =2-22 \lambda \\
x_{3}^{1}=x_{3}^{0}+\lambda s_{3}^{0} & =1+8 \lambda \\
x_{4}^{1}=x_{4}^{0}+\lambda s_{4}^{0} & =\lambda
\end{aligned}
$$

which stay non-negative if $\lambda \leq \bar{\lambda}=\frac{2}{22} \approx 0.09$.
We proceed by solving

$$
\min (2+5 \lambda)^{2}+(2-22 \lambda)^{2}+(1+8 \lambda)^{2}+\lambda_{2}-2(2+5 \lambda)-3 \lambda
$$

This means

$$
\begin{aligned}
10(2+5 \lambda)-44(2-22 \lambda)+16(1+8 \lambda)+2 \lambda-13 & =0 \\
\lambda & =\frac{65}{1148} \approx 0.06 \quad\left(\lambda<\bar{\lambda}=\frac{2}{22}\right)
\end{aligned}
$$

The minimizer $\lambda=\frac{65}{1148}$ is smaller than $\bar{\lambda}=\frac{2}{22}$, so non-negativity of the variables is preserved. Thus the new iterate is $x^{1}=\left(x_{1}^{1}, x_{2}^{1}, x_{3}^{1}, x_{4}^{1}\right)^{T}=(2.28,0.75,1.45,0.06)^{T}$ with an objective value of 3.13 .

Exercise 5.4 Perform two iterations of the reduced gradient method for the following linearly constrained convex optimization problem:

$$
\begin{aligned}
\min x_{1}^{2}+x_{2}^{4}+\left(x_{3}-x_{4}\right)^{2} & \\
\text { s.t. } \quad x_{1}+2 x_{2}+3 x_{3}+4 x_{4} & =10 \\
x & \geq 0 .
\end{aligned}
$$

Let the initial point be given as $x^{0}=(1,1,1,1)$ and use $x_{1}$ as the initial basic variable. $\triangleleft$

### 5.2 Generalized reduced gradient (GRG) method

The reduced gradient method can be generalized to nonlinearly constrained optimization problems. Similarly to the linearly constrained case we consider the problem with equality constraints and nonnegative variables as follows.

$$
\begin{align*}
(N C) \min & f(x) \\
\text { s.t. } & h_{j}(x)=0, \quad j=1, \cdots, m  \tag{5.7}\\
& x \geq 0
\end{align*}
$$

where the functions $f, h_{1}, \cdots, h_{m}$ supposed to be continuously differentiable. ${ }^{1}$
The basis idea is to replace the nonlinear equations by their linear Taylor approximation at the current value of $x$, and then apply the reduced gradient algorithm to the resulting problem.
We assume that the gradients of the constraint functions $h_{j}$ are linearly independent at every point $x \geq 0$, and that each feasible $x$ has at least $m$ positive components. These assumptions ensure that we can always apply the reduced gradient algorithm to the linearized problem. The extra difficulty here is that - since the feasible region $\mathcal{F}$ is not convex - this procedure may produce iterates that lie outside $\mathcal{F}$, and then some extra effort is needed to restore feasibility.

Let a feasible solution $x^{k} \geq 0$ with $h_{j}\left(x^{k}\right)=0$ for all $j$ be given. By assumption the Jacobian matrix of the constraints $H(x)=\left(h_{1}(x), \cdots, h_{m}(x)\right)^{T}$ at each $x \geq 0$ has full

[^3]rank and, for simplicity at the point $x^{k}$ will be denoted by
$$
A=J H\left(x^{k}\right) .
$$

Let us assume that a basis $B$, where $x_{B}^{k}>0$ is given. Then a similar construction as in the linear case apply. We generate a reduced gradient search direction by virtually keeping the linearized constraints valid. This direction by construction will be in the null space of $A$. More specifically for the linearized constraints we have

$$
H\left(x^{k}\right)+J H\left(x^{k}\right)\left(x-x^{k}\right)=0+A\left(x-x^{k}\right)=0
$$

From this one has

$$
B x_{B}+N x_{N}=A x^{k}
$$

and by introducing the notation $b=A x^{k}$ we have

$$
x_{B}=B^{-1} b-B^{-1} N x_{N}
$$

hence the basic variables $x_{B}$ can be eliminated from the linearization of the problem (5.7) to result

$$
\begin{array}{cl}
\min & f_{N}\left(x_{N}\right) \\
\text { s.t. } & B^{-1} b-B^{-1} N x_{N} \geq 0, \\
& x_{N} \geq 0 .
\end{array}
$$

where $f_{N}\left(x_{N}\right)=f(x)=f\left(B^{-1} b-B^{-1} N x_{N}, x_{N}\right)$. Using the notation

$$
\nabla f(x)^{T}=\left(\left(\nabla_{B} f(x)\right)^{T},\left(\nabla_{N} f(x)\right)^{T}\right)
$$

the gradient of $f_{N}$, namely the reduced gradient can be expressed as

$$
\nabla_{N} f(x)^{T}=-\left(\nabla_{B} f(x)\right)^{T} B^{-1} N+\left(\nabla_{N} f(x)\right)^{T}
$$

From this point on the generation of the search direction $s$ proceeds in exactly the same way as in the linearly constrained case. Due to the nonlinearity of the constraints $H\left(x^{k+1}\right)=H\left(x^{k}+\lambda s\right)=0$ will not hold in general. Hence something more has to be done to restore feasibility.
Special care has to be taken to control the step size. A larger step size might allow larger improvement of the objective but, on the other hand results in larger infeasibility of the constraints. A good compromise must be made.
In old versions of the GRG method Newton's method is applied to the nonlinear equality system $H(x)=0$ from the initial point $x^{k+1}$ to produce a next feasible iterate. In more recent implementations the reduced gradient direction is combined by a direction from the orthogonal subspace (the range space of $A^{T}$ ) and then a modified (nonlinear, discrete) line search is performed. These schemes are quite complicated and not discussed here in more detail.

Example 5.6 [Generalized reduced gradient method 1] We consider the following problem:

$$
\min \quad \begin{aligned}
x_{1}^{2}+x_{2}^{2}+12 x_{1}-4 x_{2} & \\
\text { s.t. } x_{1}^{2}-2 x_{2} & =0 \\
x_{1}, x_{2} & \geq 0
\end{aligned}
$$

We perform two steps of the Generalized Reduced Gradient Method starting from the point $x^{0}=$ $\left(x_{1}^{0}, x_{2}^{0}\right)^{T}=(4,8)^{T}$ with an objective value of 96 . We will plot the progress of the algorithm in Figure 5.1. At the point $x^{0}$ we consider $x_{2}$ as the basic variable. First we have to linearize the nonlinearly constraint:

$$
A=(N, B)=J H\left(x^{0}\right)=\left(\begin{array}{ll}
2 x_{1}^{0} & -2
\end{array}\right)=\left(\begin{array}{ll}
-8 & -2
\end{array}\right) . \quad b=A x^{0}=\left(\begin{array}{ll}
-8 & -2
\end{array}\right)\binom{4}{8}=16
$$

Now we eliminate the basic variable:

$$
f_{N}\left(x_{N}\right)=f\left(B^{-1} b-B^{-1} N x_{N}, x_{N}\right)=f\left(x_{1},-\frac{1}{2} \cdot 16+\frac{1}{2} x_{1}\right)
$$

This leads us to the following problem:

$$
\min \quad \begin{aligned}
x_{1}^{2}+\left(4 x_{1}-8\right)^{2}+12 x_{1}-4\left(4 x_{1}-8\right) & \\
\text { s.t. } 4 x_{1}-8 & \geq 0 \\
x_{1} & \geq 0
\end{aligned}
$$

## Iteration 1

The search direction is:

$$
\begin{aligned}
& \left.s_{N}^{0}=s_{1}^{0}=-\frac{\delta f_{N}\left(x_{1}^{0}\right)}{\delta x_{1}}=-\left(2 x_{1}^{0}+8\left(4 x_{1}-8\right)+12-16\right)\right)=-68 \\
& s_{B}^{0}=s_{1}^{0}=-B^{-1} N s_{N}^{0}=\frac{1}{2} \cdot 8 \cdot-68=-272
\end{aligned}
$$

The new variables as a function of the step length $\lambda$ are:

$$
\begin{aligned}
& x_{1}^{1}=x_{1}^{0}+\lambda s_{1}^{0}=4-68 \lambda \\
& x_{2}^{1}=x_{2}^{0}+\lambda s_{2}^{0}=8-272 \lambda
\end{aligned}
$$

which stay non-negative if $\lambda \leq \bar{\lambda}=\frac{1}{34}$.
We do this by solving

$$
\min (4-68 \lambda)^{2}+(8-272)^{2}+12(4-68 \lambda)-4(8-272 \lambda)
$$

This means

$$
\begin{aligned}
-136(4-68 \lambda)-544(8-272 \lambda)-816+1088 & =0 \\
\lambda & =\frac{1}{34} \quad(\lambda=\bar{\lambda})
\end{aligned}
$$

This results in $x^{1}=\left(x_{1}^{1}, x_{2}^{1}\right)^{T}=(2,0)^{T}$. But due to the nonlinearity of the constraint, the constraint will not hold with these values. To find a solution for which the constraint will hold, we consider the $x_{N}$ as a fixed variable. The $x_{B}$ will change in a value for which the constraint holds, this means $x_{B}=2$. The objective value is 24 .

## Iteration 2

Because $x_{2}^{1}$ stayed positive we now use $x_{1}^{1}$ as basic variable again. But first we have to linearize the nonlinearly constraint with the values of iteration 1 :

$$
A=J H\left(x^{1}\right)=\left(\begin{array}{ll}
2 x_{1}^{1} & -2
\end{array}\right)=\left(\begin{array}{ll}
4 & -2
\end{array}\right) . \quad b=A x^{1}=\left(\begin{array}{ll}
4 & -2
\end{array}\right)\binom{2}{2}=4
$$

We eliminate the basic variable:

$$
f_{N}\left(x_{N}\right)=f\left(B^{-1} b-B^{-1} N x_{N}, x_{N}\right)=f\left(x_{1},-\frac{1}{2} \cdot 4+\frac{1}{2} \cdot 4 \cdot x_{1}\right)=f\left(x_{1}, 2 x_{1}-2\right)
$$

This gives us the following problem:

$$
\min \quad \begin{aligned}
x_{1}^{2}+\left(2 x_{1}-2\right)^{2}+12 x_{1}-4\left(2 x_{1}-2\right) & \\
\text { s.t. } 2 x_{1}-2 & \geq 0 \\
x_{1} & \geq 0 .
\end{aligned}
$$

The search direction is:

$$
\begin{aligned}
& s_{N}^{1}=s_{1}^{1}=-\frac{\delta f_{N}\left(x_{1}^{1}\right.}{\delta x_{1}}=-\left(2 x_{1}^{1}+4\left(2 x_{1}^{1}-2\right)+12-8\right)=-16 \\
& s_{B}^{1}=s_{2}^{1}=-B^{-1} N s_{N}^{1}=-32
\end{aligned}
$$

The new variables as a function of the step length $\lambda$ are:

$$
\begin{aligned}
& x_{1}^{2}=x_{1}^{1}+\lambda s_{1}^{1}=2-16 \lambda \\
& x_{2}^{2}=x_{2}^{1}+\lambda s_{2}^{1}=2-32 \lambda
\end{aligned}
$$

which stay non-negative if $\lambda \leq \bar{\lambda}=\frac{2}{32}$.
Now we have to solve

$$
\min \quad 2(2-16 \lambda)^{2}+(2-32 \lambda)^{2}+12(2-16 \lambda)-4(2-32 \lambda)
$$

This means:

$$
\begin{aligned}
-32(2-16 \lambda)-64(2-32 \lambda)-192+128 & =0 \\
\lambda & =\frac{1}{10} \quad\left(\lambda>\bar{\lambda}=\frac{1}{16}\right)
\end{aligned}
$$

As we can see has $\lambda=\frac{1}{10}$ a larger value than $\bar{\lambda}=\frac{1}{16}$. In order to get non-negative values for the variables we have to use the value $\frac{1}{16}$ as step length. This gives us $x^{2}=\left(x_{1}^{2}, x_{2}^{2}\right)^{T}=(1,0)^{T}$. To get variables for which the constraint holds, we take the $x_{N}$ as fixed variable. This leads to $x^{2}=\left(1, \frac{1}{2}\right)^{T}$ with an objective value of $11 \frac{1}{4}$.

Example 5.7 [Generalized reduced gradient method 2] We consider the following problem:

$$
\begin{array}{rrl}
\min & 2 x_{1}^{2}+3 x_{2}^{2} \\
\text { s.t. } & 3 x_{1}^{2}+2 x_{2}^{2} & =20 \\
& x_{1}, x_{2} \geq 0 .
\end{array}
$$



Figure 5.1: Illustration of Example 5.6.

We solve this problem by using three steps of the Generalized Reduced Gradient Method starting from the point $x^{0}=\left(x_{1}^{0}, x_{2}^{0}\right)^{T}=(2,2)^{T}$ with an objective value of 20 . At this point $x^{0}$ we consider $x_{1}$ as basic variable. First we have to linearize the nonlinearly constraint:

$$
A=(B, N)=J H\left(x^{0}\right)=\left(\begin{array}{ll}
6 x_{1}^{0} & 4 x_{2}^{0}
\end{array}\right)=\left(\begin{array}{ll}
12 & 8
\end{array}\right) . \quad b=A x^{0}=\left(\begin{array}{ll}
12 & 8
\end{array}\right)\binom{2}{2}
$$

Now we eliminate the basic variables:

$$
f_{N}\left(x_{N}\right)=f\left(B^{-1} b-B^{-1} N x_{N}, x_{N}\right)=f\left(\frac{40}{12}-\frac{8}{12} x_{2}, x_{2}\right)
$$

This leads us to the following problem:

$$
\begin{aligned}
\min \quad 2\left(\frac{40}{12}-\frac{8}{12} x_{2}\right)^{2}+3 x_{2}^{2} & \\
\text { s.t. } \frac{40}{12}-\frac{8}{12} x_{2} & \geq 0 \\
x_{2} & \geq 0
\end{aligned}
$$

Iteration 1 The search direction is:

$$
s_{N}^{0}=s_{2}^{0}=-\frac{\delta f_{N}\left(x_{2}^{0}\right)}{\delta x_{2}}=-\left(-\frac{32}{12}\left(\frac{40}{12}-\frac{8}{12} x_{2}^{0}\right)+6 x_{2}^{0}\right)=-\frac{20}{3}
$$

$$
s_{B}^{0}=s_{1}^{0}=-B^{-1} N s_{N}^{0}=-\frac{1}{12} \cdot 8 \cdot-\frac{20}{3}=\frac{40}{9}
$$

The new variables as a function of $\lambda$ are:

$$
\begin{aligned}
& x_{1}^{1}=x_{1}^{0}+\lambda s_{1}^{0}=2+\frac{40}{9} \lambda \\
& x_{2}^{1}=x_{2}^{0}+\lambda s_{2}^{0}=2-\frac{20}{3} \lambda
\end{aligned}
$$

which are non-negative as long as $\lambda \leq \bar{\lambda}=\frac{2}{\frac{20}{3}}=\frac{3}{10}$.
We do this by solving

$$
\min \quad 2\left(2+\frac{40}{9} \lambda\right)^{2}+3\left(2-\frac{20}{3} \lambda\right)^{2} .
$$

This means

$$
\begin{aligned}
\frac{160}{9}\left(2+\frac{40}{9} \lambda\right)-\frac{120}{3}\left(2-\frac{20}{3} \lambda\right) & =0 \\
\lambda & =\frac{9}{70} \quad\left(\lambda<\bar{\lambda}=\frac{3}{10}\right)
\end{aligned}
$$

This results in $x^{1}=\left(x_{1}^{1}, x_{2}^{1}\right)^{T}=\left(\frac{18}{7}, \frac{8}{7}\right)^{T}$. But due to the nonlinearity of the constraint, the constraint will not hold with these values. To find a solution for which the constraint will hold, we consider the $x_{N}$ as a fixed variable. The $x_{B}$ will change in a value for which the constraint holds, this means $x_{B}=2.41$. The objective value is 15.52 .

## Iteration 2

Because $x_{1}^{1}$ stayed positive we use $x_{1}$ as basic variable again. First with the values of iteration 1 we linearize the nonlinearly constraint again:

$$
A=J H\left(x^{1}\right)=\left(\begin{array}{ll}
6 x_{1}^{1} & 4 x_{2}^{1}
\end{array}\right)=\left(\begin{array}{ll}
14.45 & 4.57
\end{array}\right) . \quad b=A x^{1}=\left(\begin{array}{ll}
14.45 & 4.57
\end{array}\right)\binom{2.41}{1.14}=40
$$

We eliminate the basic variable:

$$
f_{N}\left(x_{N}\right)=f\left(B^{-1} b-B^{-1} N x_{N}, x_{N}\right)=f\left(2.77-0.32 x_{2}, x_{2}\right)
$$

This gives us the following problem:

$$
\begin{aligned}
\min \quad 2\left(2.77-0.32 x_{2}\right)^{2}+3 x_{2}^{2} & \\
\text { s.t. } 2.77-0.32 x_{2} & \geq 0 \\
x_{2} & \geq 0 .
\end{aligned}
$$

The search direction is:

$$
\begin{aligned}
& s_{N}^{1}=s_{2}^{1}=-\frac{\delta f_{N}\left(x_{2}^{1}\right.}{\delta x_{2}}=-\left(-4 \cdot 0.32\left(2.77-0.32 x_{2}^{1}\right)+6 x_{2}^{1}\right)=-3.78 \\
& s_{B}^{1}=s_{1}^{1}=-B^{-1} N s_{N}^{1}=1.2
\end{aligned}
$$

The new variables, depending on the step length $\lambda$, are:

$$
\begin{aligned}
& x_{1}^{2}=x_{1}^{1}+\lambda s_{1}^{1}=2.41+1.20 \lambda \\
& x_{2}^{2}=x_{2}^{1}+\lambda s_{2}^{1}=1.14-3.78 \lambda
\end{aligned}
$$

which stay non-negative if $\lambda \leq \bar{\lambda}=\frac{1.14}{3.78} \approx 0.30$.
Now we have to solve

$$
\min \quad 2(2.41+1.20 \lambda)^{2}+3(1.14-3.78 \lambda)^{2}
$$

This means:

$$
\begin{aligned}
4.80(2.41+1.20 \lambda)-22.68(1.14-3.78 \lambda) & =0 \\
\lambda & =0.156 \quad(\lambda<\bar{\lambda} \approx 0.3)
\end{aligned}
$$

This gives us $x^{2}=\left(x_{1}^{2}, x_{2}^{2}\right)^{T}=(2.6,0.55)^{T}$. To get variables for which the constraint holds, we take the $x_{N}$ as fixed variable. This leads to $x^{2}=(2.52,0.55)^{T}$ with an objective value of 13.81 .

## Iteration 3

Again we can use $x_{1}$ as basic variable. We start this iteration with linearization of the constraint:

$$
A=J H\left(x^{2}\right)=\left(\begin{array}{ll}
6 x_{1}^{2} & 4 x_{2}^{2}
\end{array}\right)=\left(\begin{array}{ll}
15.24 & 2.2
\end{array}\right) . \quad b=A x^{2}=\left(\begin{array}{ll}
15.24 & 2.2
\end{array}\right)\binom{2.54}{0.55}=39.9
$$

Eliminating the basic variable:

$$
f_{N}\left(x_{N}\right)=f\left(2.62-0.14 x_{2}, x_{2}\right)
$$

This gives us the following problem:

$$
\min \quad \begin{aligned}
2\left(2.62-0.14 x_{2}\right)^{2}+3 x_{2}^{2} & \\
\text { s.t. } 2.62-0.14 x_{2} & \geq 0 \\
x_{2} & \geq 0
\end{aligned}
$$

Search directions:

$$
\begin{aligned}
& s_{N}^{2}=s_{2}^{2}=-\frac{\delta f_{N}\left(x_{2}^{2}\right.}{\delta x_{2}}=-\left(-0.56\left(2.62-0.14 x_{2}^{2}\right)+6 x_{2}^{2}\right)=-1.88 \\
& s_{B}^{2}=s_{1}^{2}=-B^{-1} N s_{N}^{2}=0.27
\end{aligned}
$$

New variables as a function of $\lambda$ :

$$
\begin{aligned}
& x_{1}^{3}=x_{1}^{2}+\lambda s_{1}^{2}=2.52+0.27 \lambda \\
& x_{2}^{3}=x_{2}^{2}+\lambda s_{2}^{2}=0.55-1.88 \lambda
\end{aligned}
$$

which stay non-negative if $\lambda \leq \bar{\lambda}=\frac{0.55}{1.88} \approx 0.293$.
Now we solve

$$
\min \quad 2(2.54+0.27 \lambda)^{2}+3(0.55-1.88 \lambda)^{2}
$$

This means:

$$
\begin{aligned}
1.08(2.52+0.27 \lambda)-5.64(0.55-1.88 \lambda) & =0 \\
\lambda & =0.161
\end{aligned}
$$

This gives us the variables $x_{1}^{3}=2.58$ and $x_{2}^{3}=0.25$. Correcting the $x_{B}$ results in $x^{3}=(2.57,0.25)^{T}$ with objective value 13.39 .

Exercise 5.5 Perform one iteration of the generalized reduced gradient method to solve the following nonlinearly constrained convex optimization problem:

$$
\begin{aligned}
\min x_{1}^{2}+x_{2}^{4}+\left(x_{3}-x_{4}\right)^{2} & \\
\text { s.t. } \quad x_{1}^{2}+x_{2}^{2}+x_{3}^{2}+x_{4}^{2} & \leq 4 \\
x & \geq 0 .
\end{aligned}
$$

Let the initial point be given as $x^{0}=(1,1,1,1)$.
(You might need MAPLE or MATLAB to make the necessary calculations.)

## Chapter 6

## The Interior Point Approach to Nonlinear Optimization

### 6.1 Introduction

In this chapter we deal with the so-called logarithmic barrier approach to convex optimization. As before we consider the CO problem in the following format:

$$
(C O) \quad \min \{f(x): x \in \mathcal{F}\},
$$

where $\mathcal{F}$ denotes the feasible region, which is given by

$$
\mathcal{F}:=\left\{x \in \mathbb{R}^{n}: g_{j}(x) \leq 0, \quad 1 \leq j \leq m\right\} ;
$$

the constraint functions $g_{j}: \mathbb{R}^{n} \rightarrow \mathbb{R}(1 \leq j \leq m)$ and the objective function $f$ : $\mathbb{R}^{n} \rightarrow \mathbb{R}$ are convex functions with continuous third order derivatives in the interior of $\mathcal{F}$. Later on, when dealing with algorithms for solving $(C O)$ we will need to assume a smoothness condition on the functions $f$ and $g_{j}(1 \leq j \leq m)$. Without loss of generality we further assume that $f(x)$ is linear, i.e. $f(x)=-c^{T} x$ for some objective vector $c \in \mathbb{R}^{n}$. If this is not true, one may introduce an additional variable $x_{n+1}$, an additional constraint $f(x)-x_{n+1} \leq 0$, and minimize $x_{n+1}$. In this way the objective function becomes linear. Thus we may assume that ( $C O$ ) has the form

$$
(C P O)\left\{\begin{array}{l}
\min -c^{T} x \\
g_{j}(x) \leq 0, j=1, \cdots, m \\
x \in \mathbb{R}^{n}
\end{array}\right.
$$

The Lagrange-Wolfe dual of ( $C P O$ ) is given by

$$
(C D O)\left\{\begin{array}{l}
\max -c^{T} x+\sum_{j=1}^{m} y_{j} g_{j}(x) \\
\sum_{j=1}^{m} y_{j} \nabla g_{j}(x)=c \\
y_{j} \geq 0, \quad j=1, \cdots, m
\end{array}\right.
$$

Here we used that $\nabla\left(-c^{T} x\right)=-c$.
The interior of the primal feasible region $\mathcal{F}$ is denoted as

$$
\mathcal{F}^{0}:=\left\{x \in \mathbb{R}^{n}: g_{j}(x)<0, j=1, \cdots, m\right\}
$$

and we say that $(C P O)$ satisfies the interior point condition (IPC) if $\mathcal{F}^{0}$ is nonempty. In other words, $(C P O)$ satisfies the IPC if and only if there exists an $x$ that is strictly primal feasible (i.e. $\left.g_{j}(x)<0, \forall j=1, \cdots, m\right)$. Similarly, we say that $(C D O)$ satisfies the IPC if there exists a strictly dual feasible solution (i.e. a dual feasible pair $(x, y)$ with $y>0)$. We will show that if $(C P O)$ and $(C D O)$ both satisfy the IPC then these problems can be solved in polynomial time provided that the above mentioned smoothness condition is fulfilled. We will also present examples of large classes of well-structured CO problems which satisfy the smoothness condition.
Let us emphasize the trivial fact that if $(C P O)$ satisfies the IPC then $(C P O)$ is Slater regular. From now the IPC for $(C P O)$ (and hence Slater regularity) and $(C D O)$ will be assumed.

### 6.2 Duality and the central path

We will apply the Karush-Kuhn-Tucker theory, as developed in Section 2.2.4, to the problem (CPO). This theory basically consists of Theorem 2.30 and its corollaries (Corollary 2.31-2.35) and Definition 2.34. When applied to ( $C P O$ ) we get the following so-called KKT-theorem.

Theorem 6.1 The vector $x$ is optimal for (CPO) if and only if there exists a vector $y \geq 0\left(y \in \mathbb{R}^{m}\right)$ such that the pair $(x, y)$ is a saddle point of the Lagrange function

$$
L(x, y):=-c^{T} x+\sum_{j=1}^{m} y_{j} g_{j}(x)
$$

In this case $(x, y)$ is a Karush-Kuhn-Tucker (KKT) point of (CPO), which means

$$
\begin{align*}
g_{j}(x) & \leq 0, \forall j=1, \cdots, m  \tag{i}\\
\sum_{j=1}^{m} y_{j} \nabla g_{j}(x) & =c, y \geq 0  \tag{ii}\\
y_{j} g_{j}(x) & =0, \forall j=1, \cdots, m \tag{iii}
\end{align*}
$$

Note that ( $i$ ) ensures primal feasibility and (ii) dual feasibility. The third condition in the KKT-system is called the complementarity condition. The complementarity condition ensures that the duality gap at optimality is zero. This follows since the difference of the primal and the dual objective value, which the duality gap, is given by

$$
-\sum_{j=1}^{m} y_{j} g_{j}(x)
$$

We relax the complementarity condition by considering the system
(ii) $\sum_{j=1}^{m} y_{j} \nabla g_{j}(x)=c, y \geq 0$,

$$
\begin{equation*}
-y_{j} g_{j}(x)=\mu, \forall j=1, \cdots, m \tag{iii}
\end{equation*}
$$

for $\mu>0$. Clearly, if the relaxed KKT-system has a solution (for some $\mu>0$ ) then we have $x$ and $y$ such that $x$ is strictly primal feasible (i.e. $g_{j}(x)<0, \forall j=1, \cdots, m$ ) and the pair $(x, y)$ is strictly dual feasible (i.e. dual feasible with $y>0$ ), whereas the duality gap equals $m \mu$. In other words, if the relaxed KKT-system has a solution (for some $\mu>0$ ) then both $(C P O)$ and $(C D O)$ satisfy the IPC. If we impose some extra condition then, similarly to the case of linear optimization (LO), we also have the converse result: if the IPC holds then the relaxed KKT-system has a solution (for every $\mu>0)$. This will be presented in the following theorem, but first we have to introduce some definitions.

Definition 6.2 Let $\bar{x}, s \in \mathbb{R}^{n}$. The ray $\mathcal{R}:=\{x \mid x=\bar{x}+\lambda s, \lambda \in \mathbb{R}\} \subset \mathbb{R}^{n}$ is called bad if every constraint function $g_{j}, j=1, \cdots, m$ is constant along the ray $\mathcal{R}$.
Let $\bar{x}, s \in \mathbb{R}^{n}$ and $\alpha^{1}, \alpha^{2} \in \mathbb{R}$. The line segment $\left\{x \mid x=\bar{x}+\lambda s, \lambda \in\left[\alpha^{1}, \alpha^{2}\right]\right\} \subset \mathbb{R}^{n}$ is called bad if every constraint function $g_{j}, j=1, \cdots, m$ is constant along the ray.

Theorem 6.3 Let us assume that for (CPO) no bad ray exists. Then the following three statements are equivalent.
(i) $(C P O)$ and $(C D O)$ satisfy the interior point condition;
(ii) For each $\mu>0$ the relaxed KKT-system (6.2) has a solution;
(iii) For each $w>0\left(w \in \mathbb{R}^{m}\right)$ there exist $y$ and $x$ such that

The proof of this important theorem can be found in Appendix A.2. From now on we assume that the IPC holds.

Lemma 6.4 Let us assume that for ( $C P O$ ) no bad line segment exists. Then the solutions of the systems (6.2) and (6.3), if they exist, are unique.

Proof: See Appendix A.2.
The solution of (6.2) is denoted as $x(\mu)$ and $y(\mu)$. The set

$$
\{x(\mu): \mu>0\}
$$

is called the central path of $(C P O)$ and the set

$$
\{(x(\mu), y(\mu)): \mu>0\}
$$

the central path of $(C D O)$.

### 6.2.1 Logarithmic barrier functions

In the sequel we need a different characterization of the central paths of $(C P O)$ and $(C D O)$ that uses the so-called primal and dual logarithmic barrier functions $\phi_{B}(x, \mu)$ and $\phi_{B}^{d}(x, y, \mu)$. These functions are defined on the interior of the primal and dual feasible regions, respectively, according to

$$
\phi_{B}(x, \mu):=\frac{-c^{T} x}{\mu}-\sum_{j=1}^{m} \log \left(-g_{j}(x)\right)
$$

and

$$
\phi_{B}^{d}(x, y, \mu):=\frac{-c^{T} x+\sum_{j=1}^{m} y_{j} g_{j}(x)}{\mu}+\sum_{j=1}^{m} \log y_{j}+n(1-\log \mu) .
$$

Lemma 6.5 We have $\phi_{B}(\bar{x}, \mu) \geq \phi_{B}^{d}(x, y, \mu)$ for all primal feasible $\bar{x}$ and dual feasible $(x, y)$. Moreover, $\phi_{B}(x(\mu), \mu)=\phi_{B}^{d}(x(\mu), y(\mu), \mu)$ and, as a consequence, $x(\mu)$ is a minimizer of $\phi_{B}(x, \mu)$ and $(x(\mu), y(\mu))$ a maximizer of $\phi_{B}^{d}(x, y, \mu)$.

Proof: The proof uses that if $h: \mathcal{D} \rightarrow \mathbb{R}$ is differentiable and convex then we have

$$
h(\bar{x})-h(x) \geq \nabla h(x)^{T}(\bar{x}-x), \quad \forall x, \bar{x} \in \mathcal{D} .
$$

We refer for this property to Lemma 1.49 in Section 1.3. Since $-c^{T} x$ and $g_{j}(x), j=$ $1, \cdots, m$, are convex on $\mathcal{F}$ it follows that for any fixed $y \geq 0$ the Lagrange function

$$
L(x, y)=-c^{T} x+\sum_{j=1}^{m} y_{j} g_{j}(x)
$$

is convex as a function of $x$. Hence, if $\bar{x}$ is primal feasible and $(x, y)$ is dual feasible then

$$
-c^{T} \bar{x}+c^{T} x+\sum_{j=1}^{m} y_{j}\left(g_{j}(\bar{x})-g_{j}(x)\right) \geq\left(-c+\sum_{j=1}^{m} y_{j} \nabla g_{j}(x)\right)^{T}(\bar{x}-x)=0
$$

the last equality follows since $(x, y)$ is dual feasible. Thus we may write

$$
\begin{aligned}
\phi_{B}(\bar{x}, \mu) & -\phi_{B}^{d}(x, y, \mu) \\
& =\frac{1}{\mu}\left(-c^{T} \bar{x}+c^{T} x\right)-\frac{1}{\mu} \sum_{j=1}^{m} y_{j} g_{j}(x)-\sum_{j=1}^{m} \log \left(-y_{j} g_{j}(\bar{x})\right)-n(1-\log \mu)
\end{aligned}
$$

$$
\begin{aligned}
& \geq \frac{-1}{\mu} \sum_{j=1}^{m} y_{j} g_{j}(\bar{x})-\sum_{j=1}^{m} \log \left(-y_{j} g_{j}(\bar{x})\right)-n(1-\log \mu) \\
& =\sum_{j=1}^{m}\left(\frac{-y_{j} g_{j}(\bar{x})}{\mu}-1-\log \frac{-y_{j} g_{j}(\bar{x})}{\mu}\right) \\
& =\sum_{j=1}^{m} \psi\left(\frac{-y_{j} g_{j}(\bar{x})}{\mu}-1\right)
\end{aligned}
$$

where the function $\psi:(-1, \infty) \rightarrow \mathbb{R}_{+}$is defined by

$$
\begin{equation*}
\psi(t)=t-\log (1+t) \tag{6.4}
\end{equation*}
$$

The function $\psi$ is strictly convex, nonnegative and $\psi(0)=0$. See [40]. Hence the inequality in the lemma follows, and equality will hold only if

$$
-c^{T} \bar{x}+\sum_{j=1}^{m} y_{j} g_{j}(\bar{x})=-c^{T} x+\sum_{j=1}^{m} y_{j} g_{j}(x)
$$

and

$$
-y_{j} g_{j}(\bar{x})=\mu, \forall j=1, \cdots, m
$$

This implies that equality holds if $\bar{x}=x=x(\mu)$ and $y=y(\mu)$.

Thus the primal central path consists of minimizers of the primal logarithmic barrier function and the dual central path of maximizers of the dual logarithmic barrier function.

### 6.2.2 Monotonicity along the paths

Theorem 6.6 If $\mu$ decreases, then the primal objective function $-c^{T} x(\mu)$ monotonically decreases and the dual objective function $-c^{T} x(\mu)+\sum_{j=1}^{m} y_{j}(\mu) g_{j}(x(\mu))$ monotonically increases.
*Proof: Suppose $0<\bar{\mu}<\mu$. Then $x(\mu)$ minimizes $\phi_{B}(x, \mu)$, and $x(\bar{\mu})$ minimizes $\phi_{B}(x, \bar{\mu})$. Thus we have

$$
\begin{aligned}
\phi_{B}(x(\mu), \mu) & \leq \phi_{B}(x(\bar{\mu}), \mu) \\
\phi_{B}(x(\bar{\mu}), \bar{\mu}) & \leq \phi_{B}(x(\mu), \bar{\mu}) .
\end{aligned}
$$

These inequalities can be rewritten as

$$
\begin{aligned}
& -\frac{c^{T} x(\mu)}{\mu}-\sum_{j=1}^{m} \log \left(-g_{j}(x(\mu))\right) \leq-\frac{c^{T} x(\bar{\mu})}{\mu}-\sum_{j=1}^{m} \log \left(-g_{j}(x(\bar{\mu}))\right), \\
& -\frac{c^{T} x(\bar{\mu})}{\bar{\mu}}-\sum_{j=1}^{m} \log \left(-g_{j}(x(\bar{\mu}))\right) \leq-\frac{c^{T} x(\mu)}{\bar{\mu}}-\sum_{j=1}^{m} \log \left(-g_{j}(x(\mu))\right) .
\end{aligned}
$$

Adding the two inequalities gives, after rearranging the terms,

$$
\left(\frac{1}{\bar{\mu}}-\frac{1}{\mu}\right)\left(c^{T} x(\mu)-c^{T} x(\bar{\mu})\right) \leq 0 .
$$

Since $0<\bar{\mu}<\mu$ this implies $c^{T} x(\mu) \leq c^{T} x(\bar{\mu})$. Hence $-c^{T} x(\bar{\mu}) \leq-c^{T} x(\mu)$, proving the first part of the lemma.

The second part follows similarly. We have that $(x(\mu), y(\mu))$ maximizes $\phi_{B}^{d}(x, y, \mu)$. Observe that the dual objective $-c^{T} x+\sum_{j=1}^{m} y_{j} g_{j}(x)$ is just the Lagrange function $L(x, y)$ of $(C P O)$. As before, let $0<\bar{\mu}<\mu$. Now $(x(\mu), y(\mu))$ maximizes $\phi_{B}^{d}(x, y, \mu)$, and $(x(\bar{\mu}), y(\bar{\mu}))$ maximizes $\phi_{B}^{d}(x, y, \bar{\mu})$, hence

$$
\begin{aligned}
& \phi_{B}^{d}(x(\mu), y(\mu), \mu) \geq \phi_{B}^{d}(x(\bar{\mu}), y(\bar{\mu}), \mu) \\
& \phi_{B}^{d}(x(\bar{\mu}), y(\bar{\mu}), \bar{\mu}) \geq \phi_{B}^{d}(x(\mu), y(\mu), \bar{\mu}) .
\end{aligned}
$$

These are equivalent with

$$
\begin{aligned}
\frac{L(x(\mu), y(\mu))}{\mu}+\sum_{j=1}^{m} \log y_{j}(\mu) & \geq \frac{L(y(\bar{\mu}), x(\bar{\mu}))}{\mu}+\sum_{j=1}^{m} \log y_{j}(\bar{\mu}) \\
\frac{L(y(\bar{\mu}), x(\bar{\mu}))}{\bar{\mu}}+\sum_{j=1}^{m} \log y_{j}(\bar{\mu}) & \geq \frac{L(y(\mu), x(\mu))}{\bar{\mu}}+\sum_{j=1}^{m} \log y_{j}(\mu)
\end{aligned}
$$

Here we omitted the term $n(1-\log \mu)$ at both sides of the first inequality and the term $n(1-\log \bar{\mu})$ at both sides of the second inequality, since these terms cancel. Adding the two inequalities gives

$$
\left(\frac{1}{\bar{\mu}}-\frac{1}{\mu}\right)(L(y(\bar{\mu}), x(\bar{\mu}))-L(y(\mu), x(\mu))) \geq 0
$$

Hence $L(y(\bar{\mu}), x(\bar{\mu})) \geq L(y(\mu), x(\mu))$. This completes the proof.

### 6.3 Logarithmic barrier method for solving (CPO)

### 6.3.1 Introduction

Let $x$ be a strictly feasible primal solution of $(C P O)$. For given $\mu>0$ it will be shown that we can compute a good approximation of the $\mu$-center $x(\mu)$. Recall that $x(\mu)$ is the (unique) minimizer of the primal barrier function $\phi_{B}(x, \mu)$; this function will be shown to be strictly convex. Note that $\phi_{B}(x, \mu)$ is defined on the open set $\mathcal{F}^{0}$ and that minimizing this function is essentially an unconstrained optimization problem; the minimizer $x(\mu)$ is characterized by the fact that the gradient $\nabla \phi_{B}(x, \mu)$ vanishes (cf. Lemma 2.6). A natural candidate for a computational method is the generic method for unconstrained optimization as described before in Section 4.1. Starting at $x$, we move into the direction of $x(\mu)$, by Newton's method for minimizing $\phi_{B}(x, \mu)$. Thus, in the next section we will compute the Newton direction (or Newton step).
In the analysis of the method we need to quantify the distance from $x$ to $x(\mu)$. A natural way to do this is provided by the method itself. The Newton step will vanish (i.e. will be equal to the zero vector) if $x=x(\mu)$ and it will be nonzero in all other cases. As a consequence, we can use the 'length' of the Newton step as a measure for the distance of $x$ to $x(\mu)$. It is crucial for the analysis of Newton's method that this 'length' is defined appropriately. We will provide some arguments why the obvious candidate for doing this, the Euclidean norm of the Newton step, is not appropriate.

We will argue that it is much more appropriate to measure the 'length' of the Newton step with respect to the norm induced by the Hessian matrix of the barrier function. Using this norm, we show that the Newton process is quadratically convergent if $x$ is 'close' to $x(\mu)$; if $x$ is 'far' from $x(\mu)$ then damped Newton steps can be used to reach the region where the Newton process is quadratically convergent.
In this way we obtain a computationally efficient method to find a good approximation for $x(\mu)$. Having such a method it becomes easy to design an efficient algorithm for solving (CPO).

### 6.3.2 Newton step for $\phi_{B}$

Recall that our aim is to find the minimizer $x(\mu)$ of the primal barrier function $\phi_{B}$, starting at some strictly primal feasible point $x$. Also recall the idea behind Newton's method, as described in Section 4.4. This is to approximate $\phi_{B}$ by its second order Taylor polynomial at $x$ and then to use the minimizer of this Taylor polynomial-which can be calculated straightforwardly - as a new approximation for $x(\mu)$.
To construct the second order Taylor polynomial at $x$ we need the value, the gradient and the Hessian of $\phi_{B}$ at $x$. These are given by

$$
\begin{aligned}
\phi_{B}(x, \mu) & =-\frac{c^{T} x}{\mu}-\sum_{j=1}^{m} \log \left(-g_{j}(x)\right) \\
\nabla \phi_{B}(x, \mu) & =-\frac{c}{\mu}+\sum_{j=1}^{m} \frac{\nabla g_{j}(x)}{-g_{j}(x)} \\
\nabla^{2} \phi_{B}(x, \mu) & =\sum_{j=1}^{m}\left(\frac{\nabla^{2} g_{j}(x)}{-g_{j}(x)}+\frac{\nabla g_{j}(x) \nabla g_{j}(x)^{T}}{g_{j}(x)^{2}}\right) .
\end{aligned}
$$

From the last expression we see that $\nabla^{2} \phi_{B}(x, \mu)$ is positive semidefinite, because the matrices $\nabla^{2} g_{j}(x)$ and $\nabla g_{j}(x) \nabla g_{j}(x)^{T}$ are positive semidefinite and $g_{j}(x)<0$. In fact, denoting

$$
\begin{align*}
H(x, \mu) & :=\nabla^{2} \phi_{B}(x, \mu)  \tag{6.5}\\
g(x, \mu) & :=\nabla \phi_{B}(x, \mu) \tag{6.6}
\end{align*}
$$

we can even show that $H(x, \mu)$ is positive definite, provided that the logarithmic barrier function satisfies some smoothness condition that will be defined later on, in Section 6.3.4. For the moment we make the following assumption.

Assumption 6.7 For each $x \in \mathcal{F}^{0}$ the matrix $H(x, \mu)$ is positive definite.

Under this assumption $\phi_{B}(x, \mu)$ is strictly convex, by Lemma 1.50. Hence the minimizer $x(\mu)$ is unique indeed. Now the second order Taylor polynomial of $\phi_{B}(x, \mu)$ at $x$ is given by

$$
t_{2}(\Delta x)=\phi_{B}(x, \mu)+\Delta x^{T} g(x, \mu)+\frac{1}{2} \Delta x^{T} H(x, \mu) \Delta x .
$$

Since $H(x, \mu)$ is positive definite, $t_{2}(\Delta x)$ is strictly convex and has a unique minimizer. At the minimizer of $t_{2}(\Delta x)$ the gradient of $t_{2}(\Delta x)$ is zero, and thus we find the minimizer from

$$
g(x, \mu)+H(x, \mu) \Delta x=0 .
$$

Therefore, the Newton step at $x$ is given by (cf. Section 4.4)

$$
\Delta x=-H(x, \mu)^{-1} g(x, \mu)
$$

and the new iterate is

$$
x:=x+\alpha \Delta x,
$$

where $\alpha$ is the step size. If $\alpha=1$ we have a so-called full Newton step and if $\alpha<1$ we have a so-called damped Newton step.

### 6.3.3 Proximity measure

We need a tool to measure how successful a Newton step is. Ideally, one full Newton step brings us at $x(\mu)$, but this can happen only if $\phi_{B}(x, \mu)$ is quadratic and this is not true, as is obvious from the definition of $\phi_{B}(x, \mu)$. Thus we need a 'distance' or 'proximity' measure which enables us to quantify the progress on the way to the minimizer $x(\mu)$. One obvious measure is the Euclidean norm

$$
\|x-x(\mu)\|,
$$

but this measure has the obvious disadvantage that we cannot calculate it because we do not know $x(\mu)$. A good alternative is the Euclidean norm of the Newton step itself:

$$
\|\Delta x\|
$$

The last norm can be considered an approximation of $\|x-x(\mu)\|$, since - hopefully$\Delta x$ is a good approximation of $x-x(\mu)$.

Instead of the Euclidean norm we use the so-called Hessian norm and measure the 'distance' from $x$ to $x(\mu)$ by the quantity

$$
\delta(x, \mu):=\|\Delta x\|_{H}:=\sqrt{\Delta x^{T} H(x, \mu) \Delta x} .
$$

Exercise 6.1 One has

$$
\delta(x, \mu)=\sqrt{g(x, \mu)^{T} H(x, \mu)^{-1} g(x, \mu)}=\|g(x, \mu)\|_{H^{-1}} .
$$

Prove this.

Remark: The choice of $\delta(x, \mu)$ as a proximity measure will be justified by the results below. At this stage it may be worth mentioning another argument for its use. Consider the function $\Phi$ defined by

$$
\Phi(z):=\phi(A z+a),
$$

where $\phi: \mathbb{R}^{m} \rightarrow \mathbb{R}$ can be any two times differentiable function, $A$ is any $m \times m$ nonsingular matrix, $a$ a vector in $\mathbb{R}^{m}$, and where $z$ runs through all vectors such that $A z+a$ is strictly primal feasible. The Newton step with respect to $\phi(x)$ at $x$ is given by the expression

$$
\Delta x=-\nabla^{2} \phi(x)^{-1} \nabla \phi(x) .
$$

Similarly, we have the Newton step with respect to $\Phi(z)$ at $z$ :

$$
\Delta z=-\nabla^{2} \Phi(z)^{-1} \nabla \Phi(z)
$$

Taking $z=A^{-1}(x-a)$ we then have $\Delta z=A^{-1} \Delta x$, as can be verified by straightforward calculations. We express this phenomenon by saying that the Newton step is affine invariant. It is clear that the norm of $\Delta x$ is not affine invariant, because $\left\|A^{-1} \Delta x\right\|$ will in general not be equal to $\|\Delta x\|$. However, $\delta(x, \mu)$ is affine invariant!

Exercise 6.2 Prove that the Newton step and $\delta(x, \mu)$ are affine invariant.

### 6.3.4 The self-concordance property

## Introduction

Let us first recall a simple example used earlier to show that slow convergence is possible when using Newton's method.
Let $f: \mathbb{R} \rightarrow \mathbb{R}$ be defined by $f(x)=x^{2 k}$, where $k \geq 1$. Clearly, $f$ has a unique minimizer, namely $x=0$. We saw in Example 4.4, that if we apply Newton's method with full Newton steps to this function, then the rate of convergence of the iterates to the minimum is only linear, unless $k=1$ ( $f$ is quadratic).
The example suggests that we cannot expect quadratic convergence behavior of Newton's method unless it is applied to a function that is 'almost' quadratic. The smoothness condition on the primal barrier function that we are going to discuss can be understood by keeping this in mind: essentially the condition defines what we mean by saying that a function is 'almost' quadratic.

## Definition of the self-concordance property

Before we define the smoothness condition we need to introduce some notation.
Let us fix $x \in \mathcal{F}^{0}$ and $h \in \mathbb{R}^{n}$. For a fixed $\mu$ we consider the function

$$
\varphi(\alpha):=\phi_{B}(x+\alpha h, \mu),
$$

where $\alpha$ runs through all real values such that $x+\alpha h \in \mathcal{F}^{0}$. Note that $\varphi$ is strictly convex because $\phi_{B}(x, \mu)$ is strictly convex. Denoting for the moment $\phi_{B}(x, \mu)$ shortly as $\phi(x)$, we have

$$
\varphi^{\prime}(0)=\sum_{i=1}^{n} h_{i} \frac{\partial \phi(x)}{\partial x_{i}}
$$

$$
\begin{aligned}
\varphi^{\prime \prime}(0) & =\sum_{i=1}^{n} \sum_{j=1}^{n} h_{i} h_{j} \frac{\partial \phi^{2}(x)}{\partial x_{i} \partial x_{j}} \\
\varphi^{\prime \prime \prime}(0) & =\sum_{i=1}^{n} \sum_{j=1}^{n} \sum_{k=1}^{n} h_{i} h_{j} h_{k} \frac{\partial \phi^{3}(x)}{\partial x_{i} \partial x_{j} \partial x_{k}} .
\end{aligned}
$$

Note that the right hand sides in the expressions, given above, are homogeneous forms in the vector $h$, of order 1, 2 and 3 respectively. It will be convenient to use short hand notations for these forms, namely $\nabla \phi(x)[h], \nabla^{2} \phi(x)[h, h]$ and $\nabla^{3} \phi(x)[h, h, h]$ respectively. We then may write

$$
\begin{aligned}
\varphi^{\prime}(0) & =\nabla \phi(x)[h]=h^{T} \nabla \phi(x) \\
\varphi^{\prime \prime}(0) & =\nabla^{2} \phi(x)[h, h]=h^{T} \nabla^{2} \phi(x) h=\|h\|_{H}^{2} \\
\varphi^{\prime \prime \prime}(0) & =\nabla^{3} \phi(x)[h, h, h]=h^{T} \nabla^{3} \phi(x)[h] h .
\end{aligned}
$$

The last expression uses that $\nabla^{3} \phi(x)[h]$ is a square matrix of size $n \times n$. Moreover, as before, $H=\nabla^{2} \phi(x)$.
Recall that the third order Taylor expansion of $\varphi$ at 0 is given by

$$
\varphi(0)+\varphi^{\prime}(0) \alpha+\frac{1}{2} \varphi^{\prime \prime}(0) \alpha^{2}+\frac{1}{6} \varphi^{\prime \prime \prime}(0) \alpha^{3} .
$$

Thus it will be clear that the following definition, which defines the so-called selfconcordance property of $\phi$, bounds the third order term in the Taylor expansion of $\varphi$ in terms of the second order term. Although our main aim is to apply this definition to the logarithmic barrier function $\phi_{B}$ above, the definition is more general; it applies to any three times differentiable convex function with open domain. In fact, after the definition we will demonstrate it on many other simple examples.

Definition 6.8 (Self-concordance) Let $\phi$ be any three times differentiable convex function with open domain. Then $\phi$ is called $\kappa$-self-concordant, where $\kappa$ is fixed and $\kappa \geq 0$, if the inequality

$$
\left|\nabla^{3} \phi(x)[h, h, h]\right| \leq 2 \kappa\left(\nabla^{2} \phi(x)[h, h]\right)^{\frac{3}{2}}
$$

holds for any $x$ in the domain of $\phi$ and for any $h \in \mathbb{R}^{n}$.

Let $\phi$ be any three times differentiable convex function with open domain. We will say that $\phi$ is self-concordant, without specifying $\kappa$, if $\phi$ is $\kappa$-self-concordant for some $\kappa \geq 0$. Obviously, this will be the case if and only if the quotient

$$
\begin{equation*}
\frac{\left(\nabla^{3} \phi(x)[h, h, h]\right)^{2}}{\left(\nabla^{2} \phi(x)[h, h]\right)^{3}} \tag{6.7}
\end{equation*}
$$

is bounded above by $4 \kappa^{2}$ when $x$ runs through the domain of $\phi$ and $h$ through all vectors in $\mathbb{R}^{n}$. Note that the condition for $\kappa$-self-concordancy is homogeneous in $h$ : if it holds for some $h$ then it holds for any $\lambda h$, with $\lambda \in \mathbb{R}$.

Exercise 6.3 In the special case where $n=1$ the $\kappa$-self-concordancy condition reduces to

$$
\left|\phi^{\prime \prime \prime}(x)\right| \leq 2 \kappa\left(\phi^{\prime \prime}(x)\right)^{\frac{3}{2}}
$$

Prove this.

The $\kappa$-self-concordancy condition bounds the third order term in terms of the second order term in the Taylor expansion. Hence, if it is satisfied, it makes that the second order Taylor expansion locally provides a good quadratic approximation of $\phi(x)$. The latter property makes that Newton's method behaves well on self-concordant functions. This will be shown later on.

Recall that the definition of the $\kappa$-self-concordance property applies to every three times differentiable convex function with an open domain. Keeping this in mind we can already give some simple examples of self-concordant functions.

Example 6.9 [Linear function] Let $\phi(x)=\gamma+a^{T} x$, with $\gamma \in \mathbb{R}$ and $a \in \mathbb{R}^{m}$. Then

$$
\nabla \phi(x)=a, \nabla^{2} \phi(x)=0, \nabla^{3} \phi(x)=0,
$$

and we conclude that $\phi$ is 0 -self-concordant.

## Example 6.10 [Convex quadratic function] Let

$$
\phi(x)=\gamma+a^{T} x+\frac{1}{2} x^{T} A x
$$

with $\gamma$ and $a$ as before and $A=A^{T}$ positive semidefinite. Then

$$
\nabla \phi(x)=a+A x, \nabla^{2} \phi(x)=A, \nabla^{3} \phi(x)=0,
$$

and we conclude that $\phi$ is 0 -self-concordant.

We may conclude from the above examples that linear and convex quadratic functions are 0 -self-concordant.

Example 6.11 Consider the convex function $\phi(x)=x^{4}$, with $x \in \mathbb{R}$. Then

$$
\phi^{\prime}(x)=4 x^{3}, \quad \phi^{\prime \prime}(x)=12 x^{2}, \quad \phi^{\prime \prime \prime}(x)=24 x
$$

Now we have

$$
\frac{\left(\phi^{\prime \prime \prime}(x)\right)^{2}}{\left(\phi^{\prime \prime}(x)\right)^{3}}=\frac{(24 x)^{2}}{\left(12 x^{2}\right)^{3}}=\frac{1}{3 x^{4}} .
$$

Clearly the right hand side expression is not bounded if $x \rightarrow 0$, hence $\phi(x)$ is not self-concordant. *

Exercise 6.5 Let $k$ be an integer and $k>1$. Prove that $\phi(x)=x^{k}$, where $x \in \mathbb{R}$, is $\kappa$-self-concordant for some $\kappa$ only if $k \leq 2$.

Example 6.12 Now consider the convex function

$$
\phi(x)=x^{4}-\log x, \quad x>0
$$

Then

$$
\phi^{\prime}(x)=4 x^{3}-\frac{1}{x}, \quad \phi^{\prime \prime}(x)=12 x^{2}+\frac{1}{x^{2}}, \quad \phi^{\prime \prime \prime}(x)=24 x-\frac{2}{x^{3}} .
$$

Therefore,

$$
\frac{\left(\phi^{\prime \prime \prime}(x)\right)^{2}}{\left(\phi^{\prime \prime}(x)\right)^{3}}=\frac{\left(24 x-\frac{2}{x^{3}}\right)^{2}}{\left(12 x^{2}+\frac{1}{x^{2}}\right)^{3}}=\frac{\left(24 x^{4}-2\right)^{2}}{\left(12 x^{4}+1\right)^{3}} \leq \frac{\left(24 x^{4}+2\right)^{2}}{\left(12 x^{4}+1\right)^{3}}=\frac{4}{12 x^{4}+1} \leq 4
$$

This proves that $\phi(x)$ is a 1-self-concordant function.

Example 6.13 [The function $-\log x$ ] Let

$$
\phi(x)=-\log x
$$

with $0<x \in \mathbb{R}$. Then

$$
\phi^{\prime}(x)=\frac{-1}{x}, \quad \phi^{\prime \prime}(x)=\frac{1}{x^{2}}, \quad \phi^{\prime \prime \prime}(x)=\frac{-2}{x^{3}}
$$

and

$$
\frac{\left(\phi^{\prime \prime \prime}(x)\right)^{2}}{\left(\phi^{\prime \prime}(x)\right)^{3}}=\frac{\left(\frac{-2}{x^{3}}\right)^{2}}{\left(\frac{1}{x^{2}}\right)^{3}}=4
$$

Hence, $\phi$ is 1-self-concordant.
Example 6.14 [The function $-\sum_{i=1}^{n} \log x_{i}$ ] We now consider

$$
\phi(x):=-\sum_{i=1}^{n} \log x_{i}
$$

with $0<x \in \mathbb{R}^{n}$. Then, with $e$ denoting the all-one vector,

$$
\nabla \phi(x)=\frac{-e}{x}, \quad \nabla^{2} \phi(x)=\operatorname{diag}\left(\frac{e}{x^{2}}\right), \quad \nabla^{3} \phi(x)[h]=\operatorname{diag}\left(\frac{-2 h}{x^{3}}\right), \forall h \in \mathbb{R}^{n}
$$

Hence we have for any $h \in \mathbb{R}^{n}$

$$
\left|\nabla^{3} \phi(x)[h, h, h]\right|=\left|\sum_{i=1}^{n} \frac{-2 h_{i}^{3}}{x_{i}^{3}}\right|
$$

and

$$
\nabla^{2} \phi(x)[h, h]=h^{T} \operatorname{diag}\left(\frac{e}{x^{2}}\right) h=\sum_{i=1}^{n} \frac{h_{i}^{2}}{x_{i}^{2}}
$$

For any $\xi \in \mathbb{R}^{n}$ one has

$$
\begin{equation*}
\left|\sum_{i=1}^{n} \xi_{i}^{3}\right| \leq \sum_{i=1}^{n}\left|\xi_{i}\right|^{3} \leq\left(\sum_{i=1}^{n}\left|\xi_{i}^{2}\right|\right)^{\frac{3}{2}} \tag{6.8}
\end{equation*}
$$

Hence, taking $\xi_{i}:=\frac{h_{i}}{x_{i}}$ we get

$$
\left|\nabla^{3} \phi(x)[h, h, h]\right| \leq 2\left(\nabla^{2} \phi(x)[h, h]\right)^{\frac{3}{2}}
$$

proving that $\phi$ is 1 -self-concordant.

Example 6.15 [The function $\psi$ ] Let

$$
\psi(x)=x-\log (1+x),
$$

with $-1<x \in \mathbb{R}$. Then

$$
\psi^{\prime}(x)=\frac{x}{1+x}, \quad \psi^{\prime \prime}(x)=\frac{1}{(1+x)^{2}}, \quad \psi^{\prime \prime \prime}(x)=\frac{-2}{(1+x)^{3}}
$$

and it easily follows that $\psi$ is 1 -self-concordant.
Example 6.16 [The function $\Psi$ ] With $\psi$ as defined in the previous example we now consider

$$
\Psi(x):=\sum_{i=1}^{n} \psi\left(x_{i}\right)
$$

with $-e<x \in \mathbb{R}^{n}$. Then

$$
\begin{gathered}
\nabla \phi(x)=\frac{x}{e+x}, \quad \nabla^{2} \phi(x)=\operatorname{diag}\left(\frac{e}{(e+x)^{2}}\right) \\
\nabla^{3} \phi(x)[h]=\operatorname{diag}\left(\frac{-2 h}{(e+x)^{3}}\right), \forall h \in \mathbb{R}^{n}
\end{gathered}
$$

Hence we have for any $h \in \mathbb{R}^{n}$

$$
\left|\nabla^{3} \phi(x)[h, h, h]\right|=\left|\sum_{i=1}^{n} \frac{-2 h_{i}^{3}}{\left(1+x_{i}\right)^{3}}\right|
$$

and

$$
\nabla^{2} \phi(x)[h, h]=h^{T} \operatorname{diag}\left(\frac{e}{(e+x)^{2}}\right) h=\sum_{i=1}^{n} \frac{h_{i}^{2}}{\left(1+x_{i}\right)^{2}}
$$

Using (6.8) with $\xi_{i}:=h_{i} /\left(1+x_{i}\right)$ we obtain

$$
\left|\nabla^{3} \phi(x)[h, h, h]\right| \leq 2\left(\nabla^{2} \phi(x)[h, h]\right)^{\frac{3}{2}}
$$

proving that $\Psi$ is 1 -self-concordant.
Example 6.17 [Barrier of the entropy function $x \log x$ ] We consider

$$
\phi(x):=x \log x-\log x=(x-1) \log x
$$

with $0<x \in \mathbb{R}$. Then

$$
\phi^{\prime}(x)=\frac{x-1}{x}+\log x, \quad \phi^{\prime \prime}(x)=\frac{x+1}{x^{2}}, \quad \phi^{\prime \prime \prime}(x)=-\frac{x+2}{x^{3}}
$$

Hence, using also $x>0$ we may write,

$$
\frac{\left(\phi^{\prime \prime \prime}(x)\right)^{2}}{\left(\phi^{\prime \prime}(x)\right)^{3}}=\frac{\left(-\frac{x+2}{x^{3}}\right)^{2}}{\left(\frac{x+1}{x^{2}}\right)^{3}}=\frac{(x+2)^{2}}{(x+1)^{3}} \leq \frac{(2 x+2)^{2}}{(x+1)^{3}}=\frac{4}{x+1} \leq 4
$$

showing that $\phi$ will be 1-self-concordant.
Later on (in Section 7.3) we will see that also the multidimensional version of the entropy function has a 1-self-concordant barrier function.

Exercise 6.6 If $\phi$ is $\kappa$-self-concordant function with $\kappa>0$, then $\phi$ can be re-scaled by a positive scalar so that it becomes 1-self-concordant. This follows because if $\lambda$ is some positive constant then $\lambda \phi$ is $\left(\frac{\kappa}{\sqrt{\lambda}}\right)$-self-concordant. Prove this.

### 6.3.5 Properties of Newton's method

From now on we assume that $\phi(x):=\phi_{B}(x, \mu)$, for some fixed $\mu>0$, and that $\phi$ is $\kappa$-self-concordant on its domain $\mathcal{F}^{0}$, with $\kappa>0$. We are ready to state the result that if $\delta(x, \mu)$ is small enough then the Newton process is quadratically convergent.

Lemma 6.18 If $x$ is strictly primal feasible and $\mu>0$ such that $\delta:=\delta(x, \mu)<\frac{1}{\kappa}$ then $x+\Delta x$ (where $\Delta x$ denotes the Newton step at $x$ ) is strictly feasible and

$$
\delta(x+\Delta x, \mu) \leq \frac{\kappa \delta^{2}}{(1-\kappa \delta)^{2}}
$$

Proof: We omit the proof here and refer to Lemma 6.39 in Section 6.4.4 and the remark following Lemma 6.39.

Corollary 6.19 If $\delta:=\delta(x, \mu) \leq \frac{1}{3 \kappa}$ then $\delta(x+\Delta x, \mu) \leq \frac{9}{4} \kappa \delta^{2}$.
In the analysis of central-path-following methods we also need to know the effect of an update of the barrier parameter on the proximity measure. Note that the following result makes clear that this effect does not depend on the parameter $\kappa$.

Lemma 6.20 Let $x$ be strictly primal feasible and $\delta:=\delta(x, \mu)$ for some $\mu>0$. If $\mu^{+}=(1-\theta) \mu$ then

$$
\delta\left(x, \mu^{+}\right) \leq \frac{\delta+\theta \sqrt{m}}{1-\theta}
$$

*Proof: ${ }^{1}$ We have, by definition,

$$
\delta(x, \mu):=\|\Delta x\|_{H}:=\sqrt{\Delta x^{T} H(x, \mu) \Delta x} .
$$

Substituting the expression

$$
\Delta x=-H(x, \mu)^{-1} g(x, \mu)
$$

for the Newton step we get (cf. Exercise 6.1)

$$
\delta(x, \mu)=\sqrt{g(x, \mu)^{T} H(x, \mu)^{-1} g(x, \mu)}=\|g(x, \mu)\|_{H^{-1}}
$$

where (cf. (6.6) and (6.5))

$$
\begin{aligned}
& g(x, \mu)=\nabla \phi_{B}(x, \mu)=-\frac{c}{\mu}+\sum_{j=1}^{m} \frac{\nabla g_{j}(x)}{-g_{j}(x)} \\
& H:=H(x, \mu)=\nabla^{2} \phi_{B}(x, \mu)=\sum_{j=1}^{m}\left(\frac{\nabla^{2} g_{j}(x)}{-g_{j}(x)}+\frac{\nabla g_{j}(x) \nabla g_{j}(x)^{T}}{g_{j}(x)^{2}}\right) .
\end{aligned}
$$

Note that $H(x, \mu)$ does not depend on $\mu$, so

$$
H\left(x, \mu^{+}\right)=H(x, \mu) .
$$

[^4]Therefore, $\delta\left(x, \mu^{+}\right)$is given by

$$
\begin{equation*}
\delta\left(x, \mu^{+}\right)=\sqrt{g\left(x, \mu^{+}\right)^{T} H(x, \mu)^{-1} g\left(x, \mu^{+}\right)}=\left\|g\left(x, \mu^{+}\right)\right\|_{H^{-1}} \tag{6.9}
\end{equation*}
$$

We proceed by calculating $g\left(x, \mu^{+}\right)$. We may write

$$
\begin{aligned}
g\left(x, \mu^{+}\right) & =-\frac{c}{\mu^{+}}+\sum_{j=1}^{m} \frac{\nabla g_{j}(x)}{-g_{j}(x)} \\
& =-\frac{c}{(1-\theta) \mu}+\sum_{j=1}^{m} \frac{\nabla g_{j}(x)}{-g_{j}(x)} \\
& =\frac{1}{1-\theta}\left(-\frac{c}{\mu}+\sum_{j=1}^{m} \frac{\nabla g_{j}(x)}{-g_{j}(x)}-\theta \sum_{j=1}^{m} \frac{\nabla g_{j}(x)}{-g_{j}(x)}\right) \\
& =\frac{1}{1-\theta}\left(g(x, \mu)-\theta \sum_{j=1}^{m} \frac{\nabla g_{j}(x)}{-g_{j}(x)}\right) \\
& =\frac{1}{1-\theta}(g(x, \mu)-\theta J e)
\end{aligned}
$$

where $J$ denotes the matrix

$$
J=\left(\frac{\nabla g_{1}(x)}{-g_{1}(x)} \cdots \frac{\nabla g_{m}(x)}{-g_{m}(x)}\right) .
$$

and $e$ the all-one vector. Substituting in (6.9) we get

$$
\delta\left(x, \mu^{+}\right)=\left\|g\left(x, \mu^{+}\right)\right\|_{H^{-1}}=\frac{1}{1-\theta}\|g(x, \mu)-\theta J e\|_{H^{-1}}
$$

Using the triangle inequality yields

$$
\delta\left(x, \mu^{+}\right) \leq \frac{1}{1-\theta}\left(\|g(x, \mu)\|_{H^{-1}}+\theta\|J e\|_{H^{-1}}\right)=\frac{\delta+\theta\|J e\|_{H^{-1}}}{1-\theta}
$$

Thus the lemma will follow if $\|J e\|_{H^{-1}} \leq \sqrt{m}$. We have

$$
\|J e\|_{H^{-1}}^{2}=e^{T} J^{T} H(x, \mu)^{-1} J e
$$

Now observe that

$$
H(x, \mu)=\sum_{j=1}^{m}\left(\frac{\nabla^{2} g_{j}(x)}{-g_{j}(x)}+\frac{\nabla g_{j}(x) \nabla g_{j}(x)^{T}}{g_{j}(x)^{2}}\right) \succeq \frac{\nabla g_{j}(x) \nabla g_{j}(x)^{T}}{g_{j}(x)^{2}}=J J^{T}
$$

where $H(x, \mu) \succeq J J^{T}$ means that the matrix $H(x, \mu)-J J^{T}$ is PSD. Hence

$$
H(x, \mu)^{-1} \preceq\left(J J^{T}\right)^{+}
$$

where $\left(J J^{T}\right)^{+}$denotes the generalized inverse of $J J^{T}$. Using this we find

$$
\|J e\|_{H^{-1}}^{2} \leq e^{T} J^{T}\left(J J^{T}\right)^{+} J e
$$

Since $J^{T}\left(J J^{T}\right)^{+} J$ is a projection matrix, we will have

$$
e^{T} J^{T}\left(J J^{T}\right)^{+} J e \leq e^{T} e=m
$$

and hence the proof is complete.

Theorem 6.21 Let $x^{+}:=x+\Delta x$ and $\mu^{+}=(1-\theta) \mu$, where $\theta=\frac{1}{30 k \sqrt{m}}$. Then

$$
\delta(x, \mu) \leq \frac{1}{3 \kappa} \Rightarrow \delta\left(x^{+}, \mu^{+}\right) \leq \frac{1}{3 \kappa}
$$

Proof: Using Lemma 6.18 and Lemma 6.20 we may write

$$
\begin{aligned}
\delta\left(x^{+}, \mu^{+}\right) & \leq \frac{9}{4} \kappa \delta\left(x, \mu^{+}\right)^{2} \\
& \leq \frac{9}{4} \kappa\left(\frac{1}{1-\frac{1}{30 \kappa \sqrt{m}}}\left(\frac{1}{3 \kappa}+\frac{1}{30 \kappa}\right)\right)^{2} \\
& \leq \frac{1}{3 \kappa} .
\end{aligned}
$$

This proves the theorem.

### 6.3.6 Logarithmic barrier algorithm with full Newton steps

We are now ready to state our first algorithm.

## Logarithmic Barrier Algorithm with full Newton steps

## Input:

A proximity parameter $\tau, 0 \leq \tau<1$;
an accuracy parameter $\epsilon>0$;
$x^{0} \in \mathcal{F}^{0}$ and $\mu^{0}>0$ such that $\delta\left(x^{0}, \mu^{0}\right) \leq \tau$;
a fixed barrier update parameter $\theta, 0<\theta<1$.
begin
$x:=x^{0} ; \mu:=\mu^{0} ;$
while $m \mu \geq \epsilon$ do
begin
$\mu:=(1-\theta) \mu ;$
$x:=x+\Delta x(\Delta x$ is the Newton step at $x)$
end
end

We prove the following theorem.

Theorem 6.22 If $\tau=\frac{1}{3 \kappa}$ and $\theta=\frac{1}{30 \kappa \sqrt{m}}$, then the Logarithmic Barrier Algorithm with full Newton steps requires at most

$$
\left\lceil 30 \kappa \sqrt{m} \log \frac{m \mu^{0}}{\epsilon}\right\rceil
$$

iterations. The output is a strictly primal feasible $x$ such that $x$ is $\epsilon$-optimal.
Proof: By Theorem 6.21 the property $\delta(x, \mu) \leq \frac{1}{3 \kappa}$ is maintained in the course of the algorithm. Thus each (full) Newton step will yield a strictly feasible point, by Lemma 6.18. At each iteration the barrier parameter is reduced by the factor $1-\theta$. Hence, after $k$ iterations we will have

$$
m \mu=(1-\theta)^{k} m \mu^{0} .
$$

Using this, one easily deduces that after no more than

$$
\begin{equation*}
\left\lceil\frac{1}{\theta} \log \frac{m \mu^{0}}{\epsilon}\right\rceil \tag{6.10}
\end{equation*}
$$

iterations the algorithm will have stopped. Substitution of the value of $\theta$ in the theorem yields the desired bound.

## Example 6.23 [Logarithmic Barrier Method with Full Newton Steps 1]

Consider the obvious minimization problem

$$
\min \{x: x \geq 0\} .
$$

We solve this problem by using the Logarithmic Barrier Algorithm with Full Newton Steps. First we transfer the function into the standard form:

$$
\min \{x:-x \leq 0\} .
$$

The logarithmic barrier function for this problem is given by

$$
\phi_{B}(x, \mu)=\frac{x}{\mu}-\log x .
$$

The function $\phi_{B}(x, \mu)$ is a 1 -self-concordant function. Therefore we take

$$
\tau=\frac{1}{3 \kappa}=\frac{1}{3}, \quad \theta=\frac{1}{30 \kappa \sqrt{m}}=\frac{1}{30} .
$$

We choose $\epsilon=0.5, \mu^{0}=0.8$ and $x^{0}=1$. Then the Logarithmic Barrier Algorithm with full Newton Steps requires at most

$$
\left\lceil 30 \log \frac{\mu^{0}}{\epsilon}\right\rceil=15
$$

iterations to reach an $\epsilon$-optimal solution $x$. We need to check if

$$
\delta\left(x^{0}, \mu^{0}\right)=\sqrt{\triangle x^{T} H\left(x^{0}, \mu^{0}\right) \triangle x} \leq \tau .
$$

For this we perform the following calculations:

$$
\begin{aligned}
g(x, \mu) & =\nabla \phi_{B}(x, \mu)=\frac{1}{\mu}-\frac{1}{x} \\
H(x, \mu) & =\nabla^{2} \phi_{B}(x, \mu)=\frac{1}{x^{2}} \\
H(x, \mu)^{-1} & =x^{2} .
\end{aligned}
$$

This implies

$$
\triangle x=-H\left(x^{0}, \mu^{0}\right)^{-1} g\left(x^{0}, \mu^{0}\right)=-1 \cdot \frac{1}{4}=-\frac{1}{4}
$$

and, as a consequence:

$$
\delta\left(x^{0}, \mu^{0}\right)=|\Delta x|=\frac{1}{4} \leq \frac{1}{3} .
$$

This means we can start the iterations.

## Iteration 1

Because

$$
m \mu^{0}=0.8 \geq \epsilon
$$

we start by computing the new $\mu$ and the new $x$ :

$$
\begin{aligned}
\mu^{1} & =(1-\theta) \mu^{0}=0.773333 \\
g\left(x^{0}, \mu^{1}\right) & =0.293103 \\
H\left(x^{0}, \mu^{1}\right) & =1 \\
H\left(x^{0}, \mu^{1}\right)^{-1} & =1 \\
x^{1} & =x^{0}+\Delta x=1-1 \cdot 0.293103=0.706896
\end{aligned}
$$

## Iteration 2

First we check

$$
m \mu^{1}=0.773333 \geq \epsilon
$$

Therefore

$$
\begin{aligned}
\mu^{2} & =(1-\theta) \mu^{1}=0.747556 \\
g\left(x^{1}, \mu^{2}\right) & =0.07694 \\
H\left(x^{1}, \mu^{2}\right) & =2.00119 \\
H\left(x^{1}, \mu^{2}\right)^{-1} & =0.499703 \\
x^{2} & =x^{1}+\triangle x=0.706896+0.038448=0.745344 .
\end{aligned}
$$

The next iterations are shown in the following tableaus:

| Iteration: | 3 | 4 | 5 | 6 |
| :--- | :---: | :---: | :---: | :---: |
| $\mu$ | 0.722637 | 0.698549 | 0.675264 | 0.652755 |
| $g(x, \mu)$ | 0.042158 | 0.04635 | 0.047759 | 0.049419 |
| $H(x, \mu)$ | 1.800057 | 1.918747 | 2.053899 | 2.19795 |
| $H(x, \mu)^{-1}$ | 0.555538 | 0.521174 | 0.486879 | 0.454969 |
| $\triangle x$ | -0.02342 | -0.02416 | -0.02325 | -0.02248 |
| $x$ | 0.721924 | 0.697767 | 0.674514 | 0.65203 |


| Iteration: | 7 | 8 | 9 | 10 |
| :--- | :---: | :---: | :---: | :---: |
| $\mu$ | 0.630997 | 0.609964 | 0.589631 | 0.569977 |
| $g(x, \mu)$ | 0.051122 | 0.052885 | 0.054709 | 0.056595 |
| $H(x, \mu)$ | 2.352149 | 2.517163 | 2.69375 | 2.882732 |
| $H(x, \mu)^{-1}$ | 0.425143 | 0.397273 | 0.371229 | 0.346893 |
| $\triangle x$ | -0.02173 | -0.02101 | -0.02031 | -0.01963 |
| $x$ | 0.630296 | 0.609286 | 0.588976 | 0.569344 |
| Iteration: | 11 | 12 | 13 | 14 |
| $\mu$ | 0.550978 | 0.532612 | 0.514858 | 0.497696 |
| $g(x, \mu)$ | 0.058547 | 0.060566 | 0.062654 | 0.064815 |
| $H(x, \mu)$ | 3.084969 | 3.301394 | 3.533002 | 3.780858 |
| $H(x, \mu)^{-1}$ | 0.324152 | 0.302902 | 0.283045 | 0.26449 |
| $\triangle x$ | -0.018989 | 0.01835 | -0.01773 | -0.01714 |
| $x$ | 0.550366 | 0.53202 | 0.514286 | 0.497143 |

We can see that after the fourteenth iteration $m \mu$ became less than $\epsilon$, hence $x^{14}=0.497143$ is $\epsilon$-optimal. *

## Example 6.24 [Logarithmic Barrier Method with Full Newton Steps 2]

We apply the Logarithmic Barrier Algorithm with Full Newton Steps to the minimization problem

$$
\min \left\{x^{4}: x \geq 0\right\}
$$

The standard form of the problem is:

$$
\min \left\{x^{4}:-x \leq 0\right\}
$$

The logarithmic barrier function for this problem is given by

$$
\phi_{B}(x, \mu)=\frac{x^{4}}{\mu}-\log x
$$

and this function is 1-self-concordant (cf. Example 6.12). Therefore we take again

$$
\tau=\frac{1}{3 \kappa}=\frac{1}{3}, \quad \theta=\frac{1}{30 \kappa \sqrt{m}}=\frac{1}{30} .
$$

We choose $\epsilon=1, \mu^{0}=3$ and $x^{0}=1$. Then the algorithm requires at most

$$
\left\lceil 30 \log \frac{\mu^{0}}{\epsilon}\right\rceil=33
$$

iterations to reach an $\epsilon$-optimal solution $x$. We have to check if:

$$
\delta\left(x^{0}, \mu^{0}\right)=\sqrt{\triangle x^{T} H\left(x^{0}, \mu^{0}\right) \triangle x} \leq \tau
$$

For this we have to calculate

$$
\begin{aligned}
g(x, \mu) & =\nabla \phi_{B}(x, \mu)=\frac{4 x^{3}}{\mu}-\frac{1}{x} \\
H(x, \mu) & =\nabla^{2} \phi_{B}(x, \mu)=\frac{12 x^{2}}{\mu}+\frac{1}{x^{2}} \\
H(x, \mu)^{-1} & =\frac{x^{2}+\mu}{12 x^{2}+1} .
\end{aligned}
$$

This implies

$$
\triangle x=-H\left(x^{0}, \mu^{0}\right)^{-1} g\left(x^{0}, \mu^{0}\right)=-\frac{1}{5} \cdot \frac{1}{3}=-\frac{1}{15},
$$

whence

$$
\delta\left(x^{0}, \mu^{0}\right)=\sqrt{\frac{1}{45}} \approx 0.15 \leq \frac{1}{3} .
$$

This means we can start the iterations.

## Iteration 1

Because

$$
m \mu^{0}=3 \geq \epsilon
$$

we start with computing the new $\mu$ and the new $x$ :

$$
\begin{aligned}
\mu^{1} & =(1-\theta) \mu^{0}=2.9 \\
g\left(x^{0}, \mu^{1}\right) & =0.37931 \\
H\left(x^{0}, \mu^{1}\right) & =5.137931 \\
H\left(x^{0}, \mu^{1}\right)^{-1} & =0.194631 \\
x^{1} & =x^{0}+\triangle x=1-0.07383=0.926174 \\
f\left(x^{1}\right) & =0.735819 .
\end{aligned}
$$

## Iteration 2

First we check

$$
m \mu^{1}=2.9 \geq \epsilon .
$$

Therefore

$$
\begin{aligned}
\mu^{2} & =(1-\theta) \mu^{1}=2.803333 \\
g\left(x^{1}, \mu^{2}\right) & =0.0539 \\
H\left(x^{1}, \mu^{2}\right) & =4.837685 \\
H\left(x^{1}, \mu^{2}\right)^{-1} & =0.20671 \\
x^{2} & =x^{1}+\triangle x=0.926174-0.01114 \\
f\left(x^{2}\right) & =0.701046 .
\end{aligned}
$$

We will also give the last two iterations.

## Iteration 32

We have

$$
m \mu^{31}=1.048818 \geq \epsilon,
$$

and thus

$$
\begin{aligned}
\mu^{32} & =(1-\theta) \mu^{31}=1.013858 \\
g\left(x^{31}, \mu^{32}\right) & =0.0484 \\
H\left(x^{31}, \mu^{32}\right) & =8.013924 \\
H\left(x^{31}, \mu^{32}\right)^{-1} & =0.124783 \\
x^{32} & =x^{31}+\triangle x=0.715609-0.00604=0.70957 \\
f\left(x^{32}\right) & =0.253502 .
\end{aligned}
$$

## Iteration 33

Still, we have

$$
m \mu^{32}=1.013858 \geq \epsilon,
$$

so we need one more iteration:

$$
\begin{aligned}
\mu^{33} & =(1-\theta) \mu^{32}=0.980063 \\
g\left(x^{32}, \mu^{33}\right) & =0.048812 \\
H\left(x^{32}, \mu^{33}\right) & =8.150924 \\
H\left(x^{32}, \mu^{33}\right)^{-1} & =0.122685 \\
x^{33} & =x^{32}+\triangle x=0.0957-0.00599=0.703582 \\
f\left(x^{33}\right) & =0.245052
\end{aligned}
$$

Now, we have

$$
m \mu^{33}=0.980063<\epsilon .
$$

Hence, $x^{33}$ is $\epsilon$-optimal.

### 6.3.7 Logarithmic barrier algorithm with damped Newton steps

It is clear that the value of the proximity parameter $\tau$ in the Logarithmic Barrier Algorithm with full Newton steps (see page 134) may be quite small. This means that the algorithm keeps the iterates very close to the central path. As a consequence, the barrier update parameter is also very small and in practice the algorithm will progress very slowly.
One obvious way to speed up the algorithm is to use larger values for the parameter $\theta$. But then, after a barrier update, the proximity value $\delta\left(x, \mu^{+}\right)$will in general be so large that we have no guarantee that a full Newton step will be feasible (cf. Lemma 6.18). Therefore, when larger barrier updates are applied we need to damp the Newton step by a factor $\alpha$ (say, with $0 \leq \alpha<1$ ) so that $x+\alpha \Delta x$ will be feasible again. We repeatedly take such damped steps until the iterate reaches the vicinity of $x\left(\mu^{+}\right)$. Then we update the barrier parameter again, and so on, until the barrier parameter reaches the threshold value $\epsilon / n$.
The analysis of an algorithm with damped Newton steps uses that the primal barrier function $\phi_{B}(x, \mu)$ is strictly convex and has $x(\mu)$ as a minimizer. We can show that if the damping factor (or step size) $\alpha$ is chosen appropriately then a damped Newton step decreases the barrier function by at least some fixed amount. Hence, after finitely many damped Newton steps we will reach the vicinity of $x(\mu)$. We have the following result.

Lemma 6.25 Let $x$ be strictly feasible, $\mu>0$ and $\delta:=\delta(x, \mu)$. If $\alpha=\frac{1}{1+\kappa \delta}$ then

$$
\phi_{B}(x, \mu)-\phi_{B}(x+\alpha \Delta x, \mu) \geq \frac{1}{\kappa^{2}} \psi(\kappa \delta) .
$$

*Proof: The proof is postponed. The result follows from Lemma 6.38.

Note that as long as $\delta(x, \mu) \geq \frac{1}{3 \kappa}$, and $x$ is outside the region around $x(\mu)$ where the Newton process is quadratically convergent (cf. Corollary 6.19), we have

$$
\frac{1}{\kappa^{2}} \psi(\kappa \delta) \geq \frac{1}{\kappa^{2}} \psi\left(\frac{1}{3}\right)=\frac{0.0457}{\kappa^{2}}>\frac{1}{22 \kappa^{2}}
$$

This shows that the barrier function decreases with at least some fixed amount, depending on $\kappa$ but not on the present iterate.
We can now state our second algorithm.

## Logarithmic Barrier Algorithm with Damped Newton Steps

```
Input:
    A proximity parameter \(\tau, 0 \leq \tau<1\);
    an accuracy parameter \(\epsilon>0\);
    \(x^{0} \in \mathcal{F}^{0}\) and \(\mu^{0}>0\) such that \(\delta\left(x^{0}, \mu^{0}\right) \leq \tau\);
    a damping factor (or step size) \(\alpha, 0 \leq \alpha<1\);
    a fixed barrier update parameter \(\theta, 0<\theta<1\).
begin
    \(x:=x^{0} ; \mu:=\mu^{0}\);
    while \(m \mu \geq \epsilon\) do
    begin
        \(\mu:=(1-\theta) \mu ;\)
        while \(\delta(x, \mu) \geq \tau\) do
        begin
            \(x:=x+\alpha \Delta x ;\)
            (The damping factor \(\alpha\) must be such that \(\phi_{B}(x, \mu)\)
            decreases sufficiently. This can be reached by tak-
            ing the default value is \(\frac{1}{1+\kappa \delta(x, \mu)}\). Larger reductions
            can be realized by performing a line search.)
        end
    end
end
```

We refer to the first while-loop in the algorithm as the outer loop and to the second while-loop as the inner loop. Each execution of the outer loop is called an outer iteration and each execution of the inner loop an inner iteration. The required number of outer iterations depends only on the dimension $n$ of the problem, on $\mu^{0}$, on $\epsilon$, and on the (fixed) barrier update parameter $\theta$. This number immediately can be bounded above by the number

$$
\left\lceil\frac{1}{\theta} \log \frac{m \mu^{0}}{\epsilon}\right\rceil
$$

given in (6.10), by using the same argument. The main task in the analysis of the algorithm is to derive an upper bound for the number of iterations in the inner loop. In this respect the following result is important.

Lemma 6.26 Each inner loop requires at most

$$
\left\lceil\frac{22 \theta}{(1-\theta)^{2}}\left(\theta \kappa^{2} m+\frac{5}{2} \kappa \sqrt{m}\right)+\frac{22}{3}\right\rceil
$$

inner iterations.
*Proof: The proof needs some other lemmas that estimate barrier function values and objective values in the region of quadratic convergence around the $\mu$-center. We refer to Theorem 2.10 (page 61) and its proof in Den Hertog [13]. For a similar result and its proof we refer to Section 6.4.6.

Combining the bounds in Lemma 6.26 and (6.10) we obtain our main result. Omitting the integer brackets we have:

Theorem 6.27 After at most

$$
\left(\frac{22}{(1-\theta)^{2}}\left(\theta \kappa^{2} m+\frac{5}{2} \kappa \sqrt{m}\right)+\frac{22}{3 \theta}\right) \log \frac{m \mu^{0}}{\epsilon}
$$

damped Newton steps the Logarithmic Barrier Algorithm with Damped Newton Steps yields a strictly primal feasible solution $x$ which is $\epsilon$-optimal.

Proof: Obvious.

If we take $\theta=\frac{\nu}{\sqrt{m}}$ for some fixed constant $\nu$ then the bound of Theorem 6.27 becomes

$$
\mathcal{O}\left(\kappa^{2} \sqrt{m} \log \frac{m \mu^{0}}{\epsilon}\right)
$$

If $\theta$ is taken independent of $n$, e.g. $\theta=\frac{1}{2}$, then the bound becomes

$$
\mathcal{O}\left(\kappa^{2} m \log \frac{m \mu^{0}}{\epsilon}\right) .
$$

## Example 6.28 [Damped Newton Steps 1]

We consider the same problem as in Example 6.23:

$$
\min \{x:-x \leq 0\} .
$$

Thus the logarithmic barrier function is given by

$$
\phi_{B}(x, \mu)=\frac{x}{\mu}-\log x .
$$

We take the same $\tau=\frac{1}{3}$, but we take a larger value for the parameter $\theta$. In this example we will use the value $\theta=0.25$ We start again from point $x^{0}=1$ with $\mu^{0}=0.8$. For $\epsilon$ we take the value 0.5 . Then the Logarithmic Barrier Algorithm with Damped Newton Steps requires at most

$$
\left\lceil\frac{1}{\theta} \log \frac{\mu^{0}}{\epsilon}\right\rceil=2
$$

iterations of the outer loop to reach an $\epsilon$-optimal $x$. Each inner loop requires at most

$$
\left\lceil\frac{22 \theta}{(1-\theta)^{2}}\left(\theta+\frac{5}{2}\right)+\frac{22}{3}\right\rceil=14
$$

iterations.
From Example 6.23 we know that

$$
\delta\left(x^{0}, \mu^{0}\right) \leq \tau
$$

so we can start with the first iteration.

## Iteration 1

Because

$$
m \mu^{0} \geq \epsilon
$$

we compute the new $\mu$ :

$$
\mu^{1}=(1-\theta) \mu^{0}=0.6
$$

First we have to know if $\delta\left(x^{0}, \mu^{1}\right) \geq \tau$. This is the case, because

$$
\delta\left(x^{0}, \mu^{1}\right)=\sqrt{\triangle x H\left(x^{0}, \mu^{1}\right) \triangle x}=0.666667
$$

We now can compute the new $x$ :

$$
\begin{aligned}
g\left(x^{0}, \mu^{1}\right) & =0.666667 \\
H\left(x^{0}, \mu^{1}\right) & =1 \\
H\left(x^{0}, \mu^{1}\right)^{-1} & =1 \\
\triangle x & =-0.66667 \\
\alpha & =\frac{1}{1+\delta\left(x^{0}, \mu^{1}\right)}=0.6 \\
x^{1} & =x^{0}+\alpha \triangle x=0.6
\end{aligned}
$$

Now we have to see if $\delta\left(x^{1}, \mu^{1}\right) \geq \tau$.

$$
\delta\left(x^{1}, \mu^{1}\right)=\sqrt{\triangle x H\left(x^{1}, \mu^{1}\right) \triangle x} \approx 0^{\prime}
$$

Because $\delta\left(x^{1}, \mu^{1}\right) \leq \tau$ we first update $\mu$ before performing iteration 2 .

## Iteration 2

We can see that $m \mu^{1} \geq \epsilon$ so we can start this iteration. First we compute the new $\mu$.

$$
\mu^{2}=(1-\theta) \mu^{1}=0.45
$$

Now we check if $\delta\left(x^{1}, \mu^{2}\right) \geq \tau$.

$$
\delta\left(x^{1}, \mu^{2}\right)=\sqrt{\triangle x H\left(x^{1}, \mu^{2}\right) \triangle x}=0.3333
$$

This means we can compute the new $x$ :

$$
\begin{aligned}
g\left(x^{1}, \mu^{2}\right) & =0.555556 \\
H\left(x^{1}, \mu^{2}\right) & =2.777778 \\
H\left(x^{1}, \mu^{2}\right)^{-1} & =0.36 \\
\triangle x & =-0.2 \\
\alpha & ==0.75 \\
x^{2} & =0.45
\end{aligned}
$$

We now reached an $\epsilon$-optimal $x$. We have used 2 iterations ("outer loops") and this is exactly the number of outer loops we could expect. Note however that the number of inner iterations is only 2 , which is far less than expected from the theory.

## Example 6.29 [Damped Newton Steps 2]

We consider the same problem as in Example 6.24:

$$
\min \left\{x^{4}:-x \leq 0\right\}
$$

Thus the logarithmic barrier function is given by

$$
\phi_{B}(x, \mu)=\frac{x^{4}}{\mu}-\log x
$$

As in the previous example we take $\tau=\frac{1}{3}$ and $\theta=0.25$. We start again from the point $x^{0}=1$ with $\mu^{0}=3$. For $\epsilon$ we take the value 1 . Then the algorithm requires at most

$$
\left\lceil\frac{1}{\theta} \log \frac{\mu^{0}}{\epsilon}\right\rceil=5
$$

iterations of the outer loop to reach an $\epsilon$-optimal $x$. Each inner loop requires at most

$$
\left\lceil\frac{22 \theta}{(1-\theta)^{2}}\left(\theta+\frac{5}{2}\right)+\frac{22}{3}\right\rceil=14
$$

iterations.
From Example 6.24 we know that

$$
\delta\left(x^{0}, \mu^{0}\right) \leq \tau
$$

so we can start with the first iteration.

## Iteration 1

Because

$$
m \mu^{0} \geq \epsilon
$$

we compute a new $\mu$ :

$$
\mu^{1}=(1-\theta) \mu^{0}=2.25
$$

First we have to compute $\delta\left(x^{0}, \mu^{1}\right)$.

$$
\delta\left(x^{0}, \mu^{1}\right)=\sqrt{\triangle x H\left(x^{0}, \mu^{1}\right) \triangle x}=0.309058<\tau
$$

So, again we can decrease the value of $\mu$.

$$
\mu^{2}=(1-\theta) \mu^{1}=1.6875
$$

Now, we have

$$
\delta\left(x^{0}, \mu^{2}\right)=\sqrt{\triangle x H\left(x^{0}, \mu^{2}\right) \triangle x}=0.481169 \geq \tau
$$

We now can compute the new $x$.

$$
\begin{aligned}
g\left(x^{0}, \mu^{2}\right) & =1.37037 \\
H\left(x^{0}, \mu^{2}\right) & =8.11111 \\
H\left(x^{0}, \mu^{2}\right)^{-1} & =0.123288 \\
\triangle x & =-0.16895 \\
\alpha & =\frac{1}{1+\delta\left(x^{0}, \mu^{2}\right)}=0.675142 \\
x^{1} & =x^{0}+\alpha \triangle x=0.885935 \\
f\left(x^{1}\right) & =0.616038
\end{aligned}
$$

Now we have to check if $\delta\left(x^{1}, \mu^{2}\right) \geq \tau$.

$$
\delta\left(x^{1}, \mu^{2}\right)=\sqrt{\triangle x H\left(x^{1}, \mu^{2}\right) \triangle x}=0.198409
$$

Because $\delta\left(x^{1}, \mu^{2}\right) \leq \tau$ we update $\mu$ and then perform outer iteration 2 .

## Iteration 2

We can see that $m \mu^{2} \geq \epsilon$ so we can start this iteration. First we compute the new $\mu$.

$$
\mu^{3}=(1-\theta) \mu^{2}=1.265625
$$

The fact that

$$
\delta\left(x^{1}, \mu^{3}\right)=\sqrt{\triangle x H\left(x^{1}, \mu^{3}\right) \triangle x}=0.362063 \geq \tau
$$

means we can compute the new $x$.

$$
\begin{aligned}
g\left(x^{1}, \mu^{3}\right) & =1.068908 \\
H\left(x^{1}, \mu^{3}\right) & =8.71591 \\
H\left(x^{1}, \mu^{3}\right)^{-1} & =0.114733 \\
\triangle x & =-0.12264 \\
\alpha & ==0.734181 \\
x^{2} & =0.795896 \\
f\left(x^{2}\right) & =0.401259
\end{aligned}
$$

Because

$$
\delta\left(x^{2}, \mu^{3}\right)=\sqrt{\triangle x H\left(x^{2}, \mu^{3}\right) \triangle x}=0.122348<\tau
$$

we start the third outer iteration.

## Iteration 3

Since $m \mu^{3} \geq \epsilon$, we compute the new $\mu$.

$$
\mu^{4}=(1-\theta) \mu^{3}=0.949219
$$

The distance

$$
\delta\left(x^{2}, \mu^{4}\right)=\sqrt{\triangle x H\left(x^{2}, \mu^{4}\right) \triangle x}=0.280366
$$

is still smaller than $\tau$, so we can decrease the value of $\mu$ again.

$$
\mu^{5}=(1-\theta) \mu^{4}=0.949219
$$

This is, in fact, already the fourth outer iteration. Now, we can see that

$$
\delta\left(x^{2}, \mu^{5}\right)=\sqrt{\triangle x H\left(x^{2}, \mu^{5}\right) \triangle x}=0.450248 \geq \tau
$$

Hence we can compute the new $x$.

$$
\begin{aligned}
g\left(x^{2}, \mu^{5}\right) & =1.576259 \\
H\left(x^{2}, \mu^{5}\right) & =12.25607 \\
H\left(x^{2}, \mu^{5}\right)^{-1} & =0.081592 \\
\triangle x & =-0.12861 \\
\alpha & ==0.689537 \\
x^{3} & =0.707214 \\
f\left(x^{3}\right) & =0.250152
\end{aligned}
$$

Because

$$
\delta\left(x^{3}, \mu^{5}\right)=\sqrt{\triangle x H\left(x^{3}, \mu^{5}\right) \triangle x}=0.177549<\tau
$$

we can end this outer iteration. Because $m \mu^{5}<\epsilon$, we reached an $\epsilon$-optimal $x$ in four outer iterations and by using only three inner iterations.

## $6.4{ }^{*}$ More on self-concordancy

### 6.4.1 Introduction

In this section we derive some properties of self-concordant functions. One of our aims is to provide proofs of some results that were stated before without proofs (especially Lemma 6.18 and Lemma 6.25) in a somewhat more general setting. We also present an efficient algorithm to find a minimizer of a $\kappa$-self-concordant function, if it exists.

We will deal with a function $\phi: \mathcal{D} \rightarrow \mathbb{R}$, where the domain $\mathcal{D}$ is an open and convex subset of $\mathbb{R}^{n}$, and we will assume that $\phi$ is closed convex and $\kappa$-self-concordant. We did not deal with the notion of a closed convex function so far, therefore we start with a definition.

Definition 6.30 A function $\phi: \mathcal{D} \rightarrow \mathbb{R}$ is called closed if its epigraph is closed. If, moreover, $\phi$ is convex then $\phi$ is called a closed convex function.

Lemma 6.31 For any point $\bar{x}$ on the boundary of the domain $\mathcal{D}$ of $\phi$ and for any sequence $\left\{x_{k}\right\}_{k=0}^{\infty}$ in the domain that converges to $\bar{x}$ we have $\phi\left(x_{k}\right) \rightarrow \infty$.

Proof: Consider the sequence $\left\{\phi\left(x_{k}\right)\right\}_{k=0}^{\infty}$. Assume that it is bounded above. Then it has a limit point $\bar{\phi}$. Of course, we can think that this is the unique limit point of the sequence. Therefore,

$$
z_{k}:=\left(x_{k}, \phi\left(x_{k}\right)\right) \rightarrow(\bar{x}, \bar{\phi}) .
$$

Note that $z_{k}$ belongs to the epigraph of $\phi$. Since $\phi$ is a closed function, then also $(\bar{x}, \bar{\phi})$ belongs to the epigraph. But this is a contradiction since $\bar{x}$ does not belong to the domain of $\phi$.

We conclude that, since the function $\phi$ considered in this section is closed convex, it has the property that $\phi(x)$ approaches infinity when $x$ approaches the boundary of the domain $\mathcal{D}$. This is also expressed by saying that $\phi$ is a barrier function on $\mathcal{D}$. In fact, the following exercise makes clear that the barrier property is equivalent to the closedness property.

Exercise 6.7 Let the function $\phi: \mathcal{D} \rightarrow \mathbb{R}$ have the property that it becomes unbounded $(+\infty)$ when approaching the boundary of its open domain $\mathcal{D}$. Then $\phi$ is closed.

We also assumed that $\phi$ is $\kappa$-self-concordant. Thus $\phi$ is three times continuously differentiable and the inequality

$$
\begin{equation*}
\left|\nabla^{3} \phi(x)[h, h, h]\right| \leq 2 \kappa\left(\nabla^{2} \phi(x)[h, h]\right)^{\frac{3}{2}} \tag{6.11}
\end{equation*}
$$

holds for any $x \in \mathcal{D}$ and for any $h \in \mathbb{R}^{n}$, where $\kappa$ is fixed and $\kappa \geq 0$.
We will denote

$$
g(x):=\nabla \phi(x), \forall x \in \mathcal{D}
$$

and

$$
H(x):=\nabla^{2} \phi(x), \forall x \in \mathcal{D}
$$

For any $v \in \mathbb{R}^{n}$, the local Hessian norm of $v$ at $x \in \mathcal{D}$ is in this section denoted as $\|v\|_{x}$. Thus

$$
\|v\|_{x}:=\sqrt{v^{T} H(x) v} .
$$

Using this notation, the inequality (6.11) can be written as

$$
\left|\nabla^{3} \phi(x)[h, h, h]\right| \leq 2 \kappa\|h\|_{x}^{3} .
$$

Let us now first point out an equivalent formulation of the self-concordance property.
Lemma 6.32 A three times differentiable closed convex function $\phi$ with open domain $\mathcal{D}$ is $\kappa$-self-concordant if and only if

$$
\left|\nabla^{3} \phi(x)\left[h_{1}, h_{2}, h_{3}\right]\right| \leq 2 \kappa\left\|h_{1}\right\|_{x}\left\|h_{2}\right\|_{x}\left\|h_{3}\right\|_{x}
$$

holds for any $x \in \mathcal{D}$ and all $h_{1}, h_{2}, h_{3} \in \mathbb{R}^{n}$.
Proof: This statement is nothing but a general property of three-linear forms. For the proof we refer to Lemma A. 2 in the Appendix.

We proceed with an interesting, and important, consequence of Lemma 6.31.
Theorem 6.33 Let the closed convex function $\phi$ with open domain $\mathcal{D}$ be $\kappa$-self-concordant. If $\mathcal{D}$ does not contain a straight line then the Hessian $\nabla^{2} \phi(x)$ is positive definite at any $x \in \mathcal{D}$.

Proof: Suppose that $H(x)$ is not positive definite for some $x \in \mathcal{D}$. Then there exists a nonzero vector $h \in \mathbb{R}^{n}$ such that $h^{T} H(x) h=0$ or, equivalently, $\|h\|_{x}=0$. For all $\alpha$ such that $x+\alpha h \in \mathcal{D}$ we consider the function

$$
k(\alpha):=h^{T} H(x+\alpha h) h=\nabla^{2} \phi(x+\alpha h)[h, h]=\|h\|_{x+\alpha h}^{2} .
$$

Then $k(0)=0$ and $k(\alpha)$ is continuously differentiable. We claim that $k(\alpha)=0$ for every $\alpha$ in the domain of $k$. Note that $k(\alpha) \geq 0$ for all $\alpha$. Assuming that the claim is not true, we may suppose without loss of generality that $k(\alpha)>0$ on some open interval $(0, \bar{\alpha})$ and, moreover, since $k^{\prime}$ is continuous, that $k(\alpha)$ is nondecreasing on this interval.

The derivative $k^{\prime}(\alpha)$ satisfies

$$
k^{\prime}(\alpha)=\nabla^{3} \phi(x+\alpha h)[h, h, h] \leq 2 \kappa\|h\|_{x+\alpha h}^{3}=2 \kappa k(\alpha)^{\frac{3}{2}} .
$$

This implies $k^{\prime}(0)=0$ and, moreover, if $\kappa=0$ then $k^{\prime}(\alpha)=0$, whence $k(\alpha)=0$ for all $\alpha$ in the domain of $k$. Thus we may further assume that $\kappa>0$. If $\alpha \in(0, \bar{\alpha})$ we may write, using $k(0)=0$ and that $k$ is nondecreasing on $(0, \bar{\alpha})$,

$$
k(\alpha)=\int_{0}^{\alpha} k^{\prime}(\beta) d \beta \leq 2 \kappa \int_{0}^{\alpha} k(\beta)^{\frac{3}{2}} d \beta \leq 2 \kappa \int_{0}^{\alpha} k(\alpha)^{\frac{3}{2}} d \beta=2 \alpha \kappa k(\alpha)^{\frac{3}{2}} .
$$

Dividing at both sides by $k(\alpha)$ we get

$$
1 \leq 2 \alpha \kappa k(\alpha)^{\frac{1}{2}}
$$

which implies

$$
k(\alpha) \geq \frac{1}{4 \alpha^{2}}, \quad \forall \alpha \in(0, \bar{\alpha})
$$

Obviously this contradicts the fact that $k$ is continuous in 0 .
Thus we have shown that $k(\alpha)=0$ for all $\alpha$ such that $x+\alpha h \in \mathcal{D}$. From this we deduce that $\phi(x+\alpha h)$ is linear in $\alpha$, because we have

$$
\begin{aligned}
\phi(x+\alpha h) & =\phi(x)+\alpha h^{T} g(x)+k(\beta), \text { for some } \beta, 0 \leq \beta \leq \alpha \\
& =\phi(x)+\alpha h^{T} g(x) .
\end{aligned}
$$

The hypothesis of the theorem implies that there exists an $\bar{\alpha}$ such that $x+\bar{\alpha} h$ belongs to the boundary of $\mathcal{D}$. Without loss of generality we may assume that $\bar{\alpha}>0$ (else we replace $h$ by $-h)$. It then follows that $\phi(x+\alpha h)$ converges to $\phi(x)+\bar{\alpha} h^{T} g(x)$ if $\alpha$ converges to $\bar{\alpha}$. However, this gives a conflict with Lemma 6.31 which implies that $\phi(x+\alpha h)$ converges if $\alpha$ converges to $\bar{\alpha}$. Thus the proof is compete.

Corollary 6.34 If $\mathcal{D}$ does not contain a straight line then $\phi(x)$ is strictly convex. As a consequence, if $\phi(x)$ has a minimizer then this minimizer is unique.

From now on it will be assumed that the hypothesis of Theorem 6.33 is satisfied. So the domain $\mathcal{D}$ does not contain a straight line. As a consequence we have

$$
\forall x \in \mathcal{D}, \forall h \in \mathbb{R}^{n}:\|h\|_{x}=0 \Leftrightarrow h=0 .
$$

### 6.4.2 Some basic inequalities

Lemma 6.35 Let $x \in \mathcal{D}$ and $\alpha \in \mathbb{R}$ and $0 \neq d \in \mathbb{R}^{n}$ such that $x+\alpha d \in \mathcal{D}$. Then

$$
\frac{\|d\|_{x}}{1+\alpha \kappa\|d\|_{x}} \leq\|d\|_{x+\alpha d} \leq \frac{\|d\|_{x}}{1-\alpha \kappa\|d\|_{x}}
$$

the left inequality holds for all $\alpha$ such that $1+\alpha \kappa\|d\|_{x}>0$ and the right for all $\alpha$ such that $1-\alpha \kappa\|d\|_{x}>0$.

Proof: With $x$ and $d$ fixed, we define, for each $\alpha$ such that $x+\alpha d \in \mathcal{D}$,

$$
q(\alpha):=\|d\|_{x+\alpha d}^{2}=d^{T} H(x+\alpha d) d
$$

Taking the derivative to $\alpha$ we get

$$
q^{\prime}(\alpha)=d^{T}\left(\nabla^{3} \phi(x+\alpha d)[d]\right) d=\nabla^{3} \phi(x+\alpha d)[d, d, d] .
$$

Hence, using the $\kappa$-self-concordancy of $\phi$,

$$
\left|q^{\prime}(\alpha)\right|=\left|\nabla^{3} \phi(x+\alpha d)[d, d, d]\right| \leq 2 \kappa\|d\|_{x+\alpha d}^{3}=2 \kappa q(\alpha)^{\frac{3}{2}} .
$$

This implies that

$$
\left|\frac{d q(\alpha)^{-\frac{1}{2}}}{d \alpha}\right|=\left|\frac{q^{\prime}(\alpha)}{2 q(\alpha)^{\frac{3}{2}}}\right| \leq \kappa .
$$

Consequently, for $0 \leq \alpha \leq 1$ we have

$$
q(0)^{-\frac{1}{2}}-\alpha \kappa \leq q(\alpha)^{-\frac{1}{2}} \leq q(0)^{-\frac{1}{2}}+\alpha \kappa
$$

Using that $q(0)^{\frac{1}{2}}=\|d\|_{x}$ and $q(\alpha)^{\frac{1}{2}}=\|d\|_{x+\alpha d}$ we get

$$
\frac{1}{\|d\|_{x}}-\alpha \kappa \leq \frac{1}{\|d\|_{x+\alpha d}} \leq \frac{1}{\|d\|_{x}}+\alpha \kappa
$$

or, equivalently,

$$
\frac{1-\alpha \kappa\|d\|_{x}}{\|d\|_{x}} \leq \frac{1}{\|d\|_{x+\alpha d}} \leq \frac{1+\alpha \kappa\|d\|_{x}}{\|d\|_{x}} .
$$

Hence, if $1+\alpha \kappa\|d\|_{x}>0$ we obtain

$$
\frac{\|d\|_{x}}{1+\alpha \kappa\|d\|_{x}} \leq\|d\|_{x+\alpha d}
$$

and if $1-\alpha \kappa\|d\|_{x}>0$ we obtain

$$
\|d\|_{x+\alpha d} \leq \frac{\|d\|_{x}}{1-\alpha \kappa\|d\|_{x}}
$$

proving the lemma.

Exercise 6.8 With $h \in \mathbb{R}^{n}$ fixed, define

$$
\beta(\alpha):=\frac{1}{\|h\|_{x+\alpha h}} .
$$

Then

$$
\beta^{\prime}(\alpha):=-\frac{\nabla^{3} \phi(x+\alpha h)[h, h, h]}{2 \nabla^{2} \phi(x+\alpha h)[h, h]^{\frac{3}{2}}},
$$

and hence $\left|\beta^{\prime}(\alpha)\right| \leq \kappa$. Derive Lemma 6.35 from this.
Lemma 6.36 Let $x$ and $d$ be such that $x \in \mathcal{D}, x+d \in \mathcal{D}$ and $\kappa\|d\|_{x}<1$. Then we have, for any nonzero $v \in \mathbb{R}^{n}$,

$$
\begin{equation*}
\left(1-\kappa\|d\|_{x}\right)\|v\|_{x} \leq\|v\|_{x+d} \leq \frac{\|v\|_{x}}{1-\kappa\|d\|_{x}} \tag{6.12}
\end{equation*}
$$

Proof: Fixing $v$, we define for $0 \leq \alpha \leq 1$,

$$
k(\alpha):=v^{T} H(x+\alpha d) v=\|v\|_{x+\alpha d}^{2} .
$$

We then have for the derivative of $k(\alpha)$ to $\alpha$ :

$$
k^{\prime}(\alpha)=v^{T}\left(\nabla^{3} \phi(x+\alpha d)[d]\right) v=\nabla^{3} \phi(x+\alpha d)[d, v, v] .
$$

Using Lemma 6.32 we obtain

$$
\left|k^{\prime}(\alpha)\right|=\left|\nabla^{3} \phi(x+\alpha d)[d, v, v]\right| \leq 2 \kappa\|d\|_{x+\alpha d}\|v\|_{x+\alpha d}^{2}=2 \kappa\|d\|_{x+\alpha d} k(\alpha),
$$

or, also using Lemma 6.35,

$$
\left|\frac{k^{\prime}(\alpha)}{k(\alpha)}\right| \leq 2 \kappa\|d\|_{x+\alpha d} \leq \frac{2 \kappa\|d\|_{x}}{1-\alpha \kappa\|d\|_{x}} .
$$

Note that $k(\alpha)>0$. Since $k^{\prime}(\alpha) / k(\alpha)$ is the derivative to $\alpha$ of $\log k(\alpha)$ we find that

$$
\left|\frac{d \log k(\alpha)}{d \alpha}\right| \leq \frac{2 \kappa\|d\|_{x}}{1-\alpha \kappa\|d\|_{x}} .
$$

Using this we may write

$$
\begin{aligned}
\log \frac{\|v\|_{x+d}}{\|v\|_{x}} & =\frac{1}{2} \log \frac{k(1)}{k(0)}=\frac{1}{2}(\log k(1)-\log k(0)) \\
& =\frac{1}{2} \int_{0}^{1}\left(\frac{d \log k(\alpha)}{d \alpha}\right) d \alpha \leq \int_{0}^{1} \frac{\kappa\|d\|_{x}}{1-\alpha \kappa\|d\|_{x}} d \alpha \\
& =-\left.\log \left(1-\alpha \kappa\|d\|_{x}\right)\right|_{\alpha=0} ^{1}=\log \left(\frac{1}{1-\kappa\|d\|_{x}}\right)
\end{aligned}
$$

and, similarly,

$$
\log \frac{\|v\|_{x+d}}{\|v\|_{x}}=\frac{1}{2} \int_{0}^{1}\left(\frac{d \log k(\alpha)}{d \alpha}\right) d \alpha \geq-\int_{0}^{1} \frac{\kappa\|d\|_{x}}{1-\alpha \kappa\|d\|_{x}} d \alpha=\log \left(1-\kappa\|d\|_{x}\right) .
$$

As an immediate consequence we obtain

$$
1-\kappa\|d\|_{x} \leq \frac{\|v\|_{x+d}}{\|v\|_{x}} \leq \frac{1}{1-\kappa\|d\|_{x}} .
$$

which implies (6.12). Thus the lemma is proved.

Exercise 6.9 If $x$ and $d$ are such that $x \in \mathcal{D}$ and $x+d \in \mathcal{D}$, then

$$
\left(1-\kappa\|d\|_{x}\right)^{2} H(x) \preceq H(x+d) \preceq \frac{H(x)}{\left(1-\kappa\|d\|_{x}\right)^{2}},
$$

Derive this from Lemma 6.36.

Lemma 6.37 Let $x \in \mathcal{D}$ and $d \in \mathbb{R}^{m}$. If $\|d\|_{x}<\frac{1}{\kappa}$ then $x+d \in \mathcal{D}$.
Proof: Since $\|d\|_{x}<\frac{1}{\kappa}$, we have from Lemma 6.36 that $H(x+\alpha d)$ is bounded for all $0 \leq \alpha \leq 1$, and thus $\phi(x+\alpha d)$ is bounded. On the other hand, $\phi$ takes infinite values on the boundary of the feasible set, by Lemma 6.31. Consequently, $x+d \in \mathcal{D}$.

### 6.4.3 Linear convergence of the damped Newton method

In this section we consider the case where $x \in \mathcal{D}$ lies outside the region where the Newton process is quadratically convergent. Performing a damped Newton step, with damping factor $\alpha$, the new iterate is given by

$$
x^{+}=x+\alpha \Delta x .
$$

The next lemma shows that with an appropriate choice of $\alpha$ we can guarantee a fixed decrease in $\phi$ after the step.

Lemma 6.38 Let $x \in \mathcal{D}$ and $\delta:=\delta(x)$. If $\alpha:=\frac{1}{1+\kappa \delta}$ then

$$
\phi(x)-\phi(x+\alpha \Delta x) \geq \frac{\psi(\kappa \delta)}{\kappa^{2}}
$$

Proof: Define

$$
\Delta(\alpha):=\phi(x)-\phi(x+\alpha \Delta x)
$$

Then

$$
\begin{aligned}
\Delta^{\prime}(\alpha) & =-g(x+\alpha \Delta x)^{T} \Delta x \\
\Delta^{\prime \prime}(\alpha) & =-\Delta x^{T} H(x+\alpha \Delta x) \Delta x=-\nabla^{2} \phi(x+\alpha \Delta x)[\Delta x, \Delta x] \\
\Delta^{\prime \prime \prime}(\alpha) & =-\nabla^{3} \phi(x+\alpha \Delta x)[\Delta x, \Delta x, \Delta x]
\end{aligned}
$$

Now using that $\phi$ is $\kappa$-self-concordant we deduce from the last expression that

$$
\Delta^{\prime \prime \prime}(\alpha) \geq-2 \kappa\|\Delta x\|_{x+\alpha \Delta x}^{3}
$$

Hence, also using Lemma 6.35,

$$
\Delta^{\prime \prime \prime}(\alpha) \geq-2 \kappa \frac{\|\Delta x\|_{x}^{3}}{\left(1-\alpha \kappa\|\Delta x\|_{x}\right)^{3}}=\frac{-2 \kappa \delta^{3}}{(1-\alpha \kappa \delta)^{3}}
$$

As a consequence we have

$$
\Delta^{\prime \prime}(\alpha)-\Delta^{\prime \prime}(0) \geq \int_{0}^{\alpha} \frac{-2 \kappa \delta^{3}}{(1-\beta \kappa \delta)^{3}} d \beta=\left.\frac{-\delta^{2}}{(1-\beta \kappa \delta)^{2}}\right|_{\beta=0} ^{\alpha}=\frac{-\delta^{2}}{(1-\alpha \kappa \delta)^{2}}+\delta^{2}
$$

Since $\Delta^{\prime \prime}(0)=-\nabla^{2} \phi(x)[\Delta x, \Delta x]=-\delta^{2}$, we obtain

$$
\Delta^{\prime \prime}(\alpha) \geq \frac{-\delta^{2}}{(1-\alpha \kappa \delta)^{2}}
$$

In a similar way, by integrating, we derive an estimate for $\Delta^{\prime}(\alpha)$ :

$$
\Delta^{\prime}(\alpha)-\Delta^{\prime}(0) \geq \int_{0}^{\alpha} \frac{-\delta^{2}}{(1-\beta \kappa \delta)^{2}} d \beta=\left.\frac{-\delta}{\kappa(1-\beta \kappa \delta)}\right|_{\beta=0} ^{\alpha}=\frac{-\delta}{\kappa(1-\alpha \kappa \delta)}+\frac{\delta}{\kappa} .
$$

Since $\Delta^{\prime}(0)=-g(x)^{T} \Delta x=\Delta x H(x) \Delta x=\delta^{2}$, we obtain

$$
\Delta^{\prime}(\alpha) \geq \frac{-\delta}{\kappa(1-\alpha \kappa \delta)}+\frac{\delta}{\kappa}+\delta^{2} .
$$

Finally, in the same way we derive an estimate for $\Delta(\alpha)$. Using that $\Delta(0)=0$ we have

$$
\Delta(\alpha) \geq \int_{0}^{\alpha}\left(\frac{-\delta}{\kappa(1-\beta \kappa \delta)}+\frac{\delta}{\kappa}+\delta^{2}\right) d \beta=\frac{1}{\kappa^{2}}\left(\log (1-\alpha \kappa \delta)+\alpha \kappa \delta+\alpha \kappa^{2} \delta^{2}\right) .
$$

The last expression is maximal for $\bar{a}=\frac{1}{1+\kappa \delta}$. Substitution of this value yields

$$
\Delta(\bar{a}) \geq \frac{1}{\kappa^{2}}\left(\log \left(1-\frac{\kappa \delta}{1+\alpha \delta}\right)+\kappa \delta\right)=\frac{1}{\kappa^{2}}(\kappa \delta-\log (1+\alpha \delta))=\frac{1}{\kappa^{2}} \psi(\kappa \delta)
$$

which is the desired inequality.

### 6.4.4 Quadratic convergence of Newton's method

Let $x^{+}:=x+\Delta x$ denote the iterate after the Newton step at $x$. Recall that the Newton step at $x$ is given by

$$
\Delta x=-H(x)^{-1} g(x)
$$

where $H(x)$ is defined as above and $g(x)=\nabla \phi(x)$. We will use (cf. Exercise 6.1)

$$
\delta(x):=\|\Delta x\|_{x}=\sqrt{g(x)^{T} H(x)^{-1} g(x)}
$$

to measure the 'length' of the Newton step. Note that if $x$ is such that $\phi(x)$ is minimal then $g(x)=0$ and hence $\delta(x)=0$; whereas in all other cases $\delta(x)$ will be positive.

After the Newton step we have

$$
\begin{equation*}
\delta\left(x^{+}\right)=\sqrt{g\left(x^{+}\right)^{T} H\left(x^{+}\right)^{-1} g\left(x^{+}\right)}=\left\|H\left(x^{+}\right)^{-1} g\left(x^{+}\right)\right\|_{x^{+}} . \tag{6.13}
\end{equation*}
$$

Lemma 6.39 If $\delta(x) \leq \frac{1}{3 \kappa}$ then $x^{+}$is feasible and

$$
\delta\left(x^{+}\right) \leq \kappa\left(\frac{\delta(x)}{1-\kappa \delta(x)}\right)^{2} \leq \frac{9 \kappa}{4} \delta(x)^{2}
$$

Proof: The feasibility of $x^{+}$follows from Lemma 6.37. Using the Mean Value Theorem we have for some $\beta, 0 \leq \beta \leq 1$,

$$
g\left(x^{+}\right)=g(x)+H(x) \Delta x+\frac{1}{2} \nabla^{3}(x+\beta \Delta x)[\Delta x, \Delta x] .
$$

Since, by the definition of $\Delta x, g(x)+H(x) \Delta x=0$ we obtain

$$
g\left(x^{+}\right)=\frac{1}{2} \nabla^{3} \phi(x+\beta \Delta x)[\Delta x, \Delta x] .
$$

Hence, for any vector $p \in \mathbb{R}^{n}$ we have

$$
p^{T} g\left(x^{+}\right)=\frac{1}{2} \nabla^{3} \phi(x+\beta \Delta x)[\Delta x, \Delta x, p] .
$$

Using Lemma 6.32 we get

$$
\begin{equation*}
\left|p^{T} g\left(x^{+}\right)\right| \leq \kappa\|\Delta x\|_{x+\beta \Delta x}^{2}\|p\|_{x+\beta \Delta x} . \tag{6.14}
\end{equation*}
$$

By Lemma 6.35 we have, using $\delta(x)=\|x\|_{x}$,

$$
\|\Delta x\|_{x+\beta \Delta x} \leq \frac{\|\Delta x\|_{x}}{1-\beta \kappa\|\Delta x\|_{x}}=\frac{\delta(x)}{1-\beta \kappa \delta(x)}
$$

and

$$
\|p\|_{x+\beta \Delta x} \leq \frac{\|p\|_{x+\Delta x}}{1-(1-\beta) \kappa\|\Delta x\|_{x+\Delta x}} \leq \frac{\|p\|_{x+\Delta x}}{1-\frac{(1-\beta) \kappa\|\Delta x\|_{x}}{1-\kappa\|\Delta x\|_{x}}}=\frac{\|p\|_{x+\Delta x}}{1-\frac{(1-\beta) \kappa \delta(x)}{1-\kappa \delta(x)}} .
$$

Substituting the last two inequalities in (6.14), while replacing $p$ by the Newton step $H\left(x^{+}\right)^{-1} g\left(x^{+}\right)$at $x^{+}$and also using $\|p\|_{x+\Delta x}=\delta\left(x^{+}\right)$, from (6.13), we obtain

$$
\delta\left(x^{+}\right)^{2}=\left|g\left(x^{+}\right)^{T} H\left(x^{+}\right)^{-1} g\left(x^{+}\right)\right| \leq \kappa\left(\frac{\delta(x)}{1-\beta \kappa \delta(x)}\right)^{2} \frac{\delta\left(x^{+}\right)}{1-\frac{(1-\beta) \kappa \delta(x)}{1-\kappa \delta(x)}} .
$$

Dividing at the left and right by $\delta\left(x^{+}\right)$we get

$$
\begin{equation*}
\delta\left(x^{+}\right) \leq \frac{\kappa \delta(x)^{2}}{h(\beta)} \tag{6.15}
\end{equation*}
$$

for some $\beta, 0 \leq \beta \leq 1$, where

$$
h(\beta)=(1-\beta \kappa \delta(x))^{2}\left(1-\frac{(1-\beta) \kappa \delta(x)}{1-\kappa \delta(x)}\right) .
$$

By elementary means it can be checked that

$$
h^{\prime \prime}(\beta)=\frac{2 \kappa^{2} \delta(x)^{2}(3 \beta \kappa \delta(x)-2 \kappa \delta(x)-1)}{1-\kappa \delta(x)} \leq \frac{2 \kappa^{2} \delta(x)^{2}(\kappa \delta(x)-1)}{1-\kappa \delta(x)}=-2 \kappa^{2} \delta(x)^{2}<0,
$$

whence $h(\beta)$ is concave. Hence, for $\beta, 0 \leq \beta \leq 1$,

$$
\begin{equation*}
h(\beta) \geq \min \{h(0), h(1)\}=\min \left\{\frac{1-2 \kappa \delta(x)}{1-\kappa \delta(x)},\left(1-\kappa \delta(x)^{2}\right)\right\}=(1-\kappa \delta(x))^{2}, \tag{6.16}
\end{equation*}
$$

where we used that

$$
(1-\kappa \delta(x))^{2} \leq \frac{1-2 \kappa \delta(x)}{1-\kappa \delta(x)}
$$

whenever

$$
\kappa \delta(x) \leq \frac{3-\sqrt{5}}{2}=0.381966
$$

By the hypothesis of the lemma we have $\kappa \delta(x) \leq \frac{1}{3}$ and hence we may conclude from (6.15) and (6.16) that

$$
\delta\left(x^{+}\right) \leq \frac{\kappa \delta(x)^{2}}{(1-\kappa \delta(x))^{2}}
$$

which proves the lemma.

Remark: Lemma 6.39 is also valid if $\delta(x)<\frac{1}{\kappa}$. This can be shown as follows. For $v \in \mathbb{R}^{n}$ and $0 \leq \alpha \leq 1$, define

$$
k(\alpha):=v^{T} g(x+\alpha \Delta x)-(1-\alpha) v^{T} g(x) .
$$

Note that if $v=H\left(x^{+}\right)^{-1} g\left(x^{+}\right)$then

$$
k(1)=g\left(x^{+}\right)^{T} H\left(x^{+}\right)^{-1} g\left(x^{+}\right)=\delta\left(x^{+}\right)^{2} .
$$

Thus our aim is to find a good estimate for $k(1)$. Taking the derivative of $k$ to $\alpha$ we get, also using that $H(x) \Delta x=-g(x)$,

$$
\begin{aligned}
k^{\prime}(\alpha) & =v^{T} H(x+\alpha \Delta x) \Delta x+v^{T} g(x) \\
& =v^{T} H(x+\alpha \Delta x) \Delta x-v^{T} H(x) \Delta x \\
& =v^{T}(H(x+\alpha \Delta x)-H(x)) \Delta x
\end{aligned}
$$

By Exercise 6.9,

$$
H(x+\alpha \Delta x)-H(x) \preceq\left(\frac{1}{\left(1-\alpha \kappa\|\Delta x\|_{x}\right)^{2}}-1\right) H(x) .
$$

Now applying the generalized Cauchy inequality of Lemma A. 1 in the Appendix we get

$$
v^{T}(H(x+\alpha \Delta x)-H(x)) \Delta x \leq\left(\frac{1}{\left(1-\alpha \kappa\|\Delta x\|_{x}\right)^{2}}-1\right)\|v\|_{x}\|\Delta x\|_{x}
$$

Hence, since $\|\Delta x\|_{x}=\delta(x)$,

$$
k^{\prime}(\alpha) \leq\left(\frac{1}{(1-\alpha \kappa \delta(x))^{2}}-1\right)\|v\|_{x} \delta(x)
$$

Therefore, since $k(0)=0$,

$$
k(1) \leq \delta(x)\|v\|_{x} \int_{0}^{1}\left(\frac{1}{(1-\alpha \kappa \delta(x))^{2}}-1\right) d \alpha
$$

We have

$$
\begin{aligned}
\int_{0}^{1}\left(\frac{1}{(1-\alpha \kappa \delta(x))^{2}}-1\right) d \alpha & =\left.\left(\frac{1}{\kappa \delta(x)(1-\alpha \kappa \delta(x))}-\alpha\right)\right|_{\alpha=0} ^{1} \\
& =\frac{1}{\kappa \delta(x)(1-\kappa \delta(x))}-1-\frac{1}{\kappa \delta(x)} \\
& =\frac{\kappa \delta(x)}{1-\kappa \delta(x)}
\end{aligned}
$$

Substitution gives

$$
k(1) \leq\|v\|_{x} \frac{\kappa \delta(x)^{2}}{1-\kappa \delta(x)}
$$

For $v=H\left(x^{+}\right)^{-1} g\left(x^{+}\right)$, we have, by Lemma 6.36,

$$
\|v\|_{x} \leq \frac{\|v\|_{x^{+}}}{1-\alpha \kappa\|\Delta x\|_{x}}=\frac{\delta\left(x^{+}\right)}{1-\alpha \kappa \delta(x)} .
$$

Since $k(1)=\delta\left(x^{+}\right)^{2}$, it follows by substitution,

$$
\delta\left(x^{+}\right)^{2}=k(1) \leq \frac{\delta\left(x^{+}\right)}{1-\alpha \kappa \delta(x)} \frac{\kappa \delta(x)^{2}}{1-\kappa \delta(x)}
$$

Dividing both sides by $\delta\left(x^{+}\right)$the claim follows.

### 6.4.5 Existence and properties of minimizer

In this section we derive a necessary and sufficient condition for the existence of a minimizer of $\phi$ and we show that, if it exists, the minimizer is unique. For the case where the minimizer exists, denoted by $x^{*}$, we also derive some estimates for $\phi(x)-\phi\left(x^{*}\right)$ and $\left\|x-x^{*}\right\|_{x}$. We start with two lemmas.

Lemma 6.40 Let $x \in \mathcal{D}$ and $0 \neq h \in \mathbb{R}^{n}$ such that $x+h \in \mathcal{D}$. Then

$$
\begin{align*}
h^{T}(g(x+h)-g(x)) & \geq \frac{\|h\|_{x}^{2}}{1+\kappa\|h\|_{x}}  \tag{6.17}\\
\phi(x+h)-\phi(x) & \geq h^{T} g(x)+\frac{\psi\left(\kappa\|h\|_{x}\right)}{\kappa^{2}} \tag{6.18}
\end{align*}
$$

Proof: Using Lemma 6.35 we write

$$
\begin{aligned}
h^{T}(g(x+h)-g(x)) & =\int_{0}^{1} h^{T} H(x+\alpha h) h d \alpha=\int_{0}^{1}\|h\|_{x+\alpha h}^{2} d \alpha \\
& \geq \int_{0}^{1} \frac{\|h\|_{x}^{2}}{\left(1+\alpha \kappa\|h\|_{x}\right)^{2}} d \alpha=\frac{\|h\|_{x}^{2}}{1+\kappa\|h\|_{x}}
\end{aligned}
$$

which proves the first inequality. Using this inequality write

$$
\begin{aligned}
\phi(x+h)-\phi(x)-h^{T} g(x) & =\int_{0}^{1} h^{T}(g(x+\alpha h)-g(x)) d \alpha \\
& \geq \int_{0}^{1} \frac{\alpha\|h\|_{x}^{2}}{1+\alpha \kappa\|h\|_{x}} d \alpha \\
& =\frac{\kappa\|h\|_{x}-\log \left(1+\kappa\|h\|_{x}\right)}{\kappa^{2}}=\frac{\psi\left(\kappa\|h\|_{x}\right)}{\kappa^{2}} .
\end{aligned}
$$

This completes the proof.

Lemma 6.41 Let $x \in \mathcal{D}$ and $0 \neq h \in \mathbb{R}^{n}$ such that $x+h \in \mathcal{D}$ and let $\|h\|_{x}<1$. Then

$$
\begin{aligned}
h^{T}(g(x+h)-g(x)) & \leq \frac{\|h\|_{x}^{2}}{1-\kappa\|h\|_{x}} \\
\phi(x+h)-\phi(x) & \leq h^{T} g(x)+\frac{\psi\left(-\kappa\|h\|_{x}\right)}{\kappa^{2}}
\end{aligned}
$$

Proof: As in the previous proof we use Lemma 6.35 and write

$$
\begin{aligned}
h^{T}(g(x+h)-g(x)) & =\int_{0}^{1} h^{T} H(x+\alpha h) h d \alpha=\int_{0}^{1}\|h\|_{x+\alpha h}^{2} d \alpha \\
& \leq \int_{0}^{1} \frac{\|h\|_{x}^{2}}{\left(1-\alpha \kappa\|h\|_{x}\right)^{2}} d \alpha=\frac{\|h\|_{x}^{2}}{1-\kappa\|h\|_{x}}
\end{aligned}
$$

which is the first inequality. Using this inequality write

$$
\begin{aligned}
\phi(x+h)-\phi(x)-h^{T} g(x) & =\int_{0}^{1} h^{T}(g(x+\alpha h)-g(x)) d \alpha \\
& \leq \int_{0}^{1} \frac{\alpha\|h\|_{x}^{2}}{1-\alpha \kappa\|h\|_{x}} d \alpha \\
& =\frac{-\kappa\|h\|_{x}-\log \left(1-\kappa\|h\|_{x}\right)}{\kappa^{2}}=\frac{\psi\left(-\kappa\|h\|_{x}\right)}{\kappa^{2}}
\end{aligned}
$$

completing the proof.

As usual, for each $x \in \mathcal{D}, \delta(x)=\|\Delta x\|_{x}$, with $\Delta x$ denoting the Newton step at $x$. We now prove that if $\delta(x)<\frac{1}{\kappa}$ for some $x \in \mathcal{D}$ then $\phi$ must have a minimizer. Note that this surprising result expresses that some local condition on $\phi$ provides us with a global property, namely the existence of a minimizer.

Theorem 6.42 Let $\delta(x)<\frac{1}{\kappa}$ for some $x \in \mathcal{D}$. Then $\phi$ has a unique minimizer $x^{*}$ in D.

Proof: The proof is based on the observation that the level set

$$
\begin{equation*}
\{y \in \mathcal{D}: \phi(y) \leq \phi(x)\} \tag{6.19}
\end{equation*}
$$

with $x$ as given in the theorem, is compact. This can be seen as follows. Let $y \in \mathcal{D}$. Writing $y=x+h$, Lemma 6.40 implies the inequality

$$
\begin{equation*}
\phi(y)-\phi(x) \geq h^{T} g(x)+\frac{\psi\left(\kappa\|h\|_{x}\right)}{\kappa^{2}}=-h^{T} H(x) \Delta x+\frac{\psi\left(\kappa\|h\|_{x}\right)}{\kappa^{2}} \tag{6.20}
\end{equation*}
$$

where we used that, by definition, the Newton step $\Delta x$ at $x$ satisfies $H(x) \Delta x=-g(x)$. Since

$$
h^{T} H(x) \Delta x \leq\|h\|_{x}\|\Delta x\|_{x}=\|h\|_{x} \delta(x)
$$

we thus have

$$
\phi(y)-\phi(x) \geq-\|h\|_{x} \delta(x)+\frac{\psi\left(\kappa\|h\|_{x}\right)}{\kappa^{2}}
$$

Hence, if $\phi(y) \leq \phi(x)$, then

$$
\begin{equation*}
\frac{\psi\left(\kappa\|h\|_{x}\right)}{\kappa\|h\|_{x}} \leq \kappa \delta(x)<1 \tag{6.21}
\end{equation*}
$$

Putting $\xi:=\kappa\|h\|_{x}$ one may easily verify that $\psi(\xi) / \xi$ is monotonically increasing for $\xi>0$ and goes to 1 if $\xi \rightarrow \infty$. Therefore, since $\delta(x)<\frac{1}{\kappa}$, we may conclude from (6.21) that $\kappa\|h\|_{x}$ is bounded above. This implies that the level set (6.19) is bounded. From this we conclude that $\phi$ has a minimizer $x^{*}$. Finally, Corollary 6.34 implies that this minimizer is unique.

The next example shows that the above result is sharp.
Example 6.43 With $\epsilon \geq 0$ fixed, consider the function $f_{\epsilon}:(0, \infty) \rightarrow \mathbb{R}$ defined by

$$
f_{\epsilon}(x)=\epsilon x-\log x, x>0 .
$$

This function is 1 -self-concordant. One has

$$
f_{\epsilon}^{\prime}(x)=\epsilon-\frac{1}{x}, \quad f_{\epsilon}^{\prime \prime}(x)=\frac{1}{x^{2}} .
$$

Therefore,

$$
\delta(x)=\sqrt{x^{2}\left(\epsilon-\frac{1}{x}\right)^{2}}=|1-\epsilon x|
$$

Thus, for $\epsilon=0$ we have $\delta(x)=1$ for each $x>0$. Since $f_{0}(x)=-\log x, f_{0}(x)$ has no minimizer. On the other hand, if $\epsilon>0$ then $\delta\left(\frac{1}{\epsilon}\right)=0<1$ and $x=\frac{1}{\epsilon}$ is a minimizer.

Exercise 6.10 Let $\delta(x) \geq \frac{1}{\kappa}$ for all $x \in \mathcal{D}$. Then $\phi$ is unbounded and, hence, has no minimizer in $\mathcal{D}$. Prove this. (Hint: use Lemma 6.38.)

Exercise 6.11 Let $\delta(x) \geq \frac{1}{\kappa}$ for all $x \in \mathcal{D}$. Then $\mathcal{D}$ is unbounded. Prove this.

The proof of the next theorem requires the result of the following exercise.
Exercise 6.12 For $s<1$ one has ${ }^{2}$

$$
\psi(-s)=\sup _{t>-1}\{s t-\psi(t)\}
$$

whence

$$
\begin{equation*}
\psi(-s)+\psi(t) \geq s t, \quad s<1, t>-1 . \tag{6.22}
\end{equation*}
$$

Prove this.

Theorem 6.44 Let $x \in \mathcal{D}$ be such that $\delta(x)<\frac{1}{\kappa}$ and let $x^{*}$ denote the unique minimizer of $\phi$. Then, with $\delta:=\delta(x)$,

$$
\left.\begin{array}{rl}
\frac{\psi(\kappa \delta)}{\kappa^{2}} & \leq \phi(x)-\phi\left(x^{*}\right)
\end{array}\right) \frac{\psi(-\kappa \delta)}{\kappa^{2}}, ~=\left\|x-x^{*}\right\|_{x} \leq \frac{\delta}{1-\kappa \delta}=-\frac{\psi^{\prime}(-\kappa \delta)}{\kappa} .
$$

Proof: The left inequality in (6.23) follows from Lemma 6.38, because $\phi$ is minimal at $x^{*}$. Furthermore, from (6.18) in Lemma 6.40, with $h=x^{*}-x$, we get the right inequality in (6.23):

$$
\begin{aligned}
\phi\left(x^{*}\right)-\phi(x) & \geq(h)^{T} g(x)+\frac{\psi\left(\kappa\|h\|_{x}\right)}{\kappa^{2}} \\
& \geq-\|h\|_{x} \delta(x)+\frac{\psi\left(\kappa\|h\|_{x}\right)}{\kappa^{2}} \\
& \geq \frac{1}{\kappa^{2}}\left(-\kappa\|h\|_{x} \kappa \delta+\psi\left(\kappa\|h\|_{x}\right)\right) \\
& \geq-\frac{\psi(-\kappa \delta)}{\kappa^{2}}
\end{aligned}
$$

where the second inequality holds since

$$
\begin{equation*}
(h)^{T} g(x)=-(h)^{T} H(x) \Delta x \leq\|h\|_{x}\|\Delta x\|_{x}=\|h\|_{x} \delta(x)=\|h\|_{x} \delta \tag{6.25}
\end{equation*}
$$

and the fourth inequality follows from (6.22).
For the proof of (6.24) we first derive from (6.25) and (6.17) in Lemma 6.40 that

$$
\frac{\|h\|_{x}^{2}}{1+\kappa\|h\|_{x}} \leq(h)^{T} g(x) \leq\|h\|_{x} \delta .
$$

Dividing by $\|h\|_{x}$ we get

$$
\frac{\|h\|_{x}}{1+\kappa\|h\|_{x}} \leq \delta
$$

[^5]which is equivalent to
$$
\|h\|_{x} \leq \frac{\delta}{1-\kappa \delta}=-\frac{\psi^{\prime}(-\kappa \delta)}{\kappa}
$$
which proves the right inequality in (6.24).
Note that the left inequality in (6.24) is trivial if $\kappa\|h\|_{x}>1$ since $\frac{\delta}{1+\kappa \delta}<\frac{1}{\delta}$. Thus we may assume that $1-\kappa\|h\|_{x}>0$. For $0 \leq \alpha \leq 1$, consider
$$
k(\alpha):=g\left(x^{*}-\alpha h\right)^{T} H(x)^{-1} g(x) .
$$

One has $k(0)=0$ and $k(1)=\delta(x)^{2}=\delta^{2}$. Using the result in Exercise 6.9 and the generalized Cauchy inequality of Lemma A. 1 in the Appendix we may write

$$
k^{\prime}(\alpha)=-h^{T} H\left(x^{*}-\alpha h\right) H(x)^{-1} g(x) \leq \frac{\|h\|_{x} \delta(x)}{\left(1-\kappa\|h\|_{x}\right)^{2}} .
$$

Hence we have

$$
\delta^{2}=k(1) \leq \int_{0}^{1} \frac{\|h\|_{x} \delta}{\left(1-\kappa\|h\|_{x}\right)^{2}} d \alpha=\frac{\|h\|_{x} \delta}{1-\kappa\|h\|_{x}}
$$

Dividing both sides by $\delta$ we obtain

$$
\delta \leq \frac{\|h\|_{x}}{1-\kappa\|h\|_{x}}
$$

which is equivalent to

$$
\|h\|_{x} \geq \frac{\delta}{1+\kappa \delta}
$$

Thus the proof is complete.

### 6.4.6 Solution strategy

Now we have all the ingredients to design an efficient method for solving the problem

$$
\min \{\phi(x): x \in \mathcal{D}\}
$$

Assuming that we have given $x^{0} \in \mathcal{D}$ we calculate $\delta^{0}:=\delta\left(x^{0}\right)$. We distinguish between the cases where $\delta^{0} \leq \frac{1}{3 \kappa}$ and $\delta^{0}>\frac{1}{3 \kappa}$ respectively.
If $\delta^{0} \leq \frac{1}{3 \kappa}$, starting at $x^{0}$ we repeatedly apply full Newton steps until the iterate $x^{k}$ satisfies $\delta\left(x^{k}\right) \leq \epsilon$, where $\epsilon>0$ is some prescribed accuracy parameter. We can estimate the required number of Newton steps by using Lemma 6.39. For $k \geq 1$ this lemma gives

$$
\delta\left(x^{k}\right) \leq \frac{9 \kappa}{4} \delta\left(x^{k-1}\right)^{2} \leq \cdots \leq \frac{4}{9 \kappa}\left(\frac{9 \kappa \delta^{0}}{4}\right)^{2^{k}}
$$

Hence we will have $\delta\left(x^{k}\right) \leq \epsilon$ if

$$
\frac{4}{9 \kappa}\left(\frac{9 \kappa \delta^{0}}{4}\right)^{2^{k}} \leq \epsilon
$$

Taking the logarithm at both sides this reduces to

$$
2^{k} \log \frac{9 \kappa \delta^{0}}{4} \leq \log \frac{9 \kappa \epsilon}{4}
$$

Note that $\frac{9 \kappa \delta^{0}}{4} \leq \frac{3}{4}<1$. Dividing by $-\log \frac{9 \kappa \delta^{0}}{4}$ we get

$$
2^{k} \geq \frac{\log \frac{9 \kappa \epsilon}{4}}{\log \frac{9 \kappa \delta^{0}}{4}},
$$

or, equivalently,

$$
k \geq{ }^{2} \log \frac{\log \frac{9 \kappa \epsilon}{4}}{\log \frac{9 \kappa \delta^{0}}{4}} .
$$

Since $-\log \frac{9 \kappa \delta^{0}}{4} \geq-\log \frac{3}{4}$ we find that after no more than

$$
{ }^{2} \log \left(\frac{\log \frac{9 \kappa \epsilon}{4}}{\log \frac{3}{4}}\right)={ }^{2} \log \left(\frac{\log \frac{9}{4}+\log \kappa \epsilon}{\log \frac{3}{4}}\right)={ }^{2} \log (-2.8188-3.4761 \log \kappa \epsilon)
$$

steps the process will stop and the output will be an $x \in \mathcal{D}$ such that $\left\|x-x^{*}\right\|_{x} \leq \epsilon$.
If $\delta^{0}>\frac{1}{3 \kappa}$ then we use damped Newton steps. By Lemma 6.38 each damped Newton step decreases $\phi$ with at least the value

$$
\frac{\psi(\kappa \delta)}{\kappa^{2}} \geq \frac{\psi\left(\frac{1}{3}\right)}{\kappa^{2}}=\frac{0.0457}{\kappa^{2}}>\frac{1}{22 \kappa^{2}}
$$

Hence, after no more than

$$
22 \kappa^{2}\left(\phi\left(x^{0}\right)-\phi\left(x^{*}\right)\right)
$$

we reach the region where $\delta(x)<\frac{1}{3 \kappa}$. Then we can proceed with full Newton steps, and after a total of

$$
\left\lceil 22 \kappa^{2}\left(\phi\left(x^{0}\right)-\phi\left(x^{*}\right)\right)+{ }^{2} \log (-2.8188-3.4761 \log \kappa \epsilon)\right\rceil
$$

steps we obtain an $x \in \mathcal{D}$ such that $\left\|x-x^{*}\right\|_{x} \leq \epsilon$. One may easily check, by using Maple or Mathematica, that for $\kappa \epsilon<10^{-6}$ this bound is less than

$$
\left\lceil 22 \kappa^{2}\left(\phi\left(x^{0}\right)-\phi\left(x^{*}\right)\right)+6(1-\kappa \epsilon)\right\rceil .
$$

Note the drawback of the above iteration bound: usually we have no a priori knowledge of $\phi\left(x^{*}\right)$ and the bound cannot be calculated at the start of the algorithm. But in many cases we can derive a good estimate for $\phi\left(x^{0}\right)-\phi\left(x^{*}\right)$ and we obtain an upper bound for the number of iterations at the start of the algorithm.

Example 6.45 Consider the function $f:(-1, \infty) \rightarrow \mathbb{R}$ defined by

$$
f(x)=x-\log (1+x), x>0
$$

We established earlier that $f$ is 1-self-concordant, in Example 6.15. One has

$$
f^{\prime}(x)=\frac{x}{1+x}, \quad f^{\prime \prime}(x)=\frac{1}{(1+x)^{2}}
$$

Therefore,

$$
\delta(x)=\sqrt{\frac{f^{\prime}(x)^{2}}{f^{\prime \prime}(x)}}=\sqrt{x^{2}}=|x|
$$

Note that $x=0$ is the unique minimizer. The Newton step at $x$ is given by

$$
\Delta x=-\frac{f^{\prime}(x)}{f^{\prime \prime}(x)}=-x(1+x)
$$

and a full Newton step yields

$$
x^{+}=x-x(1+x)=-x^{2}
$$

The Newton step is feasible only if $-x^{2}>-1$, i.e. only if $\delta(x)<1$. Note that the theory guarantees feasibility in that case. Moreover, if the Newton step is feasible then $\delta\left(x^{+}\right)=\delta(x)^{2}$, which is better than the theoretical result of Lemma 6.18. When we take a damped Newton step, with the default step size $\alpha=\frac{1}{1+\delta(x)}$, the next iterate is given by

$$
x^{+}=x-\frac{x(1+x)}{1+|x|}= \begin{cases}0, & \text { if } x>0 \\ \frac{-2 x^{2}}{1-x} & \text { if } x<0\end{cases}
$$

Thus we find in this example that the damped Newton step is exact if $x>0$. Also, if $-1<x<0$ then

$$
\frac{-2 x^{2}}{1-x}<-x^{2}
$$

and hence then the full Newton step performs better than the damped Newton step. Finally observe that if we apply Newton's method until $\delta(x) \leq \epsilon$ then the output is an $x$ such that $|x| \leq \epsilon . \quad *$

Example 6.46 Consider

$$
\phi(x):=-\sum_{i=1}^{n} \log x_{i}
$$

with $0<x \in \mathbb{R}^{n}$. We established in Example 6.14 that $\phi$ is 1 -self-concordant, and the first and second order derivatives are given by

$$
g(x)=\nabla \phi(x)=\frac{-e}{x}, \quad H(x)=\nabla^{2} \phi(x)=\operatorname{diag}\left(\frac{e}{x^{2}}\right)
$$

Therefore,

$$
\delta(x)=\sqrt{g(x)^{T} H(x)^{-1} g(x)}=\sqrt{\sum_{i=1}^{n} 1}=\|e\|=\sqrt{n}
$$

We conclude from this that $\phi$ has no minimizer (cf. Exercise 6.10).

Example 6.47 We now consider the function $\Psi$ introduced in Example 6.16:

$$
\Psi(x):=\sum_{i=1}^{n} \psi\left(x_{i}\right)=\sum_{i=1}^{n} x_{i}-\log \left(1+x_{i}\right)
$$

with $-e<x \in \mathbb{R}^{n}$. The gradient and Hessian of $\Psi$ are

$$
g(x)=\nabla \phi(x)=\frac{x}{e+x}, H(x)=\nabla^{2} \phi(x)=\operatorname{diag}\left(\frac{e}{(e+x)^{2}}\right) .
$$

We established that $\Psi$ is 1-self-concordant. One has

$$
\delta(x)=\sqrt{g(x)^{T} H(x)^{-1} g(x)}=\sqrt{\sum_{i=1}^{n} x_{i}^{2}}=\|x\|
$$

This implies that $x=0$ is the unique minimizer. The Newton step at $x$ is given by

$$
\Delta x=-H(x)^{-1} g(x)=-x(e+x)
$$

and a full Newton step yields

$$
x^{+}=x-x(e+x)=-x^{2} .
$$

The Newton step is feasible only if $-x^{2}>-e$, i.e. $x^{2}<e$; this certainly holds if $\delta(x)<1$. Note that the theory guarantees feasibility only in that case. Moreover, if the Newton step is feasible then

$$
\delta\left(x^{+}\right)=\left\|x^{2}\right\| \leq\|x\|_{\infty}\|x\| \leq \delta(x)^{2}
$$

and this is better than the theoretical result of Lemma 6.18. When we take a damped Newton step, with the default step size $\alpha=\frac{1}{1+\delta(x)}$, the next iterate is given by

$$
x^{+}=x-\frac{x(e+x)}{1+\|x\|} .
$$

If we apply Newton's method until $\delta(x) \leq \epsilon$ then the output is an $x$ such that $\|x\| \leq \epsilon$.
Example 6.48 Consider the function $f:(0, \infty) \rightarrow \mathbb{R}$ defined by

$$
f(x)=x \log x-\log x, x>0
$$

This is the barrier of the entropy function, considered in Example 6.17, where we found that $f$ is 1 -self-concordant. One has

$$
f^{\prime}(x)=\frac{x-1}{x}+\log x, \quad f^{\prime \prime}(x)=\frac{x+1}{x^{2}}
$$

Therefore,

$$
\delta(x)=\sqrt{\frac{f^{\prime}(x)^{2}}{f^{\prime \prime}(x)}}=\sqrt{\frac{(x-1+x \log x)^{2}}{1+x}}=\frac{|x-1+x \log x|}{\sqrt{1+x}}
$$

In this case $x=e$ is the unique minimizer. The Newton step at $x$ is given by

$$
\Delta x=-\frac{f^{\prime}(x)}{f^{\prime \prime}(x)}=-\frac{x(x-1+x \log x)}{1+x}
$$

and a full Newton step yields

$$
x^{+}=x-\frac{x(x-1+x \log x)}{1+x}=\frac{x(2-x \log x)}{1+x} .
$$

When we take a damped Newton step, with the default step size $\alpha=\frac{1}{1+\delta(x)}$, the next iterate is given by

$$
x^{+}=x-\frac{x(x-1+x \log x)}{(1+x)(1+\delta x)} .
$$

We conclude this example with a numerical experiment. If we start at $x=10$ we get as output the figures in the following tableau. In this tableau $k$ denotes the iteration number, $x^{k}$ the $k$-th iterate, $\delta\left(x^{k}\right)$ is the proximity value at $x^{k}$ and $\alpha_{k}$ the step size in the $k+1$-th iteration.

| $k$ | $x^{k}$ | $f\left(x^{k}\right)$ | $\delta\left(x^{k}\right)$ | $\alpha_{k}$ |
| ---: | ---: | ---: | ---: | ---: |
| 0 | 10.00000000000000 | 20.72326583694642 | 9.65615737513337 | 0.09384245791391 |
| 1 | 7.26783221086343 | 12.43198234403589 | 7.19322142387618 | 0.12205211457924 |
| 2 | 5.04872746432087 | 6.55544129967853 | 4.97000287092924 | 0.16750410705319 |
| 3 | 3.33976698811526 | 2.82152744553701 | 3.05643368252612 | 0.24652196443090 |
| 4 | 2.13180419256384 | 0.85674030296950 | 1.55140872104182 | 0.39194033937129 |
| 5 | 1.39932346194914 | 0.13416824208214 | 0.56132642454284 | 0.64048105782415 |
| 6 | 1.07453881397326 | 0.00535871156275 | 0.10538523300312 | 1.00000000000000 |
| 7 | 0.99591735745291 | 0.00001670208774 | 0.00577372342963 | 1.00000000000000 |
| 8 | 0.99998748482804 | 0.00000000015663 | 0.00001769912592 | 1.00000000000000 |
| 9 | 0.99999999988253 | 0.00000000000000 | 0.00000000016613 | 1.00000000000000 |

If we start at $x=0.1$ the output becomes

| $k$ | $x^{k}$ | $f\left(x^{k}\right)$ | $\delta\left(x^{k}\right)$ | $\alpha_{k}$ |
| ---: | ---: | ---: | ---: | ---: |
| 0 | 0.10000000000000 | 2.07232658369464 | 1.07765920479347 | 0.48131088953032 |
| 1 | 0.14945506622819 | 1.61668135596306 | 1.05829223631865 | 0.48583965986703 |
| 2 | 0.22112932596124 | 1.17532173793649 | 1.00679545093710 | 0.49830688998873 |
| 3 | 0.32152237588997 | 0.76986051286674 | 0.90755746327638 | 0.52423060340338 |
| 4 | 0.45458940014373 | 0.42998027395695 | 0.74937259761986 | 0.57163351098592 |
| 5 | 0.61604926491198 | 0.18599661844608 | 0.53678522950535 | 0.65070901307522 |
| 6 | 0.78531752299982 | 0.05188170346324 | 0.30270971353625 | 1.00000000000000 |
| 7 | 0.96323307457328 | 0.00137728412903 | 0.05199249905660 | 1.00000000000000 |
| 8 | 0.99897567517041 | 0.00000104977911 | 0.00144861398705 | 1.00000000000000 |
| 9 | 0.99999921284500 | 0.00000000000062 | 0.00000111320527 | 1.00000000000000 |
| 10 | 0.99999999999954 | 0.00000000000000 | 0.0000000000066 | 1.00000000000000 |

## Chapter 7

## Some specially structured problems

### 7.1 Introduction

In this chapter we show for some important classes of optimization problems that the logarithmic barrier function is self-concordant. Sometimes the logarithmic barrier function of the problem itself is not self-concordant, but it will be necessary to reformulate the problem. The approach in this section is based mainly on Den Hertog [13], Den Hertog et al. [15] and Jarre [21], except for the last subsection that deals with the semidefinite optimization problem. For a rich survey of applications of semidefinite optimization in control and system theory we refer to Vandenberghe and Boyd [44].

Before dealing with the problems we first present some lemmas that are quite helpful in recognizing self-concordant functions.

### 7.2 Some composition rules for self-concordant functions

The following lemma provides some composition rules for self-concordant functions.

Lemma 7.1 (Nesterov and Nemirovskii [33, 34])

- (addition and scaling) Let $\varphi_{i}$ be $\kappa_{i}$-self-concordant on $\mathcal{F}_{i}^{0}, i=1,2$, and $\rho_{1}, \rho_{2} \in \mathbb{R}^{+}$ then $\rho_{1} \varphi_{1}+\rho_{2} \varphi_{2}$ is $\kappa$-self-concordant on $\mathcal{F}_{1}^{0} \cap \mathcal{F}_{2}^{0}$, where $\kappa=\max \left\{\frac{\kappa_{1}}{\sqrt{\rho_{1}}}, \frac{\kappa_{2}}{\sqrt{\rho_{2}}}\right\}$.
- (affine invariance) Let $\varphi$ be $\kappa$-self-concordant on $\mathcal{F}^{0}$ and let $\mathcal{B}(y)=B y+b$ : $\mathbb{R}^{k} \rightarrow \mathbb{R}^{m}$ be an affine mapping such that $\mathcal{B}\left(\mathbb{R}^{k}\right) \cap \mathcal{F}^{0} \neq \emptyset$. Then $\varphi(\mathcal{B}()$.$) is$ $\kappa$-self-concordant on $\left\{y: \mathcal{B}(y) \in \mathcal{F}^{0}\right\}$.

The next lemma states that if the quotient of the third and second order derivative of $f(x)$ is bounded by the second order derivative of $-\sum_{i=1}^{n} \log x_{i}$, then the corresponding logarithmic barrier function is self-concordant. This lemma will help to simplify selfconcordance proofs in the sequel.

Lemma 7.2 Let $f(x) \in C^{3}\left(\mathcal{F}^{0}\right)$ and convex. If there exists a $\beta$ such that

$$
\begin{equation*}
\left|\nabla^{3} f(x)[h, h, h]\right| \leq \beta h^{T} \nabla^{2} f(x) h \sqrt{\sum_{i=1}^{n} \frac{h_{i}^{2}}{x_{i}^{2}}} \tag{7.1}
\end{equation*}
$$

$\forall x \in \mathcal{F}^{0}$ and $\forall h \in \mathbb{R}^{n}$, then

$$
\varphi(x):=\frac{f(x)}{\mu}-\sum_{i=1}^{n} \log x_{i}
$$

with $\mu>0$, is $\left(1+\frac{1}{3} \beta\right)$-self-concordant on $\mathcal{F}^{0}$, and

$$
\psi(t, x):=-q \log (t-f(x))-\sum_{i=1}^{n} \log x_{i}
$$

with $q \geq 1$, is $\left(1+\frac{1}{3} \beta\right)$-self-concordant on $\mathbb{R} \times \mathcal{F}^{0}$.
*Proof: We start by proving the first part of the lemma. Note that since (7.1) is scale independent, we may assume that $\mu=1$. Straightforward calculations yield

$$
\begin{align*}
\nabla \varphi(x)^{T} h & =\nabla f(x)^{T} h-\sum_{i=1}^{n} \frac{h_{i}}{x_{i}}  \tag{7.2}\\
h^{T} \nabla^{2} \varphi(x) h & =h^{T} \nabla^{2} f(x) h+\sum_{i=1}^{n} \frac{h_{i}^{2}}{x_{i}^{2}}  \tag{7.3}\\
\nabla^{3} \varphi(x)[h, h, h] & =\nabla^{3} f(x)[h, h, h]-2 \sum_{i=1}^{n} \frac{h_{i}^{3}}{x_{i}^{3}} . \tag{7.4}
\end{align*}
$$

We show that

$$
\begin{equation*}
\left(\nabla^{3} \varphi(x)[h, h, h]\right)^{2} \leq 4\left(1+\frac{1}{3} \beta\right)^{2}\left(h^{T} \nabla^{2} \varphi(x) h\right)^{3}, \tag{7.5}
\end{equation*}
$$

from which the first part of the lemma follows. Since $f$ is convex, the two terms on the right-hand side of (7.3) are nonnegative, i.e. the right-hand side can be abbreviated by

$$
\begin{equation*}
h^{T} \nabla^{2} \varphi(x) h=a^{2}+b^{2}, \tag{7.6}
\end{equation*}
$$

with $a, b \geq 0$. Because of (7.1) we have that

$$
\left|\nabla^{3} f(x)[h, h, h]\right| \leq \beta a^{2} b
$$

Using (6.8) we get

$$
\left|\sum_{i=1}^{n} \frac{h_{i}^{3}}{x_{i}^{3}}\right| \leq\left(\sum_{i=1}^{n} \frac{h_{i}^{2}}{x_{i}^{2}}\right)^{\frac{3}{2}}=b^{3} .
$$

Thus we can bound the right-hand side of (7.4) by

$$
\begin{equation*}
\left|\nabla^{3} \varphi(x)[h, h, h]\right| \leq \beta a^{2} b+2 b^{3} . \tag{7.7}
\end{equation*}
$$

It is straightforward to verify that

$$
\left(\beta a^{2} b+2 b^{3}\right)^{2} \leq 4\left(1+\frac{1}{3} \beta\right)^{2}\left(a^{2}+b^{2}\right)^{3}
$$

Together with (7.6) and (7.7) our claim (7.5) follows and hence the first part of the lemma.
Now we prove the second part of the lemma. Let

$$
\tilde{x}=\binom{t}{x}, \quad h=\left(\begin{array}{c}
h_{0}  \tag{7.8}\\
\vdots \\
h_{n}
\end{array}\right) \quad \text { and } g(\tilde{x})=t-f(x)
$$

then

$$
\begin{align*}
\psi(\tilde{x})= & -q \log g(\tilde{x})-\sum_{i=1}^{n} \log x_{i}  \tag{7.9}\\
\nabla \psi(\tilde{x})^{T} h= & -q \frac{\nabla g(\tilde{x})^{T} h}{g(\tilde{x})}-\sum_{i=1}^{n} \frac{h_{i}}{x_{i}}  \tag{7.10}\\
h^{T} \nabla^{2} \psi(\tilde{x}) h= & -q \frac{h^{T} \nabla^{2} g(\tilde{x}) h}{g(\tilde{x})}+q \frac{\left(\nabla g(\tilde{x})^{T} h\right)^{2}}{g(\tilde{x})^{2}}+\sum_{i=1}^{n} \frac{h_{i}^{2}}{x_{i}^{2}}  \tag{7.11}\\
\nabla^{3} \psi(\tilde{x})[h, h, h]= & -q \frac{\nabla^{3} g(\tilde{x})[h, h, h]}{g(\tilde{x})}+3 q \frac{\left(h^{T} \nabla^{2} g(\tilde{x}) h\right) \nabla g(\tilde{x})^{T} h}{g(\tilde{x})^{2}} \\
& -2 q \frac{\left(\nabla g(\tilde{x})^{T} h\right)^{3}}{g(\tilde{x})^{3}}-2 \sum_{i=1}^{n} \frac{h_{i}^{3}}{x_{i}^{3}} . \tag{7.12}
\end{align*}
$$

We show that

$$
\begin{equation*}
\left(\nabla^{3} \psi(\tilde{x})[h, h, h]\right)^{2} \leq 4\left(1+\frac{1}{3} \beta\right)^{2}\left(h^{T} \nabla^{2} \psi(\tilde{x}) h\right)^{3} \tag{7.13}
\end{equation*}
$$

which will prove the lemma. Since $g$ is concave, all three terms on the right-hand side of (7.11) are nonnegative, i.e. the right-hand side can be abbreviated by

$$
\begin{equation*}
h^{T} \nabla^{2} \psi(\tilde{x}) h=a^{2}+b^{2}+c^{2} \tag{7.14}
\end{equation*}
$$

with $a, b, c \geq 0$. Due to (7.1) we have

$$
\left|\frac{\nabla^{3} g(\tilde{x})[h, h, h]}{g(\tilde{x})}\right| \leq \beta a^{2} c,
$$

so that we can bound the right-hand side of (7.12) by

$$
\begin{equation*}
\left|\nabla^{3} \psi(\tilde{x})[h, h, h]\right| \leq \beta a^{2} c+3 a^{2} b+2 b^{3}+2 c^{3} \tag{7.15}
\end{equation*}
$$

It is straightforward to verify that

$$
\left(\beta a^{2} c+3 a^{2} b+2 b^{3}+2 c^{3}\right)^{2} \leq 4\left(1+\frac{1}{3} \beta\right)^{2}\left(a^{2}+b^{2}+c^{2}\right)^{3}
$$

by eliminating all odd powers in the second term via inequalities of the type $2 a b \leq a^{2}+b^{2}$. Together with (7.14) and (7.15) our claim (7.13) follows, and hence the lemma.

In the next sections we will show self-concordance for the logarithmic barrier function for several nonlinear optimization problems by showing that (7.1) is fulfilled. Below we will frequently and implicitly use the fact that $-\log \left(-g_{j}(x)\right)$ is 1-self-concordant on its domain whenever $g_{j}(x)$ is a linear or convex quadratic function. If $g_{j}(x)$ is linear this follows immediately by using Example 6.13 and the second part of Lemma 7.1.

Exercise 7.2 Prove that $-\log \left(-g_{j}(x)\right)$ is 1-self-concordant on its domain when $g_{j}(x)$ is a convex quadratic function.

### 7.3 Entropy optimization (EO)

In this and the subsequent sections the IPC is not assumed when we discuss duality results, but is assumed when we consider logarithmic barrier problems.

## Primal EO problem

The entropy optimization problem is defined as

$$
(E O P)\left\{\begin{array}{l}
\min c^{T} x+\sum_{i=1}^{n} x_{i} \log x_{i} \\
A x=b \\
x \geq 0
\end{array}\right.
$$

where $A \in \mathbb{R}^{m \times n}, b \in \mathbb{R}^{m}$ and $c \in \mathbb{R}^{n}$. Using Example 6.17 and the first composition rule in Lemma 7.1 it follows that the logarithmic barrier function is 1 -self-concordant.

## Dual EO problem

When writing the nonnegativity constraint $x \geq 0$ as

$$
x \in \mathcal{C}:=\{x: x \geq 0\},
$$

the Lagrange dual of $(E O P)$ is given by

$$
\sup \{\psi(y)\}
$$

where

$$
\psi(y)=\inf _{x \in \mathcal{C}}(L(x, y)),
$$

and where $L(x, y)$ denotes the Lagrange function of $(E O P)$ :

$$
L(x, y)=c^{T} x+\sum_{i=1}^{n} x_{i} \log x_{i}-y^{T}(A x-b)
$$

Fixing $y, L(x, y)$ is convex in $x$. Moreover, $L(x, y)$ goes to infinity if, for some $i, x_{i}$ goes to infinity, and the partial derivative to $x_{i}$ goes to infinity if $x_{i}$ approaches zero. We conclude from this that $L(x, y)$ is bounded below. Hence the minimum is attained and occurs where the gradient vanishes. The gradient of $L(x, y)$ with respect to $x$ is given by

$$
c+e+\log x-A^{T} y,
$$

and hence $L(x, y)$ is minimal if

$$
\log x_{i}=a_{i}^{T} y-c_{i}-1,1 \leq i \leq n
$$

where $a_{i}$ denotes the $i$-th column of $A$. From this we can solve $x$, namely

$$
x_{i}=e^{a_{i}^{T} y-c_{i}-1}, 1 \leq i \leq n .
$$

Substitution gives

$$
\begin{aligned}
\psi(y) & =c^{T} x+\sum_{i=1}^{n} x_{i} \log x_{i}-y^{T}(A x-b) \\
& =c^{T} x+\sum_{i=1}^{n} x_{i}\left(a_{i}^{T} y-c_{i}-1\right)-y^{T}(A x-b) \\
& =c^{T} x+y^{T} A x-c^{T} x-\sum_{i=1}^{n} x_{i}-y^{T}(A x-b) \\
& =b^{T} y-\sum_{i=1}^{n} x_{i} \\
& =b^{T} y-\sum_{i=1}^{n} e^{a_{i}^{T} y-c_{i}-1}
\end{aligned}
$$

Thus we have shown that the Lagrange dual of $(E O P)$ is given by

$$
(E O D) \quad \max b^{T} y-\sum_{i=1}^{n} e^{a_{i}^{T} y-c_{i}-1}
$$

It is quite surprising that this is an unconstrained (convex) optimization problem.

## Duality results

Here we just summarize some basic duality properties of EO problems. For details we refer to the paper of Kas and Klafszky [18].

Lemma 7.3 (Weak duality) If $x$ is feasible for $(E O P)$ and $y$ is feasible for (EOD) then

$$
c^{T} x+\sum_{i=1}^{n} x_{i} \log x_{i} \geq b^{T} y-\sum_{i=1}^{n} e^{a_{i}^{T} y-c_{i}-1}
$$

with equality if and only if for all $i=1, \cdots, n$ we have

$$
\begin{equation*}
x_{j}=e^{a_{j}^{T} y-c_{j}-1} \tag{7.16}
\end{equation*}
$$

The feasible region of $(E O P)$ is denoted by $\mathcal{P}$.
Corollary 7.4 If (7.16) holds for some $x \in \mathcal{P}$ and $y \in \mathbb{R}^{m}$ then they are both optimal and the duality gap is zero.

As we observed, both $(E O P)$ and $(E O D)$ are convex optimization problems.
Theorem 7.5 We have

1. If $(E O P)$ satisfy the Slater regularity condition and

$$
\nu=\inf \left\{\nu=c^{T} x^{*}+\sum_{i=1}^{n} x_{i}^{*} \log x_{i}^{*}: x \in \mathcal{P}\right\}>-\infty
$$

then there is an optimal $y^{*} \in \mathbb{R}^{m}$ such that

$$
\nu=b^{T} y^{*}-\sum_{i=1}^{n} e^{a_{i}^{T} y^{*}-c_{i}-1}
$$

2. If

$$
\nu=\sup \left\{b^{T} y-\sum_{i=1}^{n} e^{a_{i}^{T} y-c_{i}-1}\right\}<\infty
$$

then there is an optimal $x^{*} \in \mathcal{P}$ such that

$$
\nu=c^{T} x^{*}+\sum_{i=1}^{n} x_{i}^{*} \log x_{i}^{*}
$$

Observe that although EO problems are nonlinear, no regularity assumption is needed to guarantee zero duality gap. The primal problem $(E O P)$ always has an optimal solution if it is feasible or, equivalently if the dual objective function is bounded.

## Generalized EO (GEO)

The generalized entropy optimization problem is defined as

$$
(G E O)\left\{\begin{array}{l}
\min c^{T} x+\sum_{i=1}^{n} f_{i}\left(x_{i}\right) \\
A x=b \\
x \geq 0
\end{array}\right.
$$

where, for each $i$ there exists a positive number $\kappa_{i}$ such that the function $f_{i}:(0, \infty) \rightarrow$ $\mathbb{R}$ satisfies

$$
\left|f_{i}^{\prime \prime \prime}\left(x_{i}\right)\right| \leq \kappa_{i} \frac{f_{i}^{\prime \prime}\left(x_{i}\right)}{x_{i}}
$$

Obviously, the logarithmic barrier $x_{i} \log x_{i}-\log x_{i}$ of the entropy function $x_{i} \log x_{i}$ satisfies this condition, by Example 6.17. The class of GEO problems is studied in Ye and Potra [45] and Han et al ${ }^{1}$. [12].

[^6]
## Self-concordant barrier for (GEO)

Lemma 7.6 The logarithmic barrier function for the generalized entropy optimization problem (GEO) is $\left(1+\frac{1}{3} \max _{i} \kappa_{i}\right)$-self-concordant.

Proof: Using Lemma 7.1 it suffices to show that

$$
f_{i}\left(x_{i}\right)-\log x_{i}
$$

is $\left(1+\frac{1}{3} \kappa_{i}\right)$-self-concordant. Since

$$
\left|f_{i}^{\prime \prime \prime}\left(x_{i}\right)\right| \leq \kappa_{i} f_{i}^{\prime \prime}\left(x_{i}\right) \frac{1}{x_{i}},
$$

this immediately follows from Lemma 7.2.

## Self-concordant barrier for (EOD)

We found that the dual entropy optimization problem is

$$
(E O D) \quad \max b^{T} y-\sum_{i=1}^{n} e^{a_{i}^{T} y-c_{i}-1}
$$

and this is an unconstrained problem. Thus there is no natural logarithmic barrier function for this problem. However, we can reformulate the problem by using the substitution

$$
z_{i}=e^{a_{i}^{T} y-c_{i}-1}, 1 \leq i \leq n .
$$

Then $z_{i}>0$ and we have the equivalent problem

$$
\left(E O D^{\prime}\right) \quad\left\{\begin{array}{l}
\max b^{T} y-\sum_{i=1}^{n} z_{i} \\
a_{i}^{T} y-c_{i}-1-\log z_{i} \leq 0,1 \leq i \leq n \\
z_{i}>0,1 \leq i \leq n
\end{array}\right.
$$

Note that

$$
a_{i}^{T} y-c_{i}-1-\log z_{i} \leq 0
$$

is equivalent to

$$
e^{a_{i}^{T} y-c_{i}-1} \leq z_{i},
$$

and, since we are maximizing, at optimality we will have equality, for each $i$. The logarithmic barrier function for $\left(E O D^{\prime}\right)$ is

$$
\frac{-b^{T} y+\sum_{i=1}^{n} z_{i}}{\mu}-\sum_{j=1}^{n} \log \left(\log z_{i}-a_{i}^{T} y+c_{i}+1\right)-\sum_{j=1}^{n} \log z_{i} .
$$

It can be shown that this barrier function is 2-self-concordant. Note that the first term is linear and hence 0 -self-concordant. It follows from the next exercise and Lemma 7.1 that the second term is 2 -self-concordant. See also Jarre [21].

$$
\phi(t, x):=-\log (\log t-x)-\log t, \quad t>0, x \in \mathbb{R}, \log t-x>0,
$$

is 2-self-concordant on its domain.

### 7.4 Geometric optimization (GO)

## Primal GO problem

The primal GO problem is defined as follows. Let $\left\{I_{k}\right\}_{k=1}^{r}$ be a partition of $I=$ $\{1, \cdots, n\}$, so

$$
\bigcup_{k=1}^{r} I_{k}=\{1, \cdots, n\} \text { and } I_{k} \cap I_{\ell}=\emptyset \text { for } k \neq \ell
$$

Let $a_{i} \in \mathbb{R}^{m}, c_{i} \in \mathbb{R} \forall i \in I$ and $b \in \mathbb{R}^{m}$ be given. The primal GO problem is then given by

$$
(P G O)\left\{\begin{array}{l}
\max b^{T} y \\
G_{k}(y)=\sum_{i \in I_{k}} e^{a_{i}^{T} y-c_{i}} \leq 1, k=1, \cdots, r
\end{array}\right.
$$

N.B. Here and elsewhere in this section the symbol $e$ represents the base of the natural logarithm, and not the all-one vector!

## Posynom optimization

GO problems occur frequently in another equivalent form, in the so-called posynomial form. Note that

$$
\begin{aligned}
G_{k}(y) & =\sum_{i \in I_{k}} e^{\sum_{j=1}^{m} a_{i j} y_{j}-c_{i}} \\
& =\sum_{i \in I_{k}} e^{-c_{i}} \prod_{j=1}^{m} e^{a_{i j} y_{j}} \\
& =\sum_{i \in I_{k}} \alpha_{i} \prod_{j=1}^{m}\left(e^{y_{j}}\right)^{a_{i j}} \\
& =\sum_{i \in I_{k}} \alpha_{i} \prod_{j=1}^{m} \tau_{j}^{a_{i j}}
\end{aligned}
$$

where $\alpha_{i}=e^{-c_{i}}>0$ and $\tau_{j}=e^{y_{j}}>0$. The above polynomial is called a posynomial because all coefficients $\alpha_{i}$ and all variables $\tau_{j}$ are positive. Observe that the substitution $\tau_{j}=e^{y_{j}}$ convexifies the posynomial.
Recently the 'transistor sizing problem', which involves minimizing the active area of an electrical circuit subject to circuit delay specifications, has been modeled and solved by using a posynomial model [41].

## Dual GO problem

The Lagrange dual of (GPO) can be transformed to the form

$$
(D G O)\left\{\begin{array}{l}
\min c^{T} x+\phi(x) \\
A x=b \\
x \geq 0
\end{array}\right.
$$

where

$$
\phi(x)=\sum_{k=1}^{r}\left(\sum_{i \in I_{k}} x_{i} \log x_{i}-\left(\sum_{i \in I_{k}} x_{i}\right) \log \left(\sum_{i \in I_{k}} x_{i}\right)\right)
$$

Note that if $r=1$ and $\sum_{i \in I} x_{i}=1$ then $(D G O)$ is just the primal entropy optimization problem $(E O P)$. This case occurs in applications in statistical information theory.
Actually, to get the above formulation of the dual problem goes beyond the scope of this course. We refer to Duffin, Peterson and Zener [8] and Klafszky [24]. Let us only mention that in the dualization process is it used that $\phi(x)$ can also be written as

$$
\phi(x)=\sum_{k=1}^{r} \log \frac{\prod_{i \in I_{k}} x_{i}^{x_{i}}}{\left(\sum_{i \in I_{k}} x_{i}\right)^{\left(\sum_{i \in I_{k}} x_{i}\right)}}
$$

and this can be bounded by the Generalized Arithmetic-Geometric Mean Inequality. ${ }^{2}$
Exercise 7.4 Derive the dual geometric optimization problem (DGO) from the primal problem ( $P G O$ ) by using Lagrange duality.

## Duality results

Here we just summarize some basic duality properties of GO problems.
${ }^{2}$ The Generalized Arithmetic-Geometric Mean Inequality states that

$$
\left(\frac{\sum_{i=1}^{n} \alpha_{i}}{\sum_{i=1}^{n} \beta_{i}}\right)^{\sum_{i=1}^{n} \beta_{i}} \geq \prod_{i=1}^{n}\left(\frac{\alpha_{i}}{\beta_{i}}\right)^{\beta_{i}}
$$

where $\alpha=\left(\alpha_{1}, \cdots, \alpha_{n}\right) \geq 0$ and $\beta=\left(\beta_{1}, \cdots, \beta_{n}\right)>0$. Equality holds if and only if $\alpha=\lambda \beta$ for some nonnegative $\lambda$. The inequality is also valid for $\beta=\left(\beta_{1}, \cdots, \beta_{n}\right) \geq 0$ if we define

$$
\left(\frac{x_{i}}{0}\right)^{0}:=1
$$

Lemma 7.7 (Weak duality) If $y$ is feasible for (PGO) and $x$ is feasible for (DGO) then

$$
b^{T} y \leq c^{T} x+\phi(x)
$$

with equality if and only if for each $k=1, \cdots, r$ and $j \in I_{k}$

$$
\begin{equation*}
x_{j}=e^{a_{j}^{T} y-c_{j}} \sum_{i \in I_{k}} x_{i} . \tag{7.17}
\end{equation*}
$$

The feasible regions of ( $D G O$ ) and ( $P G O$ ) are denoted as $\mathcal{D}$ and $\mathcal{P}$ respectively.
Corollary 7.8 If (7.17) holds for some $x \in \mathcal{D}$ and $y \in \mathcal{P}$ then they are both optimal and the duality gap is zero.

Most of the following observations are trivial, some of them needs a nontrivial proof.

- $(P G O)$ is a convex optimization problem.
- $(D G O)$ is a convex optimization problem.

Proof: $\phi(x)=\sum_{k=1}^{r} \phi_{k}(x)$ where

$$
\phi_{k}(x)=\log \frac{\prod_{i \in I_{k}} x_{i}^{x_{i}}}{\left(\sum_{i \in I_{k}} x_{i}\right)^{\left(\sum_{i \in I_{k}} x_{i}\right)}}
$$

is positive homogeneous of order $1\left(\phi_{k}(\lambda x)=\lambda \phi_{k}(x)\right.$ for all $\left.\lambda>0\right)$ and subadditive (which means that $\phi_{k}\left(x^{1}+x^{2}\right) \leq \phi_{k}\left(x^{1}\right)+\phi_{k}\left(x^{2}\right)$ ).

- If $\left|I_{k}\right|=1 \forall i$ then $(P G O)$ is equivalent to an LO problem.
- If $x$ is optimal for $(D G O)$ and $x_{i}=0$ for some $i \in I_{k}$ then $x_{i}=0, \forall i \in I_{k}$.
- $G_{k}(y)$ is $\operatorname{logarithmically}$ convex (i.e. $\log G_{k}(y)$ is convex). This easily follows by using Hölder's inequality. ${ }^{3}$

Theorem 7.9 We have

1. If $(P G O)$ satisfy the Slater regularity condition and $\nu=\sup \left\{b^{T} y: y \in \mathcal{P}\right\}<\infty$ then there is an optimal $x^{*} \in \mathcal{D}$ such that $\nu=c^{T} x^{*}+\phi\left(x^{*}\right)$.
2. If $(D G O)$ satisfy the Slater regularity condition and $\nu=\inf \left\{c^{T} x+\phi(x): x \in\right.$ $\mathcal{D}\}>-\infty$ then there is an optimal $y^{*} \in \mathcal{P}$ such that $\nu=b^{T} y^{*}$.

[^7]whenever $\left(\alpha_{1}, \alpha_{2}, \cdots, \alpha_{n}\right) \geq 0,\left(\beta_{1}, \beta_{2}, \cdots, \beta_{n}\right) \geq 0$ and $0 \leq \lambda \leq 1$.
3. If both of (PGO) and (DGO) are feasible but none of them is Slater regular, then
$$
\sup \left\{b^{T} y: y \in \mathcal{P}\right\}=\inf \left\{c^{T} x+\phi(x): x \in \mathcal{D}\right\}
$$
*Proof: (Sketch)

- 1. and 2. follow by using the Convex Farkas Lemma, or the Karush-Kuhn-Tucker Theorem.
- The proof of 3. consists of several steps.
- First we reduce the dual problem to $\left(D G O_{r}\right)$ by erasing all the variables $x_{i}$ that are zero at all dual feasible solutions.
- Clearly ( $D G O$ ) and $\left(D G O_{r}\right)$ are equivalent.
- By construction $\left(D G O_{r}\right)$ is Slater regular.
- Form the primal $\left(P G O_{r}\right)$ of $\left(D G O_{r}\right)$.
- Due to 2. $\left(P G O_{r}\right)$ has optimal solution with optimal value equal to the optimal value of $\left(D G O_{r}\right)$.
- It then remains to prove that the optimal values of $(P G O)$ and $\left(P G O_{r}\right)$ are equal.


## Self-concordant barriers

We have the following result for GO.
Lemma 7.10 The logarithmic barrier functions of both (PGO) and (DGO) are 2-selfconcordant. ${ }^{4}$
*Proof: We give the proof for the dual GO problem. Because of Lemma 7.1, it suffices to verify 2-self-concordance for the following logarithmic barrier function

$$
\begin{equation*}
\varphi(x)=\sum_{i \in I_{k}} x_{i} \log x_{i}-\left(\sum_{i \in I_{k}} x_{i}\right) \log \left(\sum_{i \in I_{k}} x_{i}\right)-\sum_{i \in I_{k}} \log x_{i}, \tag{7.18}
\end{equation*}
$$

for some fixed $k$. For simplicity, we will drop the subscript $i \in I_{k}$. Now we can use Lemma 7.2 , so that we only have to verify (7.1) for

$$
f(x):=\sum x_{i} \log x_{i}-\left(\sum x_{i}\right) \log \left(\sum x_{i}\right),
$$

and $\beta=3$, which is equivalent to the following inequality:

$$
\begin{equation*}
\left|\sum \frac{h_{i}^{3}}{x_{i}^{2}}-\frac{\left(\sum h_{i}\right)^{3}}{\left(\sum x_{i}\right)^{2}}\right| \leq 3\left(\sum \frac{h_{i}^{2}}{x_{i}}-\frac{\left(\sum h_{i}\right)^{2}}{\sum x_{i}}\right) \sqrt{\sum \frac{h_{i}^{2}}{x_{i}^{2}}} . \tag{7.19}
\end{equation*}
$$

Here $x_{i}>0$ and $h_{i}$ arbitrary. To prove this inequality, let us define:

$$
\xi_{i}=x_{i}^{-\frac{1}{2}}\left(h_{i}-\frac{\sum h_{j}}{\sum x_{j}} x_{i}\right) .
$$

[^8]Note that

$$
\sum x_{i}^{\frac{1}{2}} \xi_{i}=0
$$

Using this substitution, we can rewrite the left-hand side of the inequality (7.19):

$$
\begin{align*}
\left|\sum \frac{h_{i}^{3}}{x_{i}^{2}}-\frac{\left(\sum h_{i}\right)^{3}}{\left(\sum x_{i}\right)^{2}}\right| & =\left|\sum\left(x_{i}^{-\frac{1}{2}} \xi_{i}^{3}+3 \xi_{i}^{2} \frac{\sum h_{j}}{\sum x_{j}}\right)\right| \\
& =\left|\sum \xi_{i}^{2}\left(x_{i}^{-\frac{1}{2}} \xi_{i}+3 \frac{\sum h_{j}}{\sum x_{j}}\right)\right| \\
& =\left|\sum \xi_{i}^{2}\left(\frac{h_{i}}{x_{i}}-\frac{\sum h_{j}}{\sum x_{j}}+3 \frac{\sum h_{j}}{\sum x_{j}}\right)\right| \\
& =\left|\sum \xi_{i}^{2}\left(\frac{h_{i}}{x_{i}}+2 \frac{\sum h_{j}}{\sum x_{j}}\right)\right| \\
& \leq \sum \xi_{i}^{2} \frac{h_{i} \mid}{x_{i}}+2 \sum \xi_{i}^{2} \frac{\left|\sum h_{j}\right|}{\sum x_{j}} \\
& \leq 3 \sum \xi_{i}^{2} \sqrt{\sum \frac{h_{j}^{2}}{x_{j}^{2}}} \tag{7.20}
\end{align*}
$$

where the last inequality follows because

$$
\frac{\left|h_{i}\right|}{x_{i}} \leq \sqrt{\sum \frac{h_{j}^{2}}{x_{j}^{2}}}
$$

and

$$
\begin{equation*}
\sum \frac{h_{i}^{2}}{x_{i}^{2}}\left(\sum x_{i}\right)^{2} \geq \sum \frac{h_{i}^{2}}{x_{i}^{2}} \sum x_{i}^{2} \geq\left(\sum\left|h_{i}\right|\right)^{2} \tag{7.21}
\end{equation*}
$$

(The last inequality in (7.21) follows directly from the Cauchy-Schwartz inequality.) Now note that the right-hand side of (7.19) is equal to

$$
3 \sum \xi_{i}^{2} \sqrt{\sum \frac{h_{i}^{2}}{x_{i}^{2}}}
$$

Together with (7.20), this completes the proof.

## $7.5 l_{p}$-norm optimization (NO)

## The primal $l_{p}$-norm optimization problem

As in the previous section, let $\left\{I_{k}\right\}_{k=1}^{r}$ be a partition of $I=\{1, \cdots, n\}$ :

$$
\bigcup_{k=1}^{r} I_{k}=\{1, \cdots, n\} \text { and } I_{k} \cap I_{\ell}=\emptyset \text { for } k \neq \ell
$$

and let $p_{i}>1, i=1, \cdots, n$. Moreover, let $a_{i} \in \mathbb{R}^{m}, b_{i} \in \mathbb{R}^{m}, c_{i} \in \mathbb{R}$ for all $i \in I$. Then the primal $l_{p}$-norm optimization problem $[35]-[37]$, [43] can be formulated as

$$
\left(P l_{p}\right)\left\{\begin{array}{l}
\max \eta^{T} y \\
G_{k}(y):=\sum_{i \in I_{k}} \frac{1}{p_{i}}\left|a_{i}^{T} y-c_{i}\right|^{p_{i}}+b_{k}^{T} y-d_{k} \leq 0, k=1, \cdots, r
\end{array}\right.
$$

One may easily verify that $\left(P l_{p}\right)$ is a convex optimization problem. Note that in spite of the absolute value in the problem formulation the functions $G_{k}(y)$ are infinitely many times differentiable.

## Special cases of $l_{p}$-optimization:

- if $I_{k}=\emptyset$ for all $k$ then $\left(P l_{p}\right)$ becomes an LO problem;
- the problem QO, as considered in Part I, and which has a convex quadratic objective function and linear constraints, can be reduced to the special case of $\left(P l_{p}\right)$ where all but one $I_{k}=\emptyset$ and, moreover, for the nonempty $I_{k}$ we have $p_{i}=2$ for all $i \in I_{k}$;
- if for all $i$ we have $p_{i}=2$ then $\left(P l_{p}\right)$ is a convex quadratically constrained optimization problem;
- $l_{p}$-norm approximation is the case where for all $k$ we have $b_{k}=0$;


## The dual $l_{p}$-norm optimization problem

Let $q_{i}$ be such that $\frac{1}{p_{i}}+\frac{1}{q_{i}}=1$. Moreover, let $A$ be the matrix whose columns are $a_{i}$, $i=1, \cdots, n$, and $B$ the matrix whose columns are $b_{k}, k=1, \cdots, r$. Then, the dual of the $l_{p}$-norm optimization problem $\left(P l_{p}\right)$ is (see [35]-[37], [43])

$$
\left(D l_{p}\right)\left\{\begin{array}{l}
\min \psi(x, z):=c^{T} x+d^{T} z+\sum_{k=1}^{r} z_{k} \sum_{i \in I_{k}} \frac{1}{q_{i}}\left|\frac{x_{i}}{z_{k}}\right|^{q_{i}} \\
A x+B z=\eta \\
z \geq 0
\end{array}\right.
$$

If $x_{i} \neq 0$ and $z_{k}=0$, then $z_{k}\left|\frac{x_{i}}{z_{k}}\right|^{q_{i}}$ is defined as $\infty$. If $x_{i}=0$ for all $i \in I_{k}$ and $z_{k}=0$ then we define

$$
z_{k}\left|\frac{x_{i}}{z_{k}}\right|^{q_{i}}
$$

to be zero.
Also $\left(D l_{p}\right)$ is a convex optimization problem. This follows by noting that

$$
\sum_{i \in I_{k}} \frac{1}{q_{i}} z_{k}^{q_{i}-1}\left|x_{i}\right|^{q_{i}}
$$

is positive homogeneous and subadditive, hence convex. Moreover, the dual feasible region is a polyhedron - not necessarily closed.
As in the case of geometric optimization, the derivation of the above formulation of the dual problem goes beyond the scope of this course. We refer to Peterson and Ecker [35, 36, 37] and Terlaky [43]. ${ }^{5}$

[^9]Exercise 7.5 Derive the dual $l_{p}$ optimization problem $\left(D l_{p}\right)$ from the primal problem $\left(D l_{p}\right)$ by using Lagrange duality.

## Duality results

Proposition 7.11 (Weak duality) If $y$ is feasible for $\left(P l_{p}\right)$ and $(x, z)$ for $\left(D l_{p}\right)$ then

$$
\eta^{T} y \leq \psi(x, z)
$$

with equality if and only if for each $k=1, \cdots, r$ and $i \in I_{k}, z_{k} G_{k}(y)=0$ and either $z_{k}=0$ or

$$
\frac{x_{i}}{z_{k}}=\operatorname{sign}\left(a_{i}^{T} y-c_{i}\right)\left|a_{i}^{T} y-c_{i}\right|^{p_{i}-1}
$$

or equivalently

$$
a_{i}^{T} y-c_{i}=\operatorname{sign}\left(x_{i}\right)\left|\frac{x_{i}}{z_{k}}\right|^{q_{i}-1}
$$

In fact $\left(D l_{p}\right)$ is a slightly transformed Lagrange (Wolfe) dual of $\left(P l_{p}\right)$. We have the following duality results.

Theorem 7.12 We have:

1. If $\left(P l_{p}\right)$ satisfy the Slater regularity condition and $\nu=\sup \left\{\eta^{T} y: y \in\left(P l_{p}\right)\right\}<\infty$ then there is an optimal $\left(x^{*}, z^{*}\right) \in\left(D l_{p}\right)$ with $\psi\left(x^{*}, z^{*}\right)=\nu$.
2. If $\left(D l_{p}\right)$ satisfy the Slater regularity condition and $\nu=\inf \{\psi(x, z):(x, z) \in$ $\left.\left(D l_{p}\right)\right\}>-\infty$ then there is an optimal $y^{*} \in\left(P l_{p}\right)$ such that $\nu=\eta^{T} y^{*}$.
3. If both of $\left(P l_{p}\right)$ and $\left(D l_{p}\right)$ are feasible but none of them is Slater regular, then $\max \left\{\eta^{T} y: y \in\left(P l_{p}\right)\right\}=\inf \left\{\psi(x, z):(x, z) \in\left(D l_{p}\right)\right\}$.

## *Proof: (Sketch)

- 1. and 2. follow by using the Convex Farkas Lemma, or the Karush-Kuhn-Tucker Theorem.
- The proof of 3 . is more involved. It consists of the following steps.
- Reduce the dual problem to $\left(D l_{p}\right)_{r}$ by erasing all the variables $x_{i} \in I_{k}$ and $z_{k}$ if $z_{k}$ is zero at all dual feasible solutions.
- Clearly $\left(D l_{p}\right)$ and $\left(D l_{p}\right)_{r}$ are equivalent. By construction $\left(D l_{p}\right)_{r}$ is Slater regular.
- Form the primal $\left(P l_{p}\right)_{r}$ of $\left(D l_{p}\right)_{r}$. Due to $2\left(P l_{p}\right)_{r}$ has optimal solution with optimal value equal to the optimal value of $\left(D l_{p}\right)_{r}$.
- It then remains to prove that the optimal values of $\left(P l_{p}\right)$ and $\left(P l_{p}\right)_{r}$ are equal. Moreover $\left(P l_{p}\right)$ has an optimal solution.

Then

$$
\alpha \beta \leq \frac{1}{p}|\alpha|^{p}+\frac{1}{q}|\beta|^{q}
$$

Equality holds if and only if

$$
\alpha=\operatorname{sign}(\beta)|\beta|^{q-1} \text { or } \beta=\operatorname{sign}(\alpha)|\alpha|^{p-1} .
$$

## Self-concordant barrier for the primal problem

To get a self-concordant barrier function for the primal $l_{p}$-norm optimization problem we need to reformulate it as:

$$
\left(P l_{p}^{\prime}\right)\left\{\begin{array}{l}
\max \eta^{T} y  \tag{7.22}\\
\sum_{i \in I_{k}} \frac{1}{p_{i}} t_{i}+b_{k}^{T} y-d_{k} \leq 0, k=1, \cdots, r \\
s_{i}^{p_{i}} \leq t_{i} \\
a_{i}^{T} y-c_{i} \leq s_{i} \\
-a_{i}^{T} y+c_{i} \leq s_{i} \\
s \geq 0
\end{array}\right\} i=1, \cdots, n
$$

The logarithmic barrier function for this problem can be proved to be self-concordant. Observe that in the transformed problem we have $4 n+r$ constraints, compared with $r$ in the original problem $\left(P l_{p}\right)$.

Lemma 7.13 The logarithmic barrier function for the reformulated $l_{p}$-norm optimization problem $\left(P l_{p}^{\prime}\right)$ is $\left(1+\frac{1}{3} \max _{i}\left|p_{i}-2\right|\right)$-self-concordant.
*Proof: $\quad$ Since $f\left(s_{i}\right):=s_{i}^{p_{i}}, p_{i} \geq 1$, satisfies (7.1) with $\beta=\left|p_{i}-2\right|$, we have from Lemma 7.2 that

$$
-\log \left(t_{i}-s_{i}^{p_{i}}\right)-\log s_{i}
$$

is $\left(1+\frac{1}{3}\left|p_{i}-2\right|\right)$-self-concordant. Consequently, it follows from Lemma 7.1 that the logarithmic barrier function for the reformulated primal $l_{p}$-norm optimization problem is $\left(1+\frac{1}{3} \max _{i}\left|p_{i}-2\right|\right)$-selfconcordant.

Note that the concordance parameter depends on $p_{i}$. We can improve it as follows. We replace the constraints $s_{i}^{p_{i}} \leq t_{i}$ by the equivalent constraints $s_{i} \leq t_{i}^{\pi_{i}}$, where $\pi_{i}=\frac{1}{p_{i}}$. Moreover, the redundant constraints $s \geq 0$ are replaced by $t \geq 0$. So, we obtain the following reformulated $l_{p}$-norm optimization problem:

$$
\left(P l_{p}^{\prime \prime}\right)\left\{\begin{array}{l}
\max \eta^{T} y  \tag{7.23}\\
\sum_{i \in I_{k}} \frac{1}{p_{i}} t_{i}+b_{k}^{T} y-d_{k} \leq 0, k=1, \cdots, r \\
s_{i} \leq t_{i}^{\pi_{i}} \\
a_{i}^{T} y-c_{i} \leq s_{i} \\
-a_{i}^{T} y+c_{i} \leq s_{i} \\
t \geq 0
\end{array}\right\} i=1, \cdots, n
$$

The following lemma improves the result of Lemma 7.13.

Lemma 7.14 The logarithmic barrier function for the reformulated $l_{p}$-norm optimization problem $\left(P l_{p}^{\prime \prime}\right)$ is $\frac{5}{3}$-self-concordant.
*Proof: $\quad$ Since $f\left(t_{i}\right):=-t_{i}^{\pi_{i}}, \pi_{i} \leq 1$, satisfies (7.1) with $\beta=\left|\pi_{i}-2\right|$, we have from Lemma 7.2 that

$$
-\log \left(t_{i}^{\pi_{i}}-s_{i}\right)-\log t_{i}
$$

is $\left(1+\frac{1}{3}\left|\pi_{i}-2\right|\right)$-self-concordant, where $\pi_{i} \leq 1$. Consequently, the corresponding logarithmic barrier function is $\frac{5}{3}$-self-concordant.

## Self-concordant barrier for the dual problem

The dual problem is equivalent to

$$
\left(D l_{p}^{\prime}\right)\left\{\begin{array}{l}
\min c^{T} x+d^{T} z+\sum_{i=1}^{n} \frac{1}{q_{i}} t_{i}  \tag{7.24}\\
s_{i}^{q_{i}} z_{k}^{-q_{i}+1} \leq t_{i}, \quad i \in I_{k}, \quad k=1, \cdots, r \\
x \leq s \\
-x \leq s \\
A x+B z=\eta \\
z \geq 0 \\
s \geq 0
\end{array}\right.
$$

Note that the original problem $\left(D l_{p}\right)$ has $r$ inequalities, and the reformulated problem $\left(D l_{p}^{\prime}\right) 4 n+r$. Now we prove the following lemma.

Lemma 7.15 The logarithmic barrier function of the reformulated dual $l_{p}$-norm optimization problem $\left(D l_{p}^{\prime}\right)$ is $\left(1+\frac{\sqrt{2}}{3} \max _{i}\left(q_{i}+1\right)\right)$-self-concordant.
*Proof: It suffices to show that

$$
-\log \left(t_{i}-s_{i}^{q_{i}} z_{k}^{-q_{i}+1}\right)-\log z_{k}-\log s_{i}
$$

is $\left(1+\frac{\sqrt{2}}{3}\left(q_{i}+1\right)\right.$ )-self-concordant, or equivalently by Lemma 7.2 , that (we will omit the subscript $i$ and $k$ in the sequel of this proof) $f(s, z):=s^{q} z^{-q+1}$ satisfies (7.1) for $\beta=\sqrt{2}(q+1)$, i.e. that

$$
\begin{equation*}
\left|\nabla^{3} f(s, z)[h, h, h]\right| \leq \sqrt{2}(q+1) h^{T} \nabla^{2} f(s, z) h \sqrt{\frac{\left|h_{1}\right|^{2}}{s^{2}}+\frac{\left|h_{2}\right|^{2}}{z^{2}}} \tag{7.25}
\end{equation*}
$$

where $h^{T}=\left(h_{1}, h_{2}\right)$. After doing some straightforward calculations we obtain for the second order term

$$
\begin{aligned}
h^{T} \nabla^{2} f(s, z) h & =q(q-1) s^{q-3} z^{-q-2}\left(s z^{3} h_{1}^{2}+s^{3} z h_{2}^{2}-2 s^{2} z^{2} h_{1} h_{2}\right) \\
& =q(q-1) s^{q-3} z^{-q-2}\left(z h_{1}-s h_{2}\right)^{2} s z,
\end{aligned}
$$

and for the third order term

$$
\begin{aligned}
\left|\nabla^{3} f(s, z)[h, h, h]\right|= & q(q-1) s^{q-3} z^{-q-2} \mid(q-2) z^{3} h_{1}^{3}-(q+1) s^{3} h_{2}^{3}- \\
& 3(q-1) s z^{2} h_{1}^{2} h_{2}+3 q s^{2} z h_{1} h_{2}^{2} \mid \\
= & q(q-1) s^{q-3} z^{-q-2}\left(z h_{1}-s h_{2}\right)^{2}\left|(q-2) z h_{1}-(q+1) s h_{2}\right| \\
\leq & q(q-1)(q+1) s^{q-3} z^{-q-2}\left(z h_{1}-s h_{2}\right)^{2}\left(z\left|h_{1}\right|+s\left|h_{2}\right|\right)
\end{aligned}
$$

Now we obtain

$$
\frac{\left|\nabla^{3} f(s, z)[h, h, h]\right|}{h^{T} \nabla^{2} f(s, z) h} \leq(q+1)\left(\frac{\left|h_{1}\right|}{s}+\frac{\left|h_{2}\right|}{z}\right) \leq \sqrt{2}(q+1) \sqrt{\frac{\left|h_{1}\right|^{2}}{s^{2}}+\frac{\left|h_{2}\right|^{2}}{z^{2}}}
$$

This proves (7.25) and hence the lemma.

We can improve this result as follows: the constraints $s_{i}^{q_{i}} z_{k}^{-q_{i}+1} \leq t_{i}$ are replaced by the equivalent constraints $t_{i}^{\rho_{i}} z_{k}^{-\rho_{i}+1} \geq s_{i}$, where $\rho_{i}:=\frac{1}{q_{i}}$, and the redundant constraints $s \geq 0$ are replaced by $t \geq 0$. The new reformulated dual $l_{p}$-norm optimization problem becomes:

$$
\left(D l_{p}^{\prime \prime}\right)\left\{\begin{array}{l}
\min c^{T} x+d^{T} z+\sum_{i=1}^{n} \frac{1}{q_{i}} t_{i} \\
s_{i} \leq t_{i}^{\rho_{i}} z_{k}^{-\rho_{i}+1}, \quad i \in I_{k}, \quad k=1, \cdots, r \\
x \leq s  \tag{7.26}\\
-x \leq s \\
A x+B z=\eta \\
z \geq 0 \\
t \geq 0
\end{array}\right.
$$

Lemma 7.16 The logarithmic barrier function of the reformulated dual $l_{p}$-norm optimization problem $\left(D l_{p}^{\prime \prime}\right)$ is 2-self-concordant.
*Proof: Similarly to the proof of Lemma 7.15, it can be proved that

$$
-\log \left(t_{i}^{\rho_{i}} z_{k}^{-\rho_{i}+1}-s_{i}\right)-\log t_{i},
$$

with $\rho_{i} \leq 1$, is $\left(1+\frac{\sqrt{2}}{3}\left(\rho_{i}+1\right)\right)$-self-concordant. The lemma follows now from Lemma 7.1 and from $\rho_{i} \leq 1$.

### 7.6 Semidefinite optimization (SDO)

The semidefinite optimization problem was introduced in Section 3.5. There its dual problem was derived by using the standard dualization method based on the Lagrange dual. For ease of understanding we recall the definition of the primal and the dual SDO problem.

Let $A_{0}, A_{1}, \cdots, A_{n} \in R^{m \times m}$ be symmetric matrices. Further let $c \in R^{n}$ be a given vector and $x \in R^{n}$ the vector of unknowns in which the optimization is done. The primal semidefinite optimization problem is defined as

$$
\begin{array}{rll}
(P S O) & \min & c^{T} x \\
& \text { s.t. } & -A_{0}+\sum_{k=1}^{n} A_{k} x_{k} \succeq 0,
\end{array}
$$

where $\succeq 0$ indicates that the left hand side matrix has to be positive semidefinite. Clearly the primal problem ( $P S O$ ) is a convex optimization problem since any convex combination of positive semidefinite matrices is also positive semidefinite. For convenience the notation

$$
F(x)=-A_{0}+\sum_{k=1}^{n} A_{k} x_{k}
$$

will be used.
The dual problem of semidefinite optimization is given as follows:

$$
\begin{aligned}
(D S O) \quad \max & \operatorname{Tr}\left(A_{0} Z\right) \\
\text { s.t. } & \operatorname{Tr}\left(A_{k} Z\right)=c_{k}, \quad \text { for all } \quad k=1, \cdots, n, \\
& Z \succeq 0,
\end{aligned}
$$

where $Z \in R^{m \times m}$ is the matrix of variables. Again, the problem $(D S O)$ is a convex optimization problem since the trace of a matrix is a linear function of the matrix and a convex combination of positive semidefinite matrices is positive semidefinite.

As we have seen the weak duality relation between ( $P S O$ ) and ( $D S O$ ) holds (see page 71).

Theorem 7.17 (Weak duality) If $x \in R^{n}$ is primal feasible and $Z \in R^{m \times m}$ is dual feasible, then

$$
c^{T} x \geq \operatorname{Tr}\left(A_{0} Z\right)
$$

with equality if and only if

$$
F(x) Z=0 .
$$

Because the semidefinite optimization problem is nonlinear, strong duality holds only if a certain regularity assumption, e.g. the Slater regularity assumption holds.

Exercise 7.6 Prove that (PSO) is Slater regular if and only if there is an $x \in \mathbb{R}^{n}$ such that $F(x)$ is positive definite.

Exercise 7.7 Prove that (DPSO) is Slater regular if and only if there is an $m \times m$ symmetric positive definite matrix $Z$ such that $\operatorname{Tr}\left(A_{k} Z\right)=c_{k}$, for all $k=1, \cdots, n$.

The above exercises show that the Slater regularity condition coincides with the interior point assumption in the case of semidefinite optimization.

We do not go into detailed discussion of applications and solvability of semidefinite optimization problems. For applications the interested reader is referred to [44].

Finally, we note that under the interior point assumption the function

$$
c^{T} x-\mu \log (\operatorname{det}(F(x)))
$$

is a self-concordant barrier for the problem (PSO), and the function

$$
\operatorname{Tr}\left(A_{0} Z\right)+\mu \log (\operatorname{det}(Z))
$$

is a self-concordant barrier for the problem (DSO). This results show that semidefinite optimization problems are efficiently solvable by interior point methods [34].

## Appendix A

## *Appendix

## A. 1 Some technical lemmas

We start with a slightly generalized version of the well-known Cauchy-Schwarz inequality. The classical Cauchy-Schwarz inequality follows by taking $A=M=I$ in the next lemma (where $I$ is the identity matrix).

Lemma A. 1 (Generalized Cauchy-Schwarz inequality) If $A, M$ are symmetric matrices with $\left|x^{T} M x\right| \leq x^{T} A x, \forall x \in \mathbb{R}^{n}$, then

$$
\left(a^{T} M b\right)^{2} \leq\left(a^{T} A a\right)\left(b^{T} A b\right), \forall a, b \in \mathbb{R}^{n}
$$

Proof: Note that $x^{T} A x \geq 0, \forall x \in \mathbb{R}^{n}$, so $A$ is positive semidefinite. Without loss of generality assume that $A$ is positive definite. Otherwise $A+\epsilon I$ is positive definite for all $\epsilon>0$, and we take the limit as $\epsilon \rightarrow 0$, with $a$ and $b$ fixed. The lemma is trivial if $a=0$ or $b=0$, so we assume that $a$ and $b$ are nonzero. Putting

$$
\mu:=\sqrt[4]{\frac{a^{T} A a}{b^{T} A b}}
$$

then it follows from

$$
a^{T} M b=\frac{1}{4}\left((a+b)^{T} M(a+b)-(a-b)^{T} M(a-b)\right)
$$

that

$$
\begin{aligned}
\left(a^{T} M b\right)^{2} & =\frac{1}{16}\left((a+b)^{T} M(a+b)-(a-b)^{T} M(a-b)\right)^{2} \\
& \leq \frac{1}{16}\left((a+b)^{T} A(a+b)+(a-b)^{T} A(a-b)\right)^{2} \\
& =\frac{1}{16}\left(2 a^{T} A a+2 b^{T} A b\right)^{2} \\
& =\frac{1}{4}\left(a^{T} A a+b^{T} A b\right)^{2} .
\end{aligned}
$$

When replacing $a$ by $a / \mu$ and $b$ by $b / \mu$ this implies

$$
\left(a^{T} M b\right)^{2}=\left(\left(\frac{a}{\mu}\right)^{T} M(\mu b)\right)^{2} \leq \frac{1}{4}\left(\frac{1}{\mu^{2}} a^{T} A a+\mu^{2} b^{T} A b\right)^{2}=\left(a^{T} A a\right)\left(b^{T} A b\right)
$$

which was to be shown.

Question: what if $M$ is not symmetric?
The following lemma gives an estimate for the spectral radius of a symmetric homogeneous trilinear form. The proof is due to Jarre [22].

## Lemma A. 2 (Spectral Radius for Symmetric Trilinear Forms)

Let a symmetric homogeneous trilinear form $M: \mathbb{R}^{n} \times \mathbb{R}^{n} \times \mathbb{R}^{n} \rightarrow \mathbb{R}$ be given by its coefficient matrix $M \in \mathbb{R}^{n \times n \times n}$. Let $A: \mathbb{R}^{n} \times \mathbb{R}^{n} \rightarrow \mathbb{R}$ be a symmetric bilinear form, with matrix $A \in \mathbb{R}^{n \times n}$, and $\mu>0$ a scalar such that

$$
M[x, x, x]^{2} \leq \mu A[x, x]^{3}=\mu\|x\|_{A}^{6}, \forall x \in \mathbb{R}^{n}
$$

Then

$$
|M[x, y, z]| \leq \mu\|x\|_{A}\|y\|_{A}\|z\|_{A}, \forall x, y, z \in \mathbb{R}^{n}
$$

Proof: Without loss of generality we assume that $\mu=1$. Otherwise we replace $A$ by $\sqrt[3]{\mu} A$. As in the proof of Lemma A. 1 we assume that $A$ is positive definite. Then, using the substitution

$$
M[x, y, z]:=M\left[A^{-\frac{1}{2}} x, A^{-\frac{1}{2}} y, A^{-\frac{1}{2}} z\right]
$$

we can further assume that $A=I$ is the identity matrix and we need to show that

$$
|M[x, y, z]| \leq \mu\|x\|_{2}\|y\|_{2}\|z\|_{2}, \forall x, y, z \in \mathbb{R}^{n}
$$

under the hypothesis

$$
|M[x, x, x]| \leq \mu\|x\|_{2}^{3}, \forall x \in \mathbb{R}^{n}
$$

For $x \in \mathbb{R}^{n}$ denote by $M_{x}$ the (symmetric) matrix defined by

$$
y^{T} M_{x} z:=M_{x}[y, z]:=M[x, y, z], \forall y, z \in \mathbb{R}^{n} .
$$

It is sufficient to show that

$$
|M[x, y, y]| \leq \mu\|x\|_{2}\|y\|_{2}^{2}, \forall x, y \in \mathbb{R}^{n}
$$

because the remaining part follows by applying Lemma A.1, with $M=M_{x}$, for fixed $x$.
Define

$$
\sigma:=\max \left\{M[x, y, y]:\|x\|_{2}=\|y\|_{2}=1\right\}
$$

and let $\bar{x}$ and $\bar{y}$ represent a solution of this maximization problem. The necessary optimality conditions for $\bar{x}$ and $\bar{y}$ imply that

$$
\binom{M_{\bar{y}} \bar{y}}{2 M_{\bar{y}} \bar{x}}=\alpha\binom{2 \bar{x}}{0}+\beta\binom{0}{2 \bar{y}},
$$

where $\alpha$ and $\beta$ are the Lagrange multipliers. From this we deduce that $\alpha=\sigma / 2$ and $\beta=\sigma$, by multiplying from the left with $\left(\bar{x}^{T}, 0\right)$ and $\left(0, \bar{y}^{T}\right)$, and thus we find

$$
M_{\bar{y}} \bar{y}=\sigma \bar{x}, \quad 2 M_{\bar{y}} \bar{x}=\sigma \bar{y},
$$

which implies that $M_{\bar{y}}^{2} \bar{y}=\sigma^{2} \bar{y}$. Since $M_{\bar{y}}$ is symmetric, it follows that $\bar{y}$ is an eigenvector of $M_{\bar{y}}$ with the eigenvalue $\pm \sigma$, which gives that

$$
\sigma=\left|\bar{y}^{T} M_{\bar{y}} \bar{y}\right|=M[\bar{y}, \bar{y}, \bar{y}]
$$

This completes the proof.

## A. 2 The proof of Theorem 6.3 and Lemma 6.4.

Theorem A. 3 Let us assume that for ( $C P O$ ) no bad ray exists (see Definition 6.2 on page 121). Then the following three statements are equivalent.
(i) (CPO) and (CDO) satisfy the interior point condition;
(ii) For each $\mu>0$ the relaxed KKT-system (6.2) has a solution;
(iii) For each $w>0\left(w \in \mathbb{R}^{m}\right)$ there exist $y$ and $x$ such that

$$
\begin{align*}
& \text { (ii) } \sum_{j=1}^{m} y_{j} \nabla g_{j}(x)=c, y \geq 0  \tag{i}\\
& \text { (iii) } \quad-y_{j} g_{j}(x)=w_{i}, \forall j=1, \cdots, m \tag{A.1}
\end{align*}
$$

Proof: The implications $($ iii $) \rightarrow(i i)$ and $(i i) \rightarrow(i)$ are evident. We need only to prove that (i) imply (iii).
Let us assume that the IPC holds. Let

$$
\mathcal{P}^{0}:=\left\{x \mid g_{j}(x)<0, \forall j\right\} \quad \text { and } \quad \mathcal{Y}:=\mathbb{R}_{+}^{m}=\left\{y \in \mathbb{R}^{m} \mid y>0\right\}
$$

Observe, that due to the IPC the set $\mathcal{P}^{0}$ is not empty and the solutions of (A.1) are exactly the saddle points of the function

$$
F(x, y):=-c^{T} x+\sum_{j=1}^{m} y_{j} g_{j}(x)+\sum_{j=1}^{m} w_{j} \log y_{j}
$$

on the set $\mathcal{P}^{0} \times \mathcal{Y}$, because the function $F(x, y)$ is convex in $x$ and concave in $y$. Thus we only need to prove that $F$ has a saddle point on $\mathcal{P}^{0} \times \mathcal{Y}$.

Observe, that for any fixed $x$ the function $F(x, y)$ attains its maximum in $\mathcal{Y}$ at the point

$$
y_{j}=\frac{w_{j}}{-g_{j}(x)}
$$

and this maximum, up to an additive constant, equals to

$$
\bar{F}(x):=-c^{T} x-\sum_{j=1}^{m} w_{j} \log \left(-g_{j}(x)\right)
$$

Consequently, to prove that $F(x, y)$ admits a saddle point on $\mathcal{P}^{0} \times \mathcal{Y}$ reduces to prove that $\bar{F}(x)$ attains its maximum on $\mathcal{P}^{0}$.
Next we prove that if a sequence $x^{i} \in \mathcal{P}^{0}, i=1,2, \cdots$ is minimizing for $\bar{F}(x)$, i.e

$$
\bar{F}\left(x^{i}\right) \longrightarrow \inf _{x \in \mathcal{P}^{0}} \bar{F}(x)
$$

then the sequence $\left\{x^{i}\right\}_{i=1}^{\infty}$ is bounded. ${ }^{1}$
Let us assume to the contrary that there exists a minimizing sequence $\left\{x^{i} \in \mathcal{P}^{0}\right\}_{i=1}^{\infty}$ with $\left\|x^{i}\right\| \rightarrow \infty$. Then the sequence $\tilde{x}^{i}:=\frac{x^{i}}{\left\|x^{i}\right\|}$ is bounded. We can choose a convergent subsequence of the bounded sequence $\tilde{x}^{i}$ (for simplicity denoted again the same way) with limit point

$$
s:=\lim _{i \rightarrow \infty} \tilde{x}^{i}
$$

Since $\mathcal{P}^{0}$ is convex and $x^{i} \in \mathcal{P}^{0}$ for all $i$, we conclude that $s$ is a recession direction of $\mathcal{P}^{0}$, i.e.

$$
x+\lambda s \in \mathcal{P}^{0} \quad \forall x \in \mathcal{P}^{0} \text { and } \lambda \geq 0
$$

We claim that

$$
\begin{equation*}
c^{T} s \geq 0 \tag{A.2}
\end{equation*}
$$

This is indeed true, otherwise we would have $c^{T} x^{i} \leq-\alpha\left\|x^{i}\right\|$ fore some $\alpha>0$ and $i$ large enough, and thus $\bar{F}\left(x^{i}\right)$ would go to infinity.

The convexity of the functions $g_{j}$ imply that there exists $0<\beta, \gamma \in \mathbb{R}$ such that for all $x$ we have $g_{j}(x) \geq-\beta-\gamma\|x\|$ (see Exercise A. 1 below). Then it follows that for all $i$, large enough, we have
$\bar{F}\left(x^{i}\right)=\bar{F}(x):=-c^{T} x^{i}-\sum_{j=1}^{m} w_{j} \log \left(-g_{j}\left(x^{i}\right)\right) \geq-\alpha\left\|x^{i}\right\|-\sum_{j=1}^{m} w_{j} \log (-\beta-\gamma\|x\|) \rightarrow \infty$
as $i \rightarrow \infty$. This is a contradiction, because $x^{i}$ is a minimizing sequence of $\bar{F}$.

[^10]Now let $(\bar{x}, \bar{y})$ be an interior solution of $(C D O)$. Then for all $j$ we have

$$
0 \geq g_{j}\left(x^{i}\right) \geq g_{j}(\bar{x})+\nabla g_{j}(\bar{x})^{T}\left(x^{i}-\bar{x}\right)
$$

Dividing by $\left\|x^{i}\right\|$ and taking the limit as $i \rightarrow \infty$ we obtain

$$
\begin{equation*}
\nabla g_{j}(\bar{x})^{T} s \leq 0 \tag{A.3}
\end{equation*}
$$

because $\left\|x^{i}\right\|$ goes to infinity.
On the other hand, using (A.2) and the equality constraint in ( $C D O$ ) we have

$$
0 \leq c^{T} s=\left(\sum_{j=1}^{m} \bar{y}_{j} \nabla g_{j}(\bar{x})\right)^{T} s=\sum_{j=1}^{m} \bar{y}_{j}\left(\nabla g_{j}(\bar{x})^{T} s\right) \leq 0,
$$

where the last inequality follows from (A.3). Now we have $c^{T} s=0$ and, due to $\bar{y}>0$ and (A.3) for each $j$ the equality

$$
\begin{equation*}
\nabla g_{j}(\bar{x})^{T} s=0 \tag{A.4}
\end{equation*}
$$

must hold too.
We are about concluding that the ray $\mathcal{R}=\{x \mid x=\bar{x}+\lambda s, \lambda \geq 0\}$ is a bad ray. To prove this we have to prove that $g_{j}$ is constant along $\mathcal{R}$. We already have seen that $s$ is a recession direction of $\mathcal{P}^{0}$, hence $s$ is also a recession direction of the larger set

$$
\mathcal{P}_{+}:=\left\{x \mid g_{j}(x) \leq \max \left\{0, g_{j}(\bar{x})\right\}, \text { for all } j\right\} \supseteq \mathcal{P}^{0} .
$$

Thus for each $j$ the functions $g_{j}$ are bounded from above on the ray $\mathcal{R}$, which together with (A.4) proves that $g_{j}$ is constant along $\mathcal{R}$. The proof is complete.

Exercise A. 1 Let $g(x): \mathbb{R}^{n} \rightarrow \mathbb{R}$ be a convex function. Prove that there exists $0<\beta, \gamma \in \mathbb{R}$ such that

$$
g_{j}(x) \geq-\beta-\gamma\|x\| \forall x \in \mathbb{R}^{n} .
$$

Remark: The implication $($ iii $) \rightarrow(i)$ may be false if a bad ray exists. Let us consider the following example.

Example A. 4 Let the convex set $\mathcal{Q}:=\left\{\left(x_{1}, x_{2}\right) \in \mathbb{R}^{2} \mid x_{1} \geq x_{2}^{2}\right\}$ be given and let $\pi(x)$ be the so-called Minkowski function of the set $\mathcal{Q}$ with the pole $\bar{x}=(1,0)$, i.e.

$$
\pi(x):=\min \left\{t \left\lvert\, \bar{x}+\frac{1}{t}(x-\bar{x}) \in \mathcal{Q}\right.\right\} .
$$

One easily checks that $\pi(x)$ is a nonnegative, convex function. Further, $\pi(x)=1$ if $x$ is a boundary point of the set $\mathcal{Q}$ and $\pi(x)=0$ on the ray $\mathcal{R}:=\{x \mid x=\bar{x}+\lambda(1,0), \lambda \geq 0\}$.

Setting $m:=1, c^{T} x:=x_{2}$ and $g_{1}(x):=\pi^{2}(x)+x_{2}$, we get a $(C P O)$ problem which satisfy the IPC, e.g. the point $x=(1,-0.5)$ is strictly feasible. Likewise $(C D O)$ satisfies the IPC, e.g. we may take $\tilde{x}=(1,0)$ and $\tilde{y}=1$.

However the system (A.1) never has a solution. Indeed, as we already know, saying that (A.1) has a solution is the same as to say that the function

$$
\bar{F}(x)=-c^{T} x-\mu \log \left(-g_{1}(x)\right) \equiv-x_{2}-\mu \log \left(-\pi^{2}(x)-x_{2}\right)
$$

attains its minimum in the set $\mathcal{P}^{0}=\left\{x \mid g_{1}(x)<0\right\}$. The latter is clearly not the case, since given any $\tilde{x}$ (e.g. we may take $\tilde{x}=(1,-0.5))$ with $g_{1}(\tilde{x})<0$ we always can find a better value of $\bar{F}$. One can take the points on the ray $\mathcal{R}^{\prime}:=\{x \mid x=\tilde{x}+\lambda(1,0)\}$, since $\pi(\tilde{x}+\lambda(1,0)) \rightarrow 0$ as $\lambda \rightarrow \infty$.

Lemma A. 5 Let us assume that for (CPO) the IPC holds and no bad line segment exists. Then the solutions of the systems (6.2) and (6.3), if they exist, are unique.

Proof: In the proof of Theorem A. 3 we have seen that a point $(x, s)$ solves (A.1) if and only if it is a saddle point of $F(x, y)$ on the set $\mathcal{P}^{0} \times \mathcal{Y}$. Clearly, for fixed $x$, the function $F$ is a strictly concave function of $y$, thus to prove that a saddle point is unique we only have to prove that the function

$$
\bar{F}(x):=-c^{T} x-\sum_{j=1}^{m} w_{j} \log \left(-g_{j}(x)\right)=\max _{y \in \mathcal{Y}} F(x, y)+\text { a constant }
$$

cannot attain its minimum at two different points. Let assume to the contrary that two distinct minimum points $x^{\prime}, x^{\prime \prime} \in \mathcal{P}^{0}$ exists. Due to the convexity of the function $\bar{F}$, we have that $\bar{F}(x)$ is constant on the line segment $\left[x^{\prime}, x^{\prime \prime}\right]$. This imply that both the first and second order directional derivatives of $\bar{F}(x)$ are zero on this line segment. This can only happen if the same is true for all the functions $g_{j}(x)$ separately, hence all the functions $g_{j}(x)$ are constant on the line segment $\left[x^{\prime}, x^{\prime \prime}\right]$, i.e. this line segment is bad. We have got a contradiction, thus the lemma is proved.

## Bibliography

[1] Anstreicher, K.M. (1990), A Standard Form Variant, and Safeguarded Linesearch, for the Modified Karmarkar Algorithm, Mathematical Programming 47, 337-351.
[2] M.S. Bazarraa, H.D. Sherali and C.M. Shetty, Nonlinear Programming: Theory and Algorithms, John Wiley and Sons, New York (1993).
[3] D.P. Bertsekas, Nonlinear Programming, Athena Scientific, Belmont, MA 02178-9998, (1995).
[4] V. Chvátal, Linear Programming, (W.H. Freeman and Company, 1983).
[5] R. W. Cottle and S-M. Guu. Two Characterizations of Sufficient Matrices. Linear Algebra and Its Applications, ??-??, 199?.
[6] R. W. Cottle, J.-S. Pang and V. Venkateswaran. Sufficient Matrices and the Linear Complementarity Problem. Linear Algebra and Its Applications, 114/115:231-249, 1989.
[7] R. W. Cottle, J.-S. Pang and R. E. Stone. The Linear Complementarity Problem. Academic, Boston, 1992.
[8] Duffin, R.J., Peterson, E.L. and Zener, C. (1967), Geometric Programming, John Wiley \& Sons, New York.
[9] P.E. Gill, W. Murray and M.H. Wright, Practical Optimization Academic Press, London, (1981).
[10] P.E. Gill, W. Murray and M.H. Wright, Numerical Linear Algebra and Optimization, Vol.1. Addison Wiley P.C. New York, (1991).
[11] R.A. Horn and C.R. Jonson, Matrix Analysis, Cambridge University Press, Cambridge, UK (1985).
[12] Han, C.-G., Pardalos, P.M. and Ye, Y. (1991), On Interior-Point Algorithms for Some Entropy Optimization Problems, Working Paper, Computer Science Department, The Pennsylvania State University, University Park, Pennsylvania.
[13] D. den Hertog, (1994), Interior Point Approach to Linear, Quadratic and Convex Programming: Algorithms and Complexity, Kluwer A.P.C., Dordrecht.
[14] D. den Hertog, C. Roos and T. Terlaky, The Linear Comlementarity Problem, Sufficient Matrices and the Criss-Cross Method, (1993) Linear Algebra and Its Applications, 187. 1-14.
[15] Hertog, D. den, Jarre, F., Roos, C. and Terlaky, T. (1995), A Sufficient Condition for Self-Concordance with Application to Some Classes of Structured Convex Programming Problems, Mathematical Programming 69, 75-88.
[16] J.-B. Hirriart-Urruty and C. Lemarèchal, Convex Analysis and Minimization Algorithms I and II, Springer Verlag, Berlin, Heidelberg (1993).
[17] R. Horst, P.M. Pardalos and N.V. Thoai, (1995) Introduction to Global Optimization, Kluwer A.P.C. Dordrecht.
[18] P. Kas and E. Klafszky, (1993) On the Dduality of the Mixed Entropy Programming, Optimization 27. 253-258.
[19] E. Klafszky and T. Terlaky, (1992) Some Generalizations of the Criss-Cross Method for Quadratic Programming, Mathemathische Operationsforschung und Statistics ser. Optimization 24. 127-139.
[20] Jarre, F. (1990), Interior-point Methods for Classes of Convex Programs. Technical Report SOL 90-16, Systems Optimization Laboratory, Department of Operations Research, Stanford University, Stanford, California.
[21] Jarre, F. (1996), Interior-point Methods for Convex Programming. In T. Terlaky (ed.), Interior-point Methods for Mathematical Programming, Kluwer A.P.C., Dordrecht, pp. 255-296.
[22] Jarre, F. (1994), Interior-point Methods via Self-concordance or Relative Lipschitz Condition. Habilitationsschrift, Würzburg, Germany.
[23] Karmarkar, N.K. (1984), A New Polynomial-Time Algorithm for Linear Programming, Combinatorica 4, 373-395.
[24] Klafszky, E. (1976), Geometric Programming, Seminar Notes 11.1976, Hungarian Committee for Systems Analysis, Budapest.
[25] Kortanek, K.O. and No, H. (1992), A Second Order Affine Scaling Algorithm for the Geometric Programming Dual with Logarithmic Barrier, Optimization 23, 501-507.
[26] D.C. Lay Linear Algebra and Its Applications Addision Wiley (1994).
[27] F. Lootsma, Algorithms for Unconstrained Optimization Dictaat Nr. a85A, en WI-385, Fac. TWI, TU Delft.
[28] F. Lootsma, Duality in Non-Linear Programming Dictaat Nr. a85D, Fac. TWI, TU Delft.
[29] F. Lootsma, Algorithms fo Constrained Optimization Dictaat Nr. a85B, Fac. TWI, TU Delft.
[30] H.M. Markowitz, (1956), The Optimization of a Quadratic Function Subject to Linear Constraints, Naval Research Logistics Quarterly 3, 111-133.
[31] H.M. Markowitz, (1959), Portfolio Selection, Efficient Diversification of Investments, Cowles Foundation for Research in Economics at Yale University, Monograph 16, John Wiley \& Sons, New York.
[32] Nemirovsky, A.S. (1999), Convex Optimization in Engineering, Dictaat Nr. WI-485, Fac. ITS/TWI, TU Delft.
[33] Nesterov, Y.E. and Nemirovsky, A.S. (1989), Self-Concordant Functions and Polynomial Time Methods in Convex Programming, Report, Central Economical and Mathematical Institute, USSR Academy of Science, Moscow, USSR.
[34] Nesterov, Y.E. and Nemirovsky, A.S. (1994), Interior-Point Polynomial Algorithms in Convex Programming, SIAM, Philadelphia.
[35] Peterson, E.L., and Ecker, J.G. (1970), Geometric Programming: Duality in Quadratic Programming and $l_{p}$ Approximation I, in: H.W. Kuhn and A.W. Tucker (eds.), Proceedings of the International Symposium of Mathematical Programming, Princeton University Press, New Jersey.
[36] Peterson, E.L., and Ecker, J.G. (1967), Geometric Programming: Duality in Quadratic Programming and $l_{p}$ Approximation II, SIAM Journal on Applied Mathematics 13, 317-340.
[37] Peterson, E.L., and Ecker, J.G. (1970), Geometric Programming: Duality in Quadratic Programming and $l_{p}$ Approximation III, Journal of Mathematical Analysis and Applications 29, 365-383.
[38] R.T. Rockafellar, Convex Analysis, Princeton, New Jersey, Princeton University Press (1970).
[39] C. Roos and T. Terlaky, Introduction to Linear Optimization Dictaat WI187, Fac. ITS/TWI, TU Delft (1997).
[40] C. Roos, T. Terlaky and J.-Ph. Vial (1997), Theory and Algorithms for Linear Optimization: An Interior Point Approach. John Wiley and Sons.
[41] S.S. Sapatnekar (1992), A convex programming approach to problems in VSLI designs. Ph. D. Thesis. University of Illinois at Urbana Campaign. 70-100.
[42] J. Stoer and C. Witzgall, Convexity and Optimization in Finite Dimensions I, Springer Verlag, Berlin, Heidelberg (1970).
[43] Terlaky, T. (1985), On $l_{p}$ programming, European Journal of Operational Research 22, 70-100.
[44] Vandenberghe, L. and Boyd, S. (1994) Semidefinite Programming, Report December 7, 1994, Information Systems Laboratory, Electrical Engineering Department, Stanford University, Stanford CA 94305.
[45] Ye, Y. and Potra, F. (1990), An Interior-Point Algorithm for Solving Entropy Optimization Problems with Globally Linear and Locally Quadratic Convergence Rate, Working Paper Series No. 90-22, Department of Management Sciences, The University of Iowa, Iowa City, Iowa. To Appear in SIAM Journal on Optimization.

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[^0]:    ${ }^{1}$ Lagrange multipliers were actually first introduced by Euler. (Lagrange was Euler's student.)

[^1]:    ${ }^{1}$ In the literature the recession cone is frequently referred to as the characteristic cone of the convex set $\mathcal{C}$.

[^2]:    ${ }^{1}$ Here, for the sake of simplicity, it is assumed that $\mathcal{C}=\mathbb{R}^{n}$. In other cases the negative gradient might point out of the feasible set $\mathcal{C}$.

[^3]:    ${ }^{1}$ The problem (NC) is in general not convex. It is a (CO) problem if and only if the functions $h_{j}$ are affine.

[^4]:    ${ }^{1}$ See Lemma 2.25 (page 64) in Den Hertog [13].

[^5]:    ${ }^{2}$ The property of $\psi$ below in fact means that $\psi(-t)$ is the conjugate of $\psi(t)$.

[^6]:    ${ }^{1}$ In this paper it is conjectured that these problems do not satisfy the self-concordance condition. The lemma below shows that it does satisfy the self-concordance condition.

[^7]:    ${ }^{3}$ Hölder's inequality states that

    $$
    \sum_{i=1}^{n} \alpha_{i}^{\lambda} \beta_{i}^{1-\lambda} \leq\left(\sum_{i=1}^{n} \alpha_{i}\right)^{\lambda}\left(\sum_{i=1}^{n} \beta_{i}\right)^{1-\lambda}
    $$

[^8]:    ${ }^{4}$ This contradicts the remark in [25] that the self-concordance property does not hold for the dual problem. One should realize, however, that the self-concordance property of the logarithmic barrier function depends on the representation of the dual problem. What we here show is just that the dual problem can be reformulated in such a way that its barrier function becomes self-concordant.

[^9]:    ${ }^{5}$ In the dualization process one needs the following inequality: Let $\alpha, \beta \in \mathbb{R}, p, q>1$ and

    $$
    \frac{1}{p}+\frac{1}{q}=1
    $$

[^10]:    ${ }^{1}$ Note that statement (ii) is an immediate consequence of this claim. Assuming that the sequence $\left\{x^{i}\right\}_{i=1}^{\infty}$ is bounded, then we can choose a converging subsequence with limit point $\bar{x}$. This limit point belongs to $\mathcal{P}^{0}$, because on any sequence converging to the boundary of $\mathcal{P}^{0}$ the function $\bar{F}(x)$ goes to infinity. Since $\bar{F}(x)$ is continuous, $\bar{x}$ is the desired minimizer that satisfy (ii).

