# Control and Cybernetics 

vol. 23 (1994) No. 1/2

# Regularity of solutions <br> in stability analysis of optimization and optimal control problems 

by
Kazimierz Malanowski*)
Systems Research Institute
Polish Academy of Sciences
Warsaw
Poland
Stability of solutions to cone-constrained optimization problems in Banach spaces is investigated in the situation where the so called norm discrepancy takes place. A two-norm approach is proposed, in which an additional information on regularity of solutions is used.

The obtained abstract stability results are applied to control and state constrained optimal control problems for nonlinear o.d.e..

## 1. Introduction

Stability of solutions to finite dimensional mathematical programs has been studied for many years and the obtained results are fairly complete (cf. e.g., Auslender,Cominetti,1990; Bonnans,1992; Gauvin, Janin,1988; Gollan,1981; Jittorntrum,1984; Kojima,1980; Robinson,1987).

In recent years there appeared also several papers devoted to stability analysis of cone-constrained optimization problems in Banach space (cf. Alt,1990A; Alt,1990B; Dontchev,Hager,1993; Dontchev at al.,1994; Ioffe,1991; Ito,Kunisch, 1992; Malanowski 1992; Malanowski,1993; Malanowski,1994; Shapiro,1992; Shapiro,1994; Shapiro,Bonnans,1992; Tröltzsch,1991). Some of these papers present applications to optimal control.

However, there are serious difficulties in a direct application of stability results for abstract mathematical programs to nonlinear optimal control problems. These difficulties are due to the phenomenon of the so called norm discrepancy, which is intrinsicly connected with nonlinear optimal control problems

[^0](cf. Ioffe,1979; Maurer,1981). Namely, the differentiability properties as well as constraint qualifications used in stability analysis are satisfied in a stronger norm (of $L^{\infty}$-type), whereas the second order sufficient optimality condition, which is also needed, holds in a weaker norm (of $L^{2}$-type).

A method to overcome these difficulties was proposed in Malanowski (1993), Malanowski (1994), where a modification of Robinson's generalized implicit function theorem, Robinson (1980), was used. This method is called two-norm approach since it uses both norms at the same time. The crucial role is played here by the regularity of optimal solutions, which has to be additionally studied. Therefore, the method is quite complicated technically.

The purpose of this paper is to present the very essence of the two-norm approach, which is rather simple and to emphasize the relations to finite dimensional mathematical programs.

Accordingly, only the general outlines of the proofs are given, and for technicalities the reader is refered to the original papers. Most of the presented results are contained in the papers by the author himself Malanowski (1992), Malanowski (1993), Malanowski (1994). The material related to second order sufficient optimality conditions can be found in papers by Dontchev, Hager (1993) and Dontchev et al. (1994).

The organization of the paper is the following. In Section 2 stability of solutions to linear-quadratic problems with respect to additive perturbations is analysed. In Section 3 the same class of problems is considered but in the situation of norm discrepancy.

Section 4 is devoted to nonlinear problems. It is recalled how Robinson's implicit function theorem for generalized equations, Robinson (1980), allows to reduce local stability analysis of nonlinear sytems to such an analysis for linearquadratic problems with respect to additive perturbations, which was discussed in Section 2.

Nonlinear optimization problems in case of the norm discrepancy are considered in Section 5. It is shown that, if an additional information on regularity of the solutions is available, then a modification of Robinson's theorem together with the results of Section 3 can be used in stability analysis.

Finally in Section 6 the abstract stability results obtained in Section 5 are applied to control and state constrained optimal control problems for nonlinear ordinary differential equations.

Some used notation:
$X, Y, Z, \ldots$ denotes Banach spaces and $\widehat{X}, \widehat{Y}, \widehat{Z}, \ldots$ Hilbert spaces. Asterisks denote dual spaces. $L(X, Y)$ is the Banach space of linear continuous operators from $X$ into $Y$.
The norms in Banach and Hilbert spaces are denoted by $\mathbb{\square} \cdot \square$ and $\|\cdot\|$, respectively, with a subscript refering to the space.
$(\cdot, \cdot)_{X}$ denotes the duality pairing between $X$ and $X^{*}$.
$\mathbf{R}^{n}$ is the $n$-dimensional Euclidean space with the inner product denoted by $\langle x, y\rangle$ and the norm $|x|=\langle x, x\rangle^{\frac{1}{2}}$.
$L^{2}\left(0, T ; \mathbf{R}^{n}\right)$ is the Hilbert space of square integrable vector functions, with the inner product

$$
(x, y)=\int_{0}^{T}\langle x(t), y(t)\rangle d t
$$

and the norm

$$
\|x\|_{2}=(x, x)^{\frac{1}{2}} .
$$

$M^{2}\left(0, T ; \mathbf{R}^{n}\right)$ denotes the Hilbert space of square integrable vector functions that assume a finite value at $t=0$, with the inner product

$$
(x, y)_{M}=\langle x(0), y(0)\rangle+(x, y) .
$$

$L^{\infty}\left(0, T ; \mathbf{R}^{n}\right)$ is the Banach space of essentially bounded vector functions with the norm

$$
\|x\|_{\infty}=\max _{i} \operatorname{ess} \sup _{t \in[0, T]}\left|x^{i}(t)\right| .
$$

$C^{0}\left(0, T ; \mathbf{R}^{n}\right)$ and $C^{1}\left(0, T ; \mathbf{R}^{n}\right)$ are the spaces of continuous and continuously differentiable vector functions, respectively, equipped with the usual norms.

$$
W^{1, p}\left(0, T ; \mathbf{R}^{n}\right)=\left\{x \in L^{p}\left(0, T ; \mathbf{R}^{n}\right) \mid \dot{x} \in L^{p}\left(0, T ; \mathbf{R}^{n}\right)\right\}, \quad p=2, \infty,
$$

denote the Sobolev spaces of absolutely continuous functions with the norms

$$
\|x\|_{1,2}=\left\{|x(0)|^{2}+\|\dot{x}\|_{2}^{2}\right\}^{\frac{1}{2}}
$$

and

$$
\|x\|_{1, \infty}=\max \left\{|x(0)|,\|\dot{x}\|_{\infty}\right\},
$$

respectively.

## 2. Stability analysis of linear-quadratic problem

In this section we are going to analyse stability of solutions to cone-constrained linear-quadratic problems of optimization in Banach spaces subject to linear perturbations.

Let $Z$ and $Y$ be two Banach spaces, space of arguments and constraints, respectively. Moreover there are given Hilbert spaces $\widehat{Z}$ and $\widehat{Y}$. We assume

$$
\begin{equation*}
Z \subset \widehat{Z}=\widehat{Z}^{*} \subset Z^{*}, \quad Y \subset \widehat{Y}=\widehat{Y}^{*} \subset Y^{*} \tag{2.1}
\end{equation*}
$$

with all embeddings being dense and continuous.
In the space $Y$ there is given a closed cone $K$ with the vertex at the origin, which induces a partial order in $Y$. We denote

$$
K^{+}=\left\{\lambda \in Y^{*} \mid(\lambda, y)_{Y} \geq 0 \text { for all } y \in K\right\} .
$$

Let us denote $X=Z \times Y^{*}, X^{*}=Z^{*} \times Y, \delta=(a, b) \in X^{*}$ and let $\mathcal{O}^{\rho} \subset X^{*}$ be the open ball of radius $\rho$ around the origin.

There are given two linear continuous operators

$$
Q \in L\left(Z, Z^{*}\right), \quad D \in L(Z, Y)
$$

and elements $q \in Z^{*}, e \in Y$.
We consider the family of the following optimization problems depending on the additive parameter $\delta \in X^{*}$ :

$$
\begin{align*}
& \min _{w \in Z}\left\{\frac{1}{2}(Q w, w)_{Z}+(q+a, w)_{Z}\right\} \\
& \text { subject to } \\
& D w+e+b \in K .
\end{align*}
$$

We are going to investigate local stability of solutions to $\left(L_{\delta}\right)$. More precisely, we are looking for sufficient conditions under which there exists $\rho>0$ such that for each $\delta \in \mathcal{O}^{\rho}$ Problem ( $L_{\delta}$ ) has a unique solution $w_{\delta}$, which is a Lipschitz continuous function of $\delta$.

Let us assume for a moment that the following two conditions are satisfied:

$$
\begin{array}{r}
D Z=Y \text {, i.e. } D \in L(Z, Y) \text { is surjective, } \\
\qquad(Q w, w)_{Z} \geq \alpha \rrbracket w \rrbracket_{Z}^{2}, \quad \alpha>0 \tag{2.3}
\end{array}
$$

$$
\text { for all } w \in \Phi:=\{w \in Z \mid D w \in K+(-K)\} \text {. }
$$

If some additional conditions are satisfied (see Theorem II. 10 in Brezis,1983), in particular if $Z$ and $Y$ are Hilbert spaces, then by (2.2) there exists a right invers $T \in L(Y, Z)$ of $D$. For the sake of simplicity we assume, in this section, that $Z=Z^{*}$ and $Y=Y^{*}$ are Hilbert spaces. Let us introduce the change of variables putting:

$$
\begin{equation*}
v=w+T(e+b) . \tag{2.4}
\end{equation*}
$$

In terms of $v$ Problem $\left(L_{\delta}\right)$ takes on the form:

$$
\begin{align*}
& \min _{v \in Z}\left\{\frac{1}{2}(Q v, v)_{Z}+\left(q^{\prime}+a^{\prime}, v\right)_{Z}\right\} \\
& \text { subject to }  \tag{2.5}\\
& D v \in K \text {, }
\end{align*}
$$

where

$$
\begin{equation*}
q^{\prime}=q-Q T e, \quad a^{\prime}=a-Q T b . \tag{2.6}
\end{equation*}
$$

By (2.2) the feasible set of (2.5) is non empty and by (2.3) there exists a unique solution $v_{\delta}$ of $(2.5)$, which is characterized by the following variational inequality:

$$
\left(Q v_{\delta}+q^{\prime}+a^{\prime}, v-v_{\delta}\right)_{z} \geq 0 \quad \text { for all feasible } v
$$

Let $v_{\delta_{1}}, v_{\delta_{2}}$ be the solutions to (2.5) corresponding to the values $\delta_{1}$ and $\delta_{2}$ of the parameter. We have

$$
\begin{aligned}
& \left(Q v_{\delta_{1}}+q^{\prime}+a_{1}^{\prime}, v_{\delta_{2}}-v_{\delta_{1}}\right)_{Z} \geq 0, \\
& \left(Q v_{\delta_{2}}+q^{\prime}+a_{2}^{\prime}, v_{\delta_{1}}-v_{\delta_{2}}\right)_{Z} \geq 0 .
\end{aligned}
$$

Adding these inequalities and using (2.3) we obtain

$$
\begin{aligned}
& \alpha \rrbracket v_{\delta_{1}}-v_{\delta_{2}} \rrbracket z \leq\left(Q\left(v_{\delta_{1}}-v_{\delta_{2}}\right),\left(v_{\delta_{1}}-v_{\delta_{2}}\right)\right)_{Z} \leq \\
& \leq\left(a_{2}^{\prime}-a_{1}^{\prime}, v_{\delta_{1}}-v_{\delta_{2}}\right) z \leq \llbracket a_{1}^{\prime}-a_{2}^{\prime} \rrbracket Z^{*} \llbracket v_{\delta_{1}}-v_{\delta_{2}} \rrbracket z,
\end{aligned}
$$

i.e., in view of (2.6).

$$
\begin{equation*}
\square v_{\delta_{1}}-v_{\delta_{2}} \rrbracket z \leq \beta \rrbracket \delta_{1}-\delta_{2} \rrbracket X . \tag{2.7}
\end{equation*}
$$

By (2.2) there exists a unique Lagrange multiplier $\mu_{\delta}$ such that the following Kuhn-Tucker conditions hold:

$$
\begin{align*}
& Q v_{\delta}+q^{\prime}+a^{\prime}-D^{*} \mu_{\delta}=0, \\
& \left(\mu_{\delta}, D^{*} v_{\delta}\right)_{Y}=0, \quad \mu_{\delta} \in K^{+} . \tag{2.8}
\end{align*}
$$

In view of (2.2), the first equation in (2.8) yields

$$
\mu_{\delta_{1}}-\mu_{\delta_{2}}=T^{*}\left(Q\left(v_{\delta_{1}}-v_{\delta_{2}}\right)+\left(a_{1}^{\prime}-a_{2}^{\prime}\right)\right),
$$

i.e., by (2.6) and (2.7)

$$
\begin{equation*}
\llbracket \mu_{\delta_{1}}-\mu_{\delta_{2}} \rrbracket_{Y} \leq \beta \rrbracket \delta_{1}-\delta_{2} \rrbracket X^{*} \tag{2.9}
\end{equation*}
$$

Thus we have obtained Lipschitz continuity of the solutions to $\left(L_{\delta}\right)$ and of the associated Lagrange multipliers.

The assumptions (2.2) and (2.3) are very strong and we would like to weaken them as much as possible, still preserving the above Lipschitz continuity properties.

Before doing that, let us recall some stability results for quadratic programming problems in finite dimension, i.e., let us assume for a moment that $Z=\mathbf{R}^{n}, Y=\mathbf{R}^{m}$ and $K \subset \mathbf{R}^{m}$ is the non-negative octant. Then problem $\left(L_{\delta}\right)$ takes on the form:

$$
\begin{align*}
& \min _{w \in \mathbf{R}^{n}}\left\{\frac{1}{2}\langle Q w, w\rangle+\langle q+a, w\rangle\right\} \\
& \text { subject to } \\
& \left\langle d^{i}, w\right\rangle+e^{i}+b^{i} \geq 0, \quad i=1,2, \ldots, m .
\end{align*}
$$

Let $w_{0}$ be the solution to ( $M P_{0}$ ). Denote by

$$
I=\left\{i \in\{1,2, \ldots, m\} \mid\left\langle d^{i}, w_{0}\right\rangle+e^{i}=0\right\}
$$

the set of indices of constraints active at $v_{0}$ and by $D_{I}$ the matrix with rows $d^{i}, i \in I$. We say that the condition of linear independence of the gradients of active constraints (LIC) is satisfied if

$$
\begin{equation*}
D_{I} \text { is of full rank. } \tag{2.10}
\end{equation*}
$$

If (2.10) holds, then there exists a unique Lagrange multiplier

$$
\mu_{0}=\left(\mu_{0}^{1}, \mu_{0}^{2}, \ldots, \mu_{0}^{m}\right)
$$

associated with $v_{0}$. Denote $J=\left\{i \in I \mid \mu_{0}^{i}>0\right\}$. We say that the strong second order sufficient optimality condition (SOC) holds if

$$
\begin{align*}
& \langle Q w, w\rangle \geq \alpha|w|^{2}  \tag{2.11}\\
& \text { for all } w \in\left\{w \in \mathbf{R}^{n} \mid\left\langle d^{i}, w\right\rangle=0, \quad i \in J\right\} .
\end{align*}
$$

Let us modify $\left(M P_{\delta}\right)$ as follows:

$$
\begin{aligned}
& (\widetilde{M P})) \quad \min _{v \in \mathbf{R}^{n}}\left\{\frac{1}{2}\langle Q v, v\rangle+\langle q+a, v\rangle\right\} \\
& \text { subject to } \\
& \left\langle d^{i}, v\right\rangle+e^{i}+b^{i}=0, \quad i \in J, \\
& \left\langle d^{i}, v\right\rangle+e^{i}+b^{i} \geq 0, \quad i \in I \backslash J .
\end{aligned}
$$

It ie easy to see that if (LIC) and (SOC) hold then for ( $\widetilde{M P}_{\delta}$ ) conditions (2.2) and (2.3) are satisfied, so there exists a unique solution $v_{\delta}$ of $\left(\widetilde{M P}_{\delta}\right)$ and a unique associated Lagrange multiplier $\nu_{\delta}$. Both $v_{\delta}$ and $\nu_{\delta}$ are Lipschitz continuous functions of $\delta$.

Certainly $v_{0}=w_{0}$ and $\nu_{0}=\mu_{0}$. Moreover, by Lipschitz continuity of $v_{\delta}$ and $\nu_{\delta}$, for $\delta>0$ sufficiently small we have

$$
\begin{array}{ll}
\left\langle d^{i}, v_{\delta}\right\rangle+e^{i}+b^{i}>0 & \text { for } i \notin I, \\
\nu_{\delta}^{i}>0 & \text { for } i \in J . \tag{2.12}
\end{array}
$$

It implies that $v_{\delta}=w_{\delta}$ and $\nu_{\delta}=\mu_{\delta}$, i.e., $w_{\delta}$ and $\mu_{\delta}$ are locally Lipschitz continuous functions of $\delta$.

An idea similar to that presented above will be also used in stability analysis of solutions to infinite dimensional problems $\left(L_{\delta}\right)$.

Let us start with such a reformulation of the definitions of (LIC) and (SOC), which would fit to infinite dimensional situations.

For $y \in K$ and $\mu \in K^{+}$define the subspaces:

$$
\begin{align*}
& M_{y}=(K+[y]) \cap(-K+[y]) \subset Y, \\
& N_{\mu}=\left(K^{+}+[\mu]\right) \cap\left(-K^{+}+[\mu]\right) \subset Y^{*}, \tag{2.13}
\end{align*}
$$

where $[y]=\left\{\kappa y \mid \kappa \in \mathbf{R}^{1}\right\}$, and denote $M_{0}=M_{y_{0}}, N_{0}=N_{\mu_{0}}$, where $y_{0}=D w_{0}+$ $e$, and $N_{0}^{\perp}=\left\{y \in Y \mid(\mu, y)_{Y}=0\right.$ for all $\left.\mu \in N_{0}\right\}$. It can be easily checked
that in case of ( $M P_{0}$ ), $M_{0}$ is the subspace generated by nonzero components $y_{0}^{i}$ of the vector $y_{0}$, while $N_{0}^{\perp}=\left\{y \in Y \mid\left\langle\mu^{i}, y\right\rangle=0, i \in J\right\}$. Hence conditions (2.10) and (2.11) can be rewritten in the form
(LIC) $\quad D Z+M_{0}=Y$,
(SOC) $\quad(Q w, w)_{Z} \geq \alpha \rrbracket w \rrbracket_{Z}^{2} \quad$ for all $w \in\left\{w \in Z \mid D w \in N_{0}^{\perp}\right\}$.
In finite dimensional case conditions (LIC) and (SOC) are stable under small perturbations, in the sense that they are preserved if we substitute $M_{0}$ and $N_{0}$ by $M_{y}$ and $N_{\lambda}$ with $y \in K$ and $\mu \in K^{+}$sufficiently close to $y_{0}$ and $\mu_{0}$, respectively. This fact has been reflected in (2.12). Such a stability is not automatic in infinite dimensional spaces. To assure it we will strengthen conditions (LIC) and (SOC) introducing some "margin of freedome". Namely we assume:
(I 1) There exists a closed subspace $M \subset M_{0}$ and a linear continuous mapping II : Y $\mapsto M$ such that

$$
\begin{equation*}
D Z+\Pi Y=Y \tag{2.14}
\end{equation*}
$$

and moreover there exists a neighborhood $\mathcal{Y}_{0}$ of $y_{0}$ such that

$$
\begin{equation*}
M \subset M_{y} \quad \text { for all } y \in \mathcal{Y}_{0} \cap K \tag{2.15}
\end{equation*}
$$

(I 2) There exists a closed subspace $N \subset N_{0}$ and a constant $\alpha>0$ such that

$$
\begin{equation*}
(Q w, w)_{Z} \geq \alpha \rrbracket w \rrbracket_{Z}^{2} \quad \text { for all } z \in\left\{z \in Z \mid D z \in N^{\perp}\right\} \tag{2.16}
\end{equation*}
$$

and moreover there exists a neighborhood $\Lambda_{0}$ of $\mu_{0}$ such that

$$
\begin{equation*}
\Lambda_{0} \cap\left(K^{+}+N\right) \subset K^{+} . \tag{2.17}
\end{equation*}
$$

The modified problem ( $\widetilde{L}_{\delta}$ ) analogous to ( $\widetilde{M P}_{\delta}$ ) and corresponding to $\left(L_{\delta}\right)$ is defined as follows:
( $\tilde{L}_{\delta}$ )

$$
\begin{aligned}
& \min _{v \in Z}\left\{\frac{1}{2}(Q v, v)_{Z}+(q+a, v)_{Z}\right\} \\
& \text { subject to } \\
& D v+e+b \in \widetilde{K}:=\overline{\left(K \cap N^{\perp}\right)+M} .
\end{aligned}
$$

In order to assure that condition (2.2) is satisfied, we enlarge the space of arguments to $Z \times Y$ and modify $\left(\widetilde{L}_{\delta}\right)$ as follows:

$$
\begin{align*}
& \min _{(v, u) \in Z \times Y}\left\{\frac{1}{2}\left[(Q v, v)_{Z}+(u, u)_{Y}\right]+(q+a, v)_{Z}\right\}  \tag{L}\\
& \text { subject to } \\
& D v+\Pi u+e+b \in \widetilde{K} .
\end{align*}
$$

Due to the additional term in the cost functional the coercivity condition (2.3) is satisfied. Hence by (2.7) and (2.9) the solutions ( $v_{\delta}, u_{\delta}$ ) of ( $\left.\widetilde{L}_{\delta}^{\prime}\right)$ and the associated Lagrange multiplers $\nu_{\delta}$ are Lipschitz continuous functions of $\delta$.

On the other hand, in view of the definitions of $\Pi$ and $\widetilde{K}$ it can be shown (cf. Lemma 3.7 in Malanowski,1992 as well as Lemmata 4.3 and 4.4 in Malanowski,
1993) that a pair $(v, u) \in Z \times Y$ is feasible for $\left(\widetilde{L}_{\delta}^{\prime}\right)$ if and only if $v \in Z$ is feasible for $\left(\widetilde{L}_{\delta}\right)$. Hence, it follows from the form of the cost functional in $\left(\widetilde{L}_{\delta}^{\prime}\right)$ that for any $\delta \in X^{*}$ the solution to $\left(\tilde{L}_{\delta}^{\prime}\right)$ has the form $\left(v_{\delta}, 0\right)$, where $v_{\delta}$ is the solution to $\left(\tilde{L}_{\delta}\right)$. Certainly also the Lagrange multiplier $\nu_{\delta}$ is the same for $\left(\widetilde{L}_{\delta}^{\prime}\right)$ and $\left(\tilde{L}_{\delta}\right)$.

At the reference point $\delta=0$ we have

$$
v_{0}=w_{0}, \quad \nu_{0}=\mu_{0}
$$

where $w_{0}$ is the solution to $\left(L_{0}\right)$ and $\mu_{0}$ - the associated Lagrange multiplier. It is shown in Malanowski (1993) (cf. Lemmata 4.3 and 4.4) that for $\delta$ sufficiently small this property is preserved, i.e., we have

$$
\begin{equation*}
v_{\delta}=w_{\delta}, \quad \nu_{\delta}=w_{\delta} \tag{2.18}
\end{equation*}
$$

Hence we finally obtain:
THEOREM 2.1 If $Z$ and $Y$ are Hilbert spaces and assumptions (I 1), (I 2) are satisfied, then there exists a constant $\rho>0$ such that for all $\delta \in \mathcal{O}^{\rho}$ there is a unique solution $w_{\delta}$ of $\left(L_{\delta}\right)$ and a unique associated Lagrange multiplier $\mu_{\delta}$, and both $w_{\delta}$ and $\mu_{\delta}$ are Lipschitz continuous functions of $\delta$.

In order to illustrate the nature of assumptions (I 1) and (I 2) as well as the reason why for $\delta$ sufficiently small the solutions to $\left(L_{\delta}\right)$ and $\left(\widetilde{L}_{\delta}\right)$ coincide, let us consider the following simple example.

Example 2.2 Let $Z=Y=L^{\infty}(0,1), K \subset Y$ be the cone of non-negative functions,

$$
(Q w, w)_{Z}=\int_{0}^{1} g(t) w^{2}(t) d t, \quad(D w)(t)=d(t) w(t), \quad \text { where } g, d \in C^{0}(0,1)
$$

Assume that $w_{0}$ and $\mu_{0}$ are given by functions of class $C^{0}(0,1)$.
For $\epsilon \geq 0$ and $\eta \geq 0$ we define the sets (see Figure 1):

$$
\begin{aligned}
\Omega^{\epsilon} & =\left\{t \in[0,1] \mid y_{0}(t):=D w_{0}(t)+e(t)>\epsilon\right\}, \\
\Xi^{\eta} & =\left\{t \in[0,1] \mid \mu_{0}(t)>\eta\right\} .
\end{aligned}
$$

Let us introduce the following subspaces:

$$
\begin{array}{ll}
M^{\epsilon}=\left\{y \in L^{\infty}(0, T) \mid y(t)=0\right. & \text { for all } \left.t \in[0,1] \backslash \Omega^{\epsilon}\right\}, \\
N^{\eta}=\left\{\mu \in L^{\infty}(0, T) \mid \mu(t)=0\right. & \text { for all } \left.t \in[0,1] \backslash \Xi^{\eta}\right\} .
\end{array}
$$

It is easy to see that for any $\epsilon>0$ and $\eta>0$ we have $M^{\epsilon} \subset M_{0}$ and $N^{\eta} \subset N_{0}$.
Let us choose any $\epsilon>0$ and $\eta>0$ and define $M=M^{\epsilon}, N=N^{\eta}$. Note that $M$ and $N$ are closed in $L^{\infty}(0,1)$, but $N$ is not closed in the space $Y^{*}=$ $\left(L^{\infty}(0,1)\right)^{*}$ - this fact will be discussed later on.

We define

$$
\Pi y(t)= \begin{cases}y(t) & \text { if } t \in M^{\epsilon} \\ 0 & \text { if } t \notin M^{\epsilon}\end{cases}
$$

It is easy to see that $(2.15)$ is satisfied if we choose

$$
\mathcal{Y}_{0}=\left\{y \in L^{\infty}(0,1)| | y(t)-y_{0}(t) \left\lvert\,<\frac{\epsilon}{2} \quad\right. \text { for all } t \in[0,1]\right\}
$$

while (2.14) holds if $d(t) \neq 0$ for all $t \in[0,1] \backslash \Omega^{2 \epsilon}$.
Coercivity condition (2.16) is never satisfied in the norm of the original space $L^{\infty}(0,1)$. However, it is satisfied in the weaker norm of the space $L^{2}(0,1)$, provided that $g(t)>0$ for all $t \in[0,1] \backslash \Xi^{\eta}$.


Figure 1

Constraints in $\left(\tilde{L}_{\delta}\right)$ have the form:

$$
d(t) v(t)+e(t)+b(t)\left\{\begin{array}{lll}
=0 & \text { for } t \in \Xi^{\eta}  \tag{2.19}\\
\geq 0 & \text { for } t \in\left([0,1] \backslash \Xi^{\eta}\right) \backslash \Omega^{\epsilon} \\
\text { arbitrary } & \text { for } \quad t \in \Omega^{\epsilon}
\end{array}\right.
$$

To illustrate the reason why for $\delta$ sufficiently small the solutions to $\left(L_{\delta}\right)$ and $\left(\widetilde{L}_{\delta}\right)$ coincide, let us note that since $v_{\delta} \rightarrow v_{0}$, for $\delta=(a, b)$ sufficiently small

$$
d(t) v_{\delta}(t)+e(t)+b(t)=d(t) v_{0}(t)+e(t)+b(t)+d(t)\left(v_{\delta}(t)-v_{0}(t)\right)>0
$$

for $t \in \Omega^{\epsilon}$. Hence constraints (2.19) can be changed as follows:

$$
d(t) v(t)+e(t)+b(t)\left\{\begin{array}{lll}
=0 & \text { for } & t \in \Xi^{\eta}  \tag{2.20}\\
\geq 0 & \text { for } & t \in[0,1] \backslash \Xi^{\eta} .
\end{array}\right.
$$

Certainly this mechanism has worked because we assured a "margin of freedom" choosing $\epsilon>0$ in the definiton of $M$. For $\epsilon=0$ it would not work anymore.

Assume for a moment that (2.17) holds in the sense of $L^{\infty}(0,1)$-topology, i.e., $\Lambda_{0}$ is a $L^{\infty}$-neighborhood of $\mu_{0}$ and $K^{+}=\left\{\mu \in L^{\infty}(0,1) \mid \mu(t) \geq 0\right\}$. If $\nu_{\delta} \rightarrow \mu_{0}$ in $L^{\infty}(0,1)$ (which, does not follow from (2.9)!) then, for $\delta$ sufficiently small we would have $\nu_{\delta}>0$ for $t \in \Xi^{\eta}$. This means that on the set $\Xi^{\eta}$ the equality type constraints could be substituted by inequality constraints and (2.20) could be rewritten as

$$
\begin{equation*}
d(t) v(t)+e(t)+b(t) \geq 0 \quad \text { for } t \in[0,1] \tag{2.21}
\end{equation*}
$$

which shows that for $\delta$ sufficientlly small problems $\left(L_{\delta}\right)$ and $\left(\widetilde{L}_{\delta}\right)$ coincide.

## 3. Linear-quadratic problem with norm discrepancy

Example 2.2 illustrates the fundamental role of the topology in which Problem $\left(L_{\delta}\right)$ is considered. This role can be seen in several points:
(i) Stability of constraint qualifications (2.14) is satisfied in $L^{\infty}(0,1)$, but it would fail in $L^{2}(0,1)$.
(ii) Coercivity condition (2.16) never holds in $L^{\infty}(0,1)$, but it could be satisfied in $L^{2}(0,1)$.
(iii) Condition (2.17) (stability of coercivity) is never satisfied in the space $\left(L^{\infty}(0,1)\right)^{*}$, but it could be satisfied in $L^{\infty}(0,1)$.
The phenomenon that constraint qualifications are satisfied in a stronger topology and coercivity condition in a weaker one is called the norm discrepancy and it is typical for optimal control problems (cf. Ioffe,1979; Maurer,1981). The presence of the norm discrepancy creates a serious difficulty in stability analysis of optimal control problems, especially for nonlinear systems. We will discuss it in details in this and next sections.

To cope with this difficulty we will work with two norm simultaneously, using the so called two-norm approach, Malanowski (1993), Malanowski (1994), in which regularity of the solutions plays a crucial role.

In this section we will apply this approach to linear-quadratic problems. To do so we have to reformulate Problem ( $L_{\delta}$ ).
Let us denote by $\widehat{K}$ the closure of $K$ in $\widehat{Y}$. Furthermore we assume that

$$
\begin{equation*}
Q \in L(Z, Z) \cap L(\widehat{Z}, \widehat{Z}), \quad D \in L(Z, Y) \cap L(\widehat{Z}, \widehat{Y}), \quad q \in Z, \quad e \in Y \tag{3.1}
\end{equation*}
$$

For $\delta \in \widehat{X}=\widehat{Z} \times \widehat{Y}$ we consider along $\left(L_{\delta}\right)$ the following problem:

$$
\begin{aligned}
\left(\widehat{L}_{\delta}\right) \quad & \min _{w \in \widehat{Z}}\left\{\frac{1}{2}(Q w, w)_{\widehat{Z}}+(q+a, w)_{\widehat{Z}}\right\} \\
& \text { subject to } \\
& D w+e+b \in \widehat{K},
\end{aligned}
$$

which is identical with $\left(L_{\delta}\right)$ except that it is considered in $\widehat{Z}, \widehat{Y}$ rather than in $Z, Y$ spaces. Certainly each solution to $\left(\widehat{L}_{\delta}\right)$ is also a solution to $\left(L_{\delta}\right)$, provided that $w_{\delta} \in Z$.

Assume that there exists a unique solution $w_{0} \in Z$ of $\left(\widehat{L}_{0}\right)$, i.e., also of $\left(L_{0}\right)$, and the following conditions analogous to (I 1) hold:
(II 1) There exists a subspace $M \subset M_{0}$ closed in $Z$, and a linear continuous mapping

$$
\Pi \in L(Y, Y) \cap L(\widehat{Y}, \widehat{Y}), \quad \Pi: Y \rightarrow M
$$

such that

$$
\begin{align*}
& D Z+\Pi Y=Y,  \tag{3.2}\\
& D \widehat{Z}+\Pi \widehat{Y}=\widehat{Y} . \tag{3.3}
\end{align*}
$$

Moreover there exists a $Y$-neighborhood $\mathcal{Y}_{0}$ of $y_{0}$ such that $M \subset M_{y}$ for all $y \in \mathcal{Y}_{0} \cap K$.
By (3.2) there exists a unique Lagrange multiplier $\mu_{0} \in Y^{*}$ associated with $w_{0}$. We assume that $\mu_{0}$ is more regular namely (II 2) $\mu_{0} \in Y$.

Instead of (I 1) we introduce the following condition, which reflects the norm discrepancy:
(II 3) There exists a subspace $N \subset N_{0} \subset Y^{*}$ and a constant $\alpha>0$ such that

$$
\begin{equation*}
(Q w, w)_{\widehat{z}} \geq \alpha\|w\|_{\widehat{Z}}^{2} \quad \text { for all } w \in\left\{w \in \widehat{Z} \mid D w \in N^{\perp}\right\} \tag{3.5}
\end{equation*}
$$

and moreover there exists a $Y$-neighborhood $\Lambda_{0}$ of $\mu_{0}$ such that

$$
\begin{equation*}
\Lambda_{0} \cap\left(K^{+}+N\right) \subset K^{+} \tag{3.6}
\end{equation*}
$$

By the same reasoning as in Section 2 we can find that the solution $w_{\delta}$ of $\left(\widehat{L}_{\delta}\right)$ and the associated Lagrange multipliers $\mu_{\delta}$ are Lipschitz continuous in the sense of $\widehat{Z}$-norm, and $\widehat{Y}$-norm, repectively:

$$
\begin{equation*}
\left\|w_{\delta_{1}}-w_{\delta_{2}}\right\|_{\widehat{Z}},\left\|\mu_{\delta_{1}}-\mu_{\delta_{2}}\right\|_{\widehat{Y}} \leq c\left\|\delta_{1}-\delta_{2}\right\|_{\widehat{X}}, \tag{3.7}
\end{equation*}
$$

provided that

$$
\begin{equation*}
w_{\delta_{i}} \in \mathcal{Y}_{0} \quad \text { and } \mu_{\delta_{i}} \in \Lambda_{0} \quad i=1,2 \tag{3.8}
\end{equation*}
$$

Certainly (3.7) does not assure that (3.8) is satisfied for all variations $\delta$ from any arbitrary small ball $\widehat{\mathcal{O}}^{\rho}$ in $\widehat{X}$. In order to get (3.8) we should assure that convergence in the weaker norm $\widehat{X}$ implies convergence in the stronger norm $X$. This would hold if $\left(w_{\delta}, \mu_{\delta}\right)$ belong to a certain set $\Gamma$ compact in $X$. We can not
expect it for all variations $\delta$ from any ball in $\widehat{X}$, because it would mean that the mapping

$$
\begin{equation*}
\mathcal{P}: \widehat{X} \mapsto \widehat{X}, \quad \mathcal{P}(\delta)=\left(w_{\delta}, \mu_{\delta}\right) \tag{3.9}
\end{equation*}
$$

is compact, which is not true.
Therefore, to assure that $\left(w_{\delta}, \mu_{\delta}\right)$ belong to a compact set we must restrict ourselves to a certain set $\Delta$ of more regular variations. We assume that such a set exists, namely:
(II 4) There exists a closed convex set $\Delta \subset \widehat{X}$ and a convex set $\tilde{\Gamma}$ compact in $X$ such that for any $\delta \in \Delta,\left(w_{\delta}, \mu_{\delta}\right) \in \widetilde{\Gamma}$.

From the above considerations we obtain:
Theorem 3.1 If assumptions (II 1)-(II 4) hold then there exist constants $\rho>0$ and $c>0$ such that for all $\delta \in \Delta \cap \widehat{\mathcal{O}}^{\rho}$ there exists a unique solution $w_{\delta}$ of $\left(L_{\delta}\right)$ and a unique associated Lagrange multiplier $\mu_{\delta}$ such that

$$
\begin{equation*}
\left\|w_{\delta_{1}}-w_{\delta_{2}}\right\|_{\widehat{Z}},\left\|\mu_{\delta_{1}}-\mu_{\delta_{2}}\right\|_{\widehat{Y}} \leq c\left\|\delta_{1}-\delta_{2}\right\|_{\widehat{X}} \tag{3.10}
\end{equation*}
$$

## 4. Stability analysis for nonlinear problems

This section is devoted to stability analysis of nonlinear cone-constrained optimization problems. In this analysis the stability results for linear-quadratic problems presented in Section 2 will be applied.

The same notation as in the previous section will be used. In addition, let us introduce a Banach space $H$, called the space of parameters, and let $G \subset H$ be an open set of admissible parameters. On $Z \times G$ there are defined two functions

$$
F: Z \times G \mapsto \mathbf{R}^{1}, \quad \varphi: Z \times G \mapsto Y
$$

We consider the family of the following optimization problems depending on the parameter $h \in G$ :

$$
\begin{align*}
& \min _{z \in Z} F(z, h)  \tag{h}\\
& \text { subject to } \\
& \varphi(z, h) \in K \subset Y
\end{align*}
$$

Let $\bar{h}$ be a fixed reference value of the parameter.
We assume:
(III 1) There exists a (local) solution $z_{\bar{h}}$ of $\left(P_{\bar{h}}\right)$,
(III 2) $F(\cdot, \cdot)$ and $\varphi(\cdot, \cdot)$ are two times Fréchet differentiable,
(III 3) The point $z_{\tilde{h}_{-}}$is regular in the sense of Robinson, i.e.,

$$
D_{z} \varphi\left(z_{\bar{h}}, \bar{h}\right) Z-K+\left[\varphi\left(z_{\bar{h}}, \bar{h}\right)\right]=Y
$$

Let us introduce the following Lagrangian associated with $\left(P_{h}\right)$ :

$$
\mathcal{L}(\cdot, \cdot, \cdot): Z \times Y^{*} \times G \mapsto \mathbf{R}^{1}, \quad \mathcal{L}(z, \lambda, h)=F(z, h)-(\lambda, \varphi(z, h))_{Y}
$$

By (III 3) there exists a Lagrange multiplier $\lambda_{\bar{h}}$ associated with $z_{\bar{h}}$ such that the following Kuhn-Tucker conditions hold:

$$
\begin{aligned}
& D_{z} \mathcal{L}\left(z_{\bar{h}}, \lambda_{\bar{h}}, \bar{h}\right)=D_{z} F\left(z_{\bar{h}}, \bar{h}\right)-D_{z} \varphi^{*}\left(z_{\bar{h}}, \bar{h}\right) \lambda_{\bar{h}}=0, \\
& \left(\lambda_{\bar{h}}, \varphi\left(z_{\bar{h}}, \bar{h}\right)\right)_{Y}=0, \quad \lambda_{\bar{h}} \in K^{+} \subset Y^{*}
\end{aligned}
$$

We are going to investigate the existence, local uniquness and stability with respect to $h$ of Kuhn-Tucker points $x_{h}=\left(z_{h}, \lambda_{h}\right) \in X$ of $\left(P_{h}\right)$, i.e., of points which satisfy conditions:

$$
\begin{align*}
& D_{z} \mathcal{L}\left(z_{h}, \lambda_{h}, h\right)=D_{z} F\left(z_{h}, h\right)-D_{z} \varphi^{*}\left(z_{h}, h\right) \lambda_{h}=0  \tag{4.1}\\
& \left(\lambda_{h}, \varphi\left(z_{h}, h\right)\right)_{Y}=0, \quad \lambda_{h} \in K^{+} \subset Y^{*}
\end{align*}
$$

To this end we will use Robinson's implicit function theorem for generalized equations, Robinson (1980). To apply this theorem we need (cf. Malanowski, 1992) the following linear-quadratic approximation to $\left(P_{\bar{h}}\right)$, perturbed by the parameter $\delta=(a, b) \in X^{*}$ :

$$
\begin{aligned}
\left(L P_{\delta}\right) \quad & \min _{w \in Z}\left\{\frac{1}{2}(Q w, w)_{Z}+(q+a, w)_{Z}\right\} \\
& \operatorname{subject~to~} \\
& D w+e+b \in K
\end{aligned}
$$

where

$$
\begin{array}{ll}
Q=D_{z z}^{2} \mathcal{L}\left(z_{\bar{h}}, \lambda_{\bar{h}}, \bar{h}\right), & D=D_{z} \varphi\left(z_{\bar{h}}, \bar{h}\right) \\
q=-Q z_{\bar{h}}+D_{z} F\left(z_{\bar{h}}, \bar{h}\right), & e=\varphi\left(z_{\bar{h}}, \bar{h}\right)-D z_{\bar{h}} \tag{4.2}
\end{array}
$$

Following Robinson (1980), we call $\left(P_{\bar{h}}\right)$ strongly regular at the point $\left(z_{\bar{h}}, \lambda_{\bar{h}}\right)$ if:
(A) There exists $\rho>0$ such that for all $\delta \in \mathcal{O}^{\rho}$ there is a unique pair $\left(w_{\delta}, \mu_{\delta}\right) \in Z \times Y^{*}$, where $w_{\delta}$ is the solution to $\left(L P_{\delta}\right)$ and $\mu_{\delta}$ - the associated Lagrange multiplier and $\left(w_{\delta}, \mu_{\delta}\right)$ is a Lipschitz continuous function of $\delta$.
The following fundamental theorem is due to Robinson (1980):
THEOREM 4.1 Suppose that (A) is satisfied, then there exist a neighborhood $G_{\bar{h}} \subset G$ of $\bar{h}$ and a neighborhood $\mathcal{X}_{\bar{h}} \subset X$ of $x_{\bar{h}}$ such that for any $h \in G_{\bar{h}}$, (4.1) has a solution $x_{h}$ unique in $\mathcal{X}_{\bar{h}}$, which is a Lipschitz continuous function of $h$.

Theorem 4.1 allows us to reduce stability analysis of nonlinear problem $\left(P_{h}\right)$ to such an analysis of the linearized problem $\left(L P_{\delta}\right)$. Cerainly $\left(L P_{\delta}\right)$ coincides with

Problem $\left(L_{\delta}\right)$ considered in Section 2. By Theorem 2.1 we find that condition (A) is satisfied if assumptions (I 1) and (I 2) hold.

It is easy to see that $w_{0}=z_{\bar{h}}, \mu_{0}=\lambda_{\bar{h}}, y_{0}=D w_{0}+e=\varphi\left(z_{\bar{h}}, \bar{h}\right)$. Hence, by (2.13) the subspaces $M_{0} \subset Y$ and $N_{0} \subset Y^{*}$ are given by

$$
\begin{align*}
& M_{0}=\left(K+\left[\varphi\left(z_{\bar{h}}, \bar{h}\right)\right]\right) \cap\left(-K+\left[\varphi\left(z_{\bar{h}}, \bar{h}\right)\right]\right), \\
& N_{0}=\left(K^{+}+\left[\lambda_{\bar{h}}\right]\right) \cap\left(-K^{+}+\left[\lambda_{\bar{h}}\right]\right) . \tag{4.3}
\end{align*}
$$

We assume:
(III 4) Condition (I 1) holds with $M_{0}$ given in (4.3),
(III 5) Condition (I 2) holds with $N_{0}$ given in (4.3).
It follows from Theorems 2.1 and 4.1 that if (III 1)-(III 5) hold, then for $h \in G_{\bar{h}}$ the solutions to (4.1) exist and are Lipschitz continuous functions of $h$.

On the other hand, it follows from Lemma 8 in Dontchev, Hager (1993) (cf. also Lemma 1 in Dunn,Tian,1992) that (III 5) together with (III 4) constitutes a (strong) second order sufficient optimality condition. In view of Lipschitz continuity of $\left(z_{h}, \lambda_{h}\right)$, using the same argument as in the proof of Lemma 5.3 in Malanowski,1993, we find that conditions (III 4) and (III 5) are satisfied at $\left(z_{h}, \lambda_{h}, h\right)$ for $h$ sufficiently close to $\bar{h}$. This means that $z_{h}$ is a (local) solution to $\left(P_{h}\right)$ and $\lambda_{h}$ - the associated Lagrange multiplier. Thus we obtain:

Theorem 4.2 If $Z$ and $Y$ are Hilbert spaces and assumptions (III 1)-(III 5) hold, then there exist a neighborhood $G_{\bar{h}} \subset G$ of $\bar{h}$ and a neighborhood $\mathcal{Z}_{\bar{h}} \subset Z$ of $z_{\bar{h}}$ such that for each $h \in G_{\bar{h}}$ there is a unique in $\mathcal{Z}_{\bar{h}}$ solution $z_{h}$ to $\left(P_{h}\right)$ and a unique associated Lagrange multiplier $\lambda_{h}$.

Moreover, $z_{h}$ and $\lambda_{h}$ are Lipschitz continuous functions of $h$.

## 5. Nonlinear problems with norm discrepancy

The notion of norm discrepancy was introduced in Section 3. Now we are going to discuss stability of solutions to optimization problems for nonlinear systems in presence of the norm discrepancy. To illustrate better the nature of the norm discrepancy we start with the following classical example.
Example 5.1 In (2.1) let us choose $Z=L^{\infty}(0,1)$ and $\widehat{Z}=L^{2}(0,1)$. Consider the integral functional

$$
F(z)=\int_{0}^{1}\left(z^{2}(t)-1\right)^{2} d t
$$

It is easy to see that $F$ is well defined and differentiable on $Z$, but not on $\widehat{Z}$. However for $z \in Z$

$$
\left(D_{z z}^{2} F(z) y, y\right)=4 \int_{0}^{1}\left(3 z^{2}(t)-1\right) y^{2}(t) d t
$$

is a quadratic form well defined and continuous on $\widehat{Z}$.
Any function $z$, such that $|z(\cdot)| \equiv 1$ is a global minimizer of $F$ on $Z$. Let us take as the minimizer $\tilde{z}(\cdot) \equiv 1$. We have

$$
\begin{equation*}
\left(D_{z z}^{2} F(\widetilde{z}) y, y\right)=4 \int_{0}^{1}\left(3 \widetilde{z}^{2}(t)-1\right) y^{2}(t) d t=8\|y\|_{2}^{2} \tag{5.1}
\end{equation*}
$$

however does not exist $c>0$ such that

$$
\left(D_{z z}^{2} F(\tilde{z}) y, y\right) \geq c\|y\|_{\infty}^{2} \quad \text { for all } y \in Z=L^{\infty}(0, T)
$$

Just this phenomenon is called the norm discrepancy. Moreover, note that if for any $\epsilon \in(0,1)$ we put

$$
z_{\epsilon}=\left\{\begin{array}{ll}
\frac{1}{\sqrt{3}} & \text { for } t \in[0, \epsilon) \\
1 & \text { for } t \in[\epsilon, 1]
\end{array}, \quad y_{\epsilon}= \begin{cases}1 & \text { for } t \in[0, \epsilon) \\
0 & \text { for } t \in[\epsilon, 1]\end{cases}\right.
$$

then we have $\left\|\widetilde{z}-z_{\epsilon}\right\|_{2} \rightarrow 0$ for $\epsilon \rightarrow 0$ and

$$
\left(D_{z z}^{2} F\left(z_{\epsilon}\right) y_{\epsilon}, y_{\epsilon}\right)=0
$$

which shows that coercitivity condition (5.1) is not stable under small perturbations in $L^{2}(0,1)$, but it is stable under small perturbations in $L^{\infty}(0,1)$.

We are going to analyse stability of solutions to $\left(P_{h}\right)$ in the situation of the norm discrepancy using the two-norm approach introduced in Section 3. In a similar way as in Section 4 we assume:
(IV 1) There exists a (local) solution $z_{\bar{h}} \in Z$ of $\left(P_{\bar{h}}\right)$,
(IV 2) $F(\cdot, \cdot)$ and $\varphi(\cdot, \cdot)$ are two times Fréchet differentiable on $Z \times G$.

Moreover we introduce the following conditions:
(IV 3) For each $h \in G, z \in Z$ and $\lambda \in Y$ the following compatibility conditions hold:

$$
\begin{align*}
& D_{z} F(z, h) \in \widehat{Z} \text {, } \\
& D_{z z}^{2} F(z, h) \in L(\widehat{Z}, \widehat{Z}) \text {, } \\
& D_{z} \varphi(z, h) \in L(\widehat{Z}, \widehat{Y}),  \tag{5.2}\\
& D_{z} \varphi^{*}(z, h) \lambda \in \widehat{Z} \text {, } \\
& D_{z z}^{2} \varphi^{*}(z, h) \lambda \in L(\widehat{Z}, \widehat{Z}) \text {. } \\
& \lim \left\|D_{z z}^{2} F\left(z_{1}, h_{1}\right)-D_{z z}^{2} F\left(z_{2}, h_{2}\right)\right\|_{L(\widehat{Z}, \widehat{Z})}=0, \\
& \lim \left\|D_{z} \varphi\left(z_{1}, h_{1}\right)-D_{z} \varphi\left(z_{2}, h_{2}\right)\right\|_{L(\widehat{Z}, \widehat{Y})}=0 \text {, }  \tag{5.3}\\
& \lim \left\|D_{z z}^{2} \varphi^{*}\left(z_{1}, h_{1}\right) \lambda_{1}-D_{z z}^{2} \varphi^{*}\left(z_{2}, h_{2}\right) \lambda_{2}\right\|_{L(\widehat{Z}, \widehat{Z})}=0, \\
& \text { for } \llbracket z_{1}-z_{2} \rrbracket z \rightarrow 0, \rrbracket \lambda_{1}-\lambda_{2} \rrbracket_{Y} \rightarrow 0, \llbracket h_{1}-h_{2} \rrbracket_{H} \rightarrow 0 \text {. }
\end{align*}
$$

Moreover
(IV 4) Let $M_{0}$ be given in (4.3). There exists a subspace $M \subset M_{0}$ closed in $Z$, and a linear continuous mapping

$$
\Pi \in L(Y, Y) \cap L(\widehat{Y}, \widehat{Y}), \quad \Pi: Y \mapsto M,
$$

such that

$$
\begin{align*}
& D_{z} \varphi\left(z_{\bar{h}}, \bar{h}\right) Z+\Pi Y=Y, \\
& D_{z} \varphi\left(z_{\bar{h}}, \bar{h}\right) \widehat{Z}+\Pi \widehat{Y}=\widehat{Y} \tag{5.4}
\end{align*}
$$

Moreover there exists a $Y$-neighborhood $\mathcal{Y}_{0}$ of $y_{\bar{h}}=\varphi\left(z_{\bar{h}}, \bar{h}\right)$ such that

$$
\begin{equation*}
M \subset M_{y} \quad \text { for all } y \in \mathcal{Y}_{0} \cap K \tag{5.5}
\end{equation*}
$$

By (IV 4) there exists a unique Lagrange multiplier $\lambda_{\bar{h}} \in Y^{*}$ associated with $z_{\bar{h}}$. We assume that $\lambda_{\bar{h}}$ is more regular, namely:
(IV 5) $\lambda_{\bar{h}} \in Y$.
(IV 6) There exist a subspace $N \subset N_{0} \subset Y^{*}$ and a constant $\alpha>0$ such that

$$
\begin{equation*}
\left(D_{z z}^{2} \mathcal{L}\left(z_{\bar{h}}, \lambda_{\bar{h}}, \bar{h}\right) w, w\right)_{\widehat{z}} \geq \alpha\|w\|_{\widehat{Z}}^{2} \tag{5.6}
\end{equation*}
$$

for all $w \in\left\{w \in \widehat{Z} \mid D_{z} \varphi\left(z_{\bar{h}}, \bar{h}\right) w \in N^{\perp}\right\}$
and moreover there exists a $Y$-neighborhood $\Lambda_{0}$ of $\lambda_{\bar{h}}$ such that

$$
\begin{equation*}
\Lambda_{0} \cap\left(K^{+}+N\right) \subset K^{+} . \tag{5.7}
\end{equation*}
$$

For $\delta=(a, b) \in \widehat{X}$ we consider the following linear-quadratic problem ( $\left.\widehat{L P}_{\delta}\right)$ associated with ( $P_{\bar{h}}$ ) and analogous to ( $\widehat{L}_{\delta}$ ):

$$
\begin{align*}
& \min _{w \in \widehat{Z}}\left\{\frac{1}{2}(Q w, w)_{\widehat{Z}}+(q+a, \dot{w})_{\widehat{Z}}\right\}  \tag{LP}\\
& \text { subject to } \\
& D w+e+b \in \widehat{K}
\end{align*}
$$

In order to be able to use Theorem 3.1 we still have to assume condition (II 4), which takes on the form:
(IV 7) There exists a closed convex set $\Delta \subset \widehat{X}$ and a convex set $\tilde{\Gamma}$ compact in $X$ such that for any $\delta \in \Delta,\left(w_{\delta}, \mu_{\delta}\right) \in \tilde{\Gamma}$, where $w_{\delta}$ is the solution to $\left(\widehat{L P}_{\delta}\right)$ and $\mu_{\delta}$ - the associated Lagrange multiplier.
Under the above conditions, (3.10) holds for ( $\widehat{L P}_{\delta}$ ). However, it is not enough to allow application of Robinson's theorem 4.1. To see it let us briefly recall the idea of the proof of that theorem.

Define the following mappings:

$$
\begin{aligned}
& \mathcal{P}: \mathcal{O}^{\rho} \mapsto X, \quad \mathcal{P}(\delta):=y_{\delta}:=\left(w_{\delta}, \mu_{\delta}\right), \\
& g(\cdot, \cdot): X \times G \mapsto X^{*}, \quad \ell(\cdot, \cdot): X \times G \mapsto X^{*} \\
& g(x, h):=\binom{D_{z} \mathcal{L}(z, \lambda, h)}{\varphi(z, h)}, \\
& \ell(x, h):=g\left(x_{\bar{h}}, \bar{h}\right)+D_{x} g\left(x_{\bar{h}}, \bar{h}\right)\left(x-x_{\bar{h}}\right)-g(x, h) .
\end{aligned}
$$

By (A) the mapping $\mathcal{P}$ is well defined and Lipschitz continu $י$ is. On the other hand, by continuity of $\ell(\cdot, \cdot)$ there exist $\sigma>0$ and $G_{\bar{h}}$ such that for all $h \in$ $G_{\bar{h}}, \ell(\cdot, h)$ maps $\mathcal{B}_{\bar{h}}^{\sigma}:=\left\{x \in X \mid \rrbracket x-x_{\bar{h}} \rrbracket_{X} \leq \sigma\right\}$ into $\mathcal{O}^{\rho}$. Hence for all $h \in G_{\bar{h}}$ there is the well defined mapping

$$
\Phi_{h}(\cdot): \mathcal{B}_{\bar{h}}^{\sigma} \mapsto X, \quad \Phi_{h}(x):=\mathcal{P}(\ell(x, h)),
$$

as it is illustrated in Figure 2.


Figure 2

It can be checked that $x=\Phi_{h}(x)$ if and only if the Kuhn-Tucker conditions (4.1) are satisfied, i.e., solving (4.1) is equivalent to finding a fixed point of $\Phi_{h}$. It is shown in Robinson (1980), that if (A) is satisfied, then for sufficiently small $\sigma>0$ and $G_{\bar{h}}, \Phi_{h}$ is a contraction self-map on $\mathcal{B}_{\bar{h}}^{\sigma}$ for all $h \in G_{\bar{h}}$.

By the contraction principle there exists a unique fixed point $x_{h}$ of $\Phi_{h}$. Lipschitz continuity of $x_{h}$ follows from the contraction estimates.

In the situation of the norm discrepancy, where instead of (A) we have (3.10), $\Phi_{h}$ is no longer a self-map on any closed ball either in $X$ or in $\widehat{X}$.

Rather than on a ball we can try to use the fixed point argument on the set $\widetilde{\Gamma}$ defined in (IV 7).. However, to be able to do so we still have to assure that $\ell(x, h) \in \Delta$ for all $x \in \widetilde{\Gamma}$ and $h \in G_{\bar{h}}$ (cf. Fig. 2).

We assume:
(IV 8) There exists $\sigma>0$ and $G_{\bar{h}}$ such that

$$
\begin{equation*}
\ell(x, h) \in \Delta \quad \text { for all } x \in \Gamma \text { and } h \in G_{\bar{h}}, \tag{5.8}
\end{equation*}
$$

where $\Gamma=\widetilde{\Gamma} \cap \widehat{\mathcal{B}}_{\bar{h}}^{\sigma}$, with $\widehat{\mathcal{B}}_{\bar{h}}^{\sigma}:=\left\{x \in \widehat{X} \mid\left\|x-x_{\bar{h}}\right\|_{\widehat{X}} \leq \sigma\right\}$.
By (IV 7) and (IV 8) we can repeat the fixed point argument on the set $\Gamma$ and we obtain (cf, Corollary 5.2 in Malanowski,1993):

Theorem 5.2 If (IV 1)-(IV 8) are satisfied then there exist a neighborhood $G_{\bar{h}} \subset G$ of $\bar{h}$ and constants $\sigma>0$ and $c>0$ such that for any $h \in G_{\bar{h}}$ there exists a solution $x_{h}=\left(z_{h}, \lambda_{h}\right)$ of (4.1), unique in $\Gamma=\widetilde{\Gamma} \cap \widehat{\mathcal{B}}_{\bar{h}}^{\sigma}$ and

$$
\begin{equation*}
\left\|z_{h_{1}}-z_{h_{2}}\right\|_{\widehat{Z}},\left\|\lambda_{h_{1}}-\lambda_{h_{2}}\right\|_{\widehat{Y}} \leq c \llbracket h_{1}-h_{2} \rrbracket_{H} \tag{5.9}
\end{equation*}
$$

for all $h_{1}, h_{2} \in G_{\bar{h}}$.
To show that $z_{h}$ and $\lambda_{h}$ are actually the solutions and the Lagrange multipliers of $\left(P_{h}\right)$ we have to show that a sufficient optimality condition is satisfied at $z_{h}$.

First of all note that by (5.9) and by compactness of the set $\Gamma \subset X$ we have:

$$
\lim _{h \rightarrow \bar{h}} \rrbracket z_{h}-z_{\bar{h}} \rrbracket_{Z}=0, \quad \lim _{h \rightarrow \bar{h}} \square \lambda_{h}-\lambda_{\bar{h}} \rrbracket_{Y}=0 .
$$

Hence, using (IV 6), (IV 7) and the argument similar to that in the proof of Lemma 5.3 in Malanowski,1993, we find that condition (5.6) is preserved under small perturbations of $h$, i.e., there exists a neighborhood $G_{\bar{h}}$ of $\bar{h}$ such that for each $h \in G_{\bar{h}}$

$$
\begin{align*}
& \left(D_{z z}^{2} \mathcal{L}\left(z_{h}, \lambda_{h}, h\right) w, w\right)_{\widehat{Z}} \geq \frac{\alpha}{2}\|w\|_{\widehat{Z}}^{2}  \tag{5.10}\\
& \text { for all } w \in\left\{w \in \widehat{Z} \mid D_{z} \varphi\left(z_{h}, h\right) w \in N^{\perp}\right\}
\end{align*}
$$

On the other hand, it follows from (5.7) that there exists a constant $\kappa>0$ such that
$\left(\lambda_{h}, y\right)_{Y} \geq \kappa \rrbracket y \rrbracket_{Y^{*}} \quad$ for all $y \in K \cap N$.
Moreover, for any $y \in Y$ we have

$$
\begin{align*}
& \|y\|_{\widehat{Y}}^{2}=(y, y)_{\widehat{Y}} \leq \rrbracket y \rrbracket_{Y} \cdot \square y \rrbracket_{Y}, \quad \text { i.e., } \\
& \|y\|_{\widehat{Y}}^{2} / \square y \rrbracket_{Y^{*}} \leq \llbracket y \rrbracket_{Y} . \tag{5.12}
\end{align*}
$$

In view of (5.11) and (5.12) it follows from Theorem 1 in Dontchev et al. (1994), that (5.10) constitutes a sufficient optimality condition and that there exist a constant $\beta>0$ and a constant $\varsigma>0$ such that

$$
\begin{equation*}
F(z, h) \geq F\left(z_{h}, h\right)+\beta\left\|z-z_{h}\right\|_{\widehat{Z}}^{2} \tag{5.13}
\end{equation*}
$$

for all $z \in \mathcal{Z}_{h}^{\varsigma}$ feasible for $\left(P_{h}\right)$ and for all $h \in G_{\bar{h}}$, where

$$
\mathcal{Z}_{h}^{\varsigma}=\left\{z \in Z \mid \square z-z_{h} \rrbracket z<\varsigma\right\} .
$$

From Theorem 5.2 and from (5.13) we obtain:
Corollary 5.3 If (IV 1)-(IV 8) are satisfied then there exist a neighborhood $G_{\bar{h}} \subset G$ of $\bar{h}$ and a neighborhood $\mathcal{Z}_{\bar{h}} \subset Z$ of $z_{\bar{h}}$ such that for each $h \in G_{\bar{h}}$ there is a unique in $\mathcal{Z}_{\bar{h}}$ solution $z_{h}$ of $\left(P_{h}\right)$ and a unique associated Lagrange multiplier $\lambda_{h}$. Moreover, there exists a constant $c>0$ such that

$$
\left\|z_{h_{1}}-z_{h_{2}}\right\|_{\widehat{Z}},\left\|\lambda_{h_{1}}-\lambda_{h_{2}}\right\|_{\widehat{Y}} \leq c \llbracket h_{1}-h_{2} \rrbracket_{H}
$$

for all $h_{1}, h_{2} \in G_{\bar{h}}$.

## 6. Optimal control problems

In this section the obtained abstract stability results will be applied to state and control constrained optimal control problems for nonlinear ordinary differential equations.

As before $G \subset H$ denotes an open set of admissible parameters and for each $h \in G$ we consider the following optimal control problem:

$$
\begin{aligned}
& \left(O_{h}\right) \\
& \begin{array}{l}
\text { find }\left(u_{h}, x_{h}\right) \in L^{\infty}\left(0, T ; \mathbf{R}^{m}\right) \times W^{1, \infty}\left(0, T ; \mathbf{R}^{n}\right) \text { such that } \\
F\left(u_{h}, x_{h}, h\right)=\min _{u, x}\left\{F(u, x, h):=\int_{0}^{T} f^{0}(u(t), x(t), h) d t\right\} \\
\begin{array}{ll}
\text { subject to } \\
\dot{x}(t)=f(u(t), x(t), h) & \text { for a.a. } t \in[0, T], \\
x(0)=x_{0}(h), & \\
\begin{array}{ll}
\theta(u(t), h) \leq 0 & \text { for a.a. } t \in[0, T], \\
v(x(t), h) \leq 0 & \text { for all } t \in[0, T],
\end{array}
\end{array} l
\end{array}
\end{aligned}
$$

where

$$
\theta(\cdot, \cdot): \mathbf{R}^{m} \times G \mapsto \mathbf{R}^{k}, \quad v(\cdot, \cdot): \mathbf{R}^{n} \times G \mapsto \mathbf{R}^{l}
$$

In order to represent $\left(O_{h}\right)$ in the form $\left(P_{h}\right)$ we put:

$$
\begin{aligned}
& Z=L^{\infty}\left(0, T ; \mathbf{R}^{m}\right) \times W^{1, \infty}\left(0, T ; \mathbf{R}^{n}\right) \\
& Y=L^{\infty}\left(0, T ; \mathbf{R}^{n}\right) \times \mathbf{R}^{n} \times L^{\infty}\left(0, T ; \mathbf{R}^{k}\right) \times W^{1, \infty}\left(0, T ; \mathbf{R}^{l}\right) \\
& K=K_{1} \times K_{2} \times K_{3} \times K_{4}, \\
& K_{1}=\{0\}, \quad K_{2}=\{0\}, \\
& K_{3}=\left\{u \in L^{\infty}\left(0, T ; \mathbf{R}^{k}\right) \mid u^{i}(t) \geq 0, i=1, \ldots, k \text { for a.a. } t \in[0, T]\right\}, \\
& K_{4}=\left\{x \in W^{1, \infty}\left(0, T ; \mathbf{R}^{l}\right) \mid x^{j}(t) \geq 0, j=1, \ldots, l \text { for all } t \in[0, T]\right\} \\
& F(z, h)=F(u, x, h), \\
& \varphi(z, h)=\left(\dot{x}-f(u, x, h), x(0)-x_{0}(h),-\theta(u, h),-v(x, h)\right)
\end{aligned}
$$

As in Section 3 we introduce the Hilbert spaces

$$
\begin{aligned}
& \widehat{Z}=L^{2}\left(0, T ; \mathbf{R}^{m}\right) \times W^{1,2}\left(0, T ; \mathbf{R}^{n}\right), \quad \text { and } \\
& \widehat{Y}=L^{2}\left(0, T ; \mathbf{R}^{n}\right) \times \mathbf{R}^{n} \times L^{2}\left(0, T ; \mathbf{R}^{k}\right) \times W^{1,2}\left(0, T ; \mathbf{R}^{l}\right)
\end{aligned}
$$

As in (2.1) we treat $\widehat{Y}$ as the pivot space putting $\widehat{Y}=\widehat{Y}^{*}$, but for technical reasons it will be more convenient not to identify $\widehat{Z}$ with $\widehat{Z}^{*}$ but to introduce another pivot space:

$$
\widehat{V}=L^{2}\left(0, T: \mathbf{R}^{m}\right) \times M^{2}\left(0, T: \mathbf{R}^{n}\right)
$$

We have

$$
Z \subset \widehat{Z} \subset \widehat{V}=\widehat{V}^{*} \subset \widehat{Z}^{*} \subset Z^{*}
$$

and the considerations of Sections 3 and 5 remain virtually intact, if in (IV 3) the spaces $\widehat{Z}, L(\widehat{Z}, \widehat{Z})$ and $L(H, \widehat{Z})$ are substituted by $\widehat{V}, L(\widehat{Z}, \widehat{V})$ and $L(H, \widehat{V})$, respectively.

In order to apply Theorem 5.2 we have to verify all assumptions of this theorem.

To simplify notation let us put:

$$
\begin{aligned}
& A(t):=D_{x} f\left(u_{\bar{h}}(t), x_{\bar{h}}(t), \bar{h}\right), \\
& B(t):=D_{u} f\left(u_{\bar{h}}(t), x_{\bar{h}}(t), \bar{h}\right), \\
& \Theta(t):=D_{u} \theta\left(u_{\bar{h}}(t), \bar{h}\right), \\
& \Upsilon(t):=D_{x} v\left(x_{\bar{h}}, \bar{h}\right) .
\end{aligned}
$$

We assume:
(V 1) There exists a (local) solution ( $u_{\bar{h}}, x_{\bar{h}}$ ) of $\left(O_{\bar{h}}\right)$, which satisfies the following regularity condition:

$$
\left(u_{\bar{h}}, x_{\bar{h}}\right) \in C^{0}\left(0, T ; \mathbf{R}^{m}\right) \times C^{1}\left(0, T ; \mathbf{R}^{n}\right)
$$

(V 2) $f^{0}(\cdot, \cdot, \cdot), f(\cdot, \cdot),, \theta(\cdot, \cdot), v(\cdot, \cdot)$ and $D_{x} v(\cdot, \cdot)$ are two times Fréchet differentiable in all arguments, and the respective derivatives are locally Lipschitz continuous in $u, x$, $x_{0}(\cdot)$ is Fréchet differentiable.
(V 3) $h$ does not depend on $t$.

Certainly (V 1) implies (IV 1). On the other hand, it is easy to check that (IV 2) and (IV 3) are satisfied by (V 2).

Let us pass to constraint qualifications. For $\epsilon>0$ we introduce the sets:

$$
\begin{array}{ll}
\Psi_{\epsilon}^{i}:=\left\{t \in[0, T] \mid-\theta^{i}\left(u_{\bar{h}}(t), \bar{h}\right)>\epsilon\right\}, & i=1, \ldots, k, \\
\Omega_{\epsilon}^{j}:=\left\{t \in[0, T] \mid-v^{j}\left(x_{\bar{h}}(t), \bar{h}\right)>\epsilon\right\}, & j=1, \ldots, l,
\end{array}
$$

and define the following continuous functions (see Figure 3):

$$
\psi_{\epsilon}^{i}(t)= \begin{cases}\theta^{i}\left(u_{\bar{h}}(t), \bar{h}\right)+\epsilon & \text { if } t \in \Psi_{\epsilon}^{i} \\ 0 & \text { if } t \notin \Psi_{\epsilon}^{i},\end{cases}
$$

and

$$
\omega_{\epsilon}^{j}(t)= \begin{cases}v^{j}\left(x_{\bar{h}}(t), \bar{h}\right)+\epsilon & \text { if } t \in \Omega_{\epsilon}^{j} \\ 0 & \text { if } t \notin \Omega_{\epsilon}^{j} .\end{cases}
$$



Figure 3

We introduce the following $(k \times k)$ and $(l \times l)$ diagonal matrices

$$
U_{\epsilon}(t)=\operatorname{diag} \psi_{\epsilon}^{i}(t) \quad T_{\epsilon}(t)=\operatorname{diag} \omega_{\epsilon}^{j}(t)
$$

and define $(k+l) \times(m+k+l)$ matrices

$$
V_{\epsilon}(t)=\left[\begin{array}{ccc}
\Theta(t) & U_{\epsilon}(t) & 0 \\
\Upsilon(t) B(t) & 0 & T_{\epsilon}(t)
\end{array}\right]
$$

In addition to (V 1)-(V 3) we assume:
(V 4) $v^{j}\left(x_{0}(\bar{h}), \bar{h}\right)<.0$, for $j=1,2, \ldots, l$,
(V 5) there exists $\eta>0$ such that

$$
\left|V_{0}(t) V_{0}^{*}(t) \zeta\right| \geq \eta|\zeta| \quad \text { for all } \zeta \in \mathbf{R}^{k+l} \text { and all } t \in[0, T]
$$

Roughly speaking condition (V 5) has the meaning that all gradients of the active control constraints $\theta^{i}\left(u_{\bar{h}}(t), \bar{h}\right)=0$ and all gradients of the active state constraints $v^{j}\left(x_{\bar{h}}(t), \bar{h}\right)=0$, transformed into the space $\mathbf{R}^{m}$ by means of the mapping $B^{*}(t): \mathbf{R}^{n} \rightarrow \mathbf{R}^{m}$, are jointly linearly independent, uniformly on $[0, T]$.

By (V 2) and (V 5) (cf. Lemma 7.2 in Malanowski,1993), there exists $\epsilon>0$ such that

$$
\begin{equation*}
\left|V_{\epsilon}(t) V_{\epsilon}^{*}(t) \zeta\right| \geq \frac{\eta}{2}|\zeta| \quad \text { for all } \zeta \in \mathbf{R}^{k+l} \text { and all } t \in[0, T] \tag{6.1}
\end{equation*}
$$

The subspace $M \subset Y$ needed in (IV 4) is define as follows:

$$
M=M_{1} \times M_{2} \times M_{3} \times M_{4},
$$

where

$$
\begin{aligned}
& M_{1}=\{0\}, \quad M_{2}=\{0\}, \\
& M_{3}=\left\{u \in L^{\infty}\left(0, T ; \mathbf{R}^{k}\right) \mid u^{i}(t)=0 \text { for a.a. } t \in[0, T] \backslash \Psi_{\frac{e}{2}}^{i}, \quad i=1, \ldots, k\right\}, \\
& M_{4}=\left\{x \in W^{1, \infty}\left(0, T ; \mathbf{R}^{l}\right) \mid x^{j}(t)=0 \text { for all } t \in[0, T] \backslash \Omega_{\frac{c}{2}}^{j}, j=1, \ldots, l\right\},
\end{aligned}
$$

and $\epsilon$ is given in (6.1).
It is easy to see that for any $\epsilon>0, M \subset M_{0}$, where $M_{0}$ is given by (4.3). Moreover, if we choose as $\mathcal{Y}_{0} \subset Y$ the open ball of radius $\epsilon / 4$ about $-\varphi\left(z_{\bar{h}}, \bar{h}\right)$, then (5.5) holds.
The mapping II : $Y \mapsto M$ is defined by

$$
\begin{array}{ll}
\Pi_{1}=0, & \Pi_{2}=0, \\
\left(\Pi_{3} y\right)(t)=U_{\frac{c}{2}}(t) u(t), & \left(\Pi_{4} y\right)(t)=T_{\frac{s}{2}}(t) x(t) .
\end{array}
$$

It follows from Lemma 4.3 in Malanowski (1994), that for the above difinition of $M$, condition (V 4) and (V5) imply (5.4), i.e., (IV 4) is fully satisfied.

Let us define the following Lagrangian associated with $\left(O_{h}\right)$ :

$$
\begin{aligned}
& \mathcal{L}: \quad L^{\infty}\left(0, T ; \mathbf{R}^{m}\right) \times W^{1, \infty}\left(0, T ; \mathbf{R}^{n}\right) \times\left(L^{\infty}\left(0, T ; \mathbf{R}^{n}\right)\right)^{*} \times \\
& \times \mathbf{R}^{n} \times\left(L^{\infty}\left(0, T ; \mathbf{R}^{k}\right)\right)^{*} \times\left(W^{1, \infty}\left(0, T ; \mathbf{R}^{l}\right)\right)^{*} \times G \mapsto \mathbf{R}^{1}, \\
& \mathcal{L}(u, x, q, \rho, \kappa, \nu, h)=F(u, x, h)+(q, \dot{x}-f(u, x, h))+ \\
&+\left\langle\rho, x(0)-x_{0}(h)\right\rangle+(\kappa, \theta(u, h))+ \\
&+\langle\nu(0), v(x(0), h)\rangle+\left(D_{x} v^{*}(x, h) \dot{\nu}, f(u, x, h)\right) .
\end{aligned}
$$

By (IV 4) there exist unique Lagrange multipliers ( $q_{\bar{h}}, \rho_{\bar{h}}, \kappa_{\bar{h}}, \nu_{\bar{h}}$ ) associated with ( $u_{\bar{h}}, x_{\bar{h}}$ ). Using (V 1) and (V 5) it can be shown (cf. Corollary 4.6 in Hager,1979), that the multipliers are more regular. Namely:

$$
q_{\bar{h}} \in C^{1}\left(0, T ; \mathbf{R}^{n}\right), \quad \rho_{\bar{h}} \in \mathbf{R}^{n}, \quad \kappa_{\bar{h}} \in C^{0}\left(0, T ; \mathbf{R}^{k}\right), \quad \nu_{\bar{h}} \in C^{1}\left(0, T ; \mathbf{R}^{l}\right),
$$

i.e., in particular (IV 5) holds.

Let us define the following augmented Hamiltonian:

$$
\begin{aligned}
& \mathcal{H}(t)=f^{0}\left(u_{\bar{h}}(t), x_{\bar{h}}(t), \bar{h}\right)-\left\langle q_{\bar{h}}(t), f\left(u_{\bar{h}}(t), x_{\bar{h}}(t), \bar{h}\right)\right\rangle+ \\
& +\left\langle\kappa_{\bar{h}}(t), \theta\left(u_{\bar{h}}(t), \bar{h}\right)\right\rangle+\left\langle D_{x} v^{*}\left(x_{\bar{h}}(t), \bar{h}\right) \dot{\mu}_{\bar{h}}(t), f\left(u_{\bar{h}}(t), x_{\bar{h}}(t), \bar{h}\right)\right\rangle .
\end{aligned}
$$

In order to introduce a second order sufficient optimality condition, for $\eta>0$ we define the sets

$$
\Xi_{\eta}^{i}=\left\{t \in[0, T] \mid \kappa_{\bar{h}}^{i}(t)>\eta\right\} .
$$

Let us introduce the following subspaces:

$$
\begin{aligned}
& U_{\eta}(t)=\left\{u \in \mathbf{R}^{m} \mid\left\langle\theta^{i}(t), u\right\rangle=0 \text { for } t \in \Xi_{\eta}^{i}, \quad i=1,2, \ldots, k\right\}, \\
& \mathcal{U}_{\eta}=\left\{u \in L^{2}\left(0, T ; \mathbf{R}^{m}\right) \mid u(t) \in U_{\eta}(t) \quad \text { for a.a. } t \in[0, T]\right\} .
\end{aligned}
$$

We assume:
(V 6) There exists $\gamma>0$ such that

$$
\begin{equation*}
\left\langle u, D_{u u}^{2} \mathcal{H}(t) u\right\rangle \geq \gamma|u|^{2} \tag{6.2}
\end{equation*}
$$

for all $u \in U_{0}(t)$ and all $t \in[0, T]$. Moreover

$$
\begin{align*}
& \int_{0}^{T}\left(\left[u^{T}(t), x^{T}(t)\right]\left[\begin{array}{ll}
D_{u u}^{2} \mathcal{H}(t) & D_{u x}^{2} \mathcal{H}(t) \\
D_{x u}^{2} \mathcal{H}(t) & D_{x x}^{2} \mathcal{H}(t)
\end{array}\right]\left[\begin{array}{l}
u(t) \\
x(t)
\end{array}\right]\right) d t \geq \\
& \quad \geq \gamma\left(\|u\|_{2}^{2}+\|x\|_{1,2}^{2}\right) \tag{6.3}
\end{align*}
$$

for all pairs $(u, x) \in \mathcal{U}_{0} \times W^{1,2}\left(0, T ; \mathbf{R}^{n}\right)$ satisfying

$$
\begin{aligned}
\dot{x}(t) & =A(t) x(t)+B(t) u(t) \\
x(0) & =0 .
\end{aligned}
$$

Note that it follows from (V6) that there exists $\eta>0$ such that (6.2) and (6.3) hold with $\gamma$ substituted by $\gamma / 2$ for the subspaces $U_{\eta}(t)$ and $\mathcal{U}_{\eta}$, respectively (see Dontchev et al.,1994). Using this fact it can be shown that condition (IV 6) is satisfied with the subspaces $N=N_{1} \times N_{2} \times N_{3} \times N_{4}$ defined as follows:

$$
\begin{array}{ll}
N_{1}=L^{\infty}\left(0, T ; \mathbf{R}^{n}\right), & N_{2}=\mathbf{R}^{n}, \\
N_{3}=\left\{u \in L^{\infty}\left(0, T ; \mathbf{R}^{m}\right) \left\lvert\, u(t) \in\left[U_{\frac{n}{2}}(t)\right]^{\perp}\right.\right\}, & N_{4}=\{0\} .
\end{array}
$$

Note that the cone $K_{4}^{+}$polar to $K_{4}$ is given by the closure in $\left(W^{1, \infty}\left(0, T ; \mathbf{R}^{l}\right)\right)^{*-}$ topology of the cone

$$
\begin{aligned}
\widehat{K}_{4}^{+}= & \left\{\nu \in W^{1,2}\left(0, T ; \mathbf{R}^{l}\right) \mid \nu^{j}(t) \geq 0, \dot{\nu}^{j}(t)-\right.\text { is non-increasing, } \\
& \text { and } \left.0 \leq \dot{\nu}^{j}(t) \leq \nu^{j}(0), \quad j=1,2, \ldots, l\right\} .
\end{aligned}
$$

Hence we can not expect that (5.7) is satisfied for a non-trivial subspace $N$. For that reason we have chosen $N_{4}=\{0\}$.

To apply Theorem 5.3 it remains to verify assumptions (IV 7) and (IV 8) concerning regularity of the solutions. To this end, for a given value of the parameter

$$
\delta=\left(a_{1}, a_{2}, b_{1}, b_{2}, b_{3}, b_{4}\right) \in \widehat{Z} \times \widehat{Y}=\widehat{X},
$$

let us introduce the linear-quadratic problem $\left(\widehat{L O}_{\delta}\right)$ corresponding to $\left(O_{\bar{h}}\right)$ and analogous to ( $\widehat{L P}{ }_{\delta}$ ):

$$
\left(\widehat{L O}_{\delta}\right) \quad \text { find }\left(w_{\delta}, z_{\delta}\right) \in L^{2}\left(0, T ; \mathbf{R}^{m}\right) \times W^{1,2}\left(0, T ; \mathbf{R}^{n}\right) \text { such that }
$$

$$
\begin{aligned}
& I\left(w_{\delta}, z_{\delta}, \delta\right)=\min _{(w, z)}\{I(w, z, \delta):= \\
& :=\int_{0}^{T}\left(\frac{1}{2}\left[w^{T}(t), z^{T}(t)\right]\left[\begin{array}{ll}
D_{u u}^{2} \mathcal{H}(t) & D_{u u}^{2} \mathcal{H}(t) \\
D_{x u}^{2} \mathcal{H}(t) & D_{x x}^{2} \mathcal{H}(t)
\end{array}\right]\left[\begin{array}{c}
w(t) \\
z(t)
\end{array}\right]+\right. \\
& \left.\left.\quad+\left(q_{1}^{T}(t)+a_{1}^{T}(t)\right) w(t)+\left(q_{2}^{T}(t)+a_{2}^{T}(t)\right) z(t)\right) d t\right\},
\end{aligned}
$$

subject to

$$
\begin{aligned}
& \dot{z}(t)-A(t) z(t)-B(t) w(t)+e_{1}(t)+b_{1}(t)=0 \\
& z(0)+e_{2}+b_{2}=0 \\
& \Theta(t) w(t)+e_{3}(t)+b_{3}(t) \leq 0 \\
& \Upsilon(t) z(t)+e_{4}(t)+b_{4}(t) \leq 0
\end{aligned}
$$

where

$$
\begin{aligned}
& q_{1}=-D_{u u}^{2} \mathcal{H} u_{\bar{h}}-D_{u x}^{2} \mathcal{H} x_{\bar{h}}+D_{u} f^{0}\left(u_{\bar{h}}, x_{\bar{h}}, \bar{h}\right), \\
& q_{2}=-D_{x u}^{2} \mathcal{H} u_{\bar{h}}-D_{x x}^{2} \mathcal{H} x_{\bar{h}}+D_{x} f^{0}\left(u_{\bar{h}}, x_{\bar{h}}, \bar{h}\right), \\
& e_{1}=-\dot{x}_{\bar{h}}+A x_{\bar{h}}+B u_{\bar{h}}, \\
& e_{2}=-x_{\bar{h}}(0), \\
& e_{3}=\theta\left(u_{\bar{h}}, \bar{h}\right)-\Theta u_{\bar{h}}, \\
& e_{4}=v\left(x_{\bar{h}}, \bar{h}\right)-\Upsilon x_{\bar{h}} .
\end{aligned}
$$

Note that conditions (V 1)-(V 6) imply the following regularity of the solution and Lagrange multipliers of $\left(O_{\bar{h}}\right)$ (cf. Lemma 5.2 in Malanowski,1994):

$$
\begin{equation*}
\left(u_{\bar{h}}, \dot{x}_{\bar{h}}, \dot{q}_{\bar{h}}, \kappa_{\bar{h}}, \dot{\nu}_{\bar{h}}\right) \quad \text { are Lipschitz continuous on }[0, T] \tag{6.4}
\end{equation*}
$$

By (V 2) and (6.4) the functions $q_{1}, q_{2}, e_{1}, e_{2}, e_{3}, \dot{e}_{4}$ are uniformly bounded and Lipschitz continuous on $[0, T]$. Let $l_{1}>0$ be the Lipschitz modulus joint for all these functions, i.e.,

$$
l_{1}=\max \left\{\left\|\dot{q}_{1}\right\|_{\infty},\left\|\dot{q}_{2}\right\|_{\infty},\left\|\dot{e}_{1}\right\|_{\infty},\left\|\dot{e}_{3}\right\|_{\infty},\left\|\ddot{e}_{4}\right\|_{\infty}\right\}
$$

Moreover let

$$
m_{1}=\max \left\{\left\|q_{1}\right\|_{\infty},\left\|q_{2}\right\|_{\infty},\left\|e_{1}\right\|_{\infty},\left|e_{2}\right|,\left\|e_{3}\right\|_{\infty},\left\|\dot{e}_{4}\right\|_{\infty}\right\}
$$

As the set of variations $\Delta$ we choose the set of uniformly bounded and Lipschitz continuous functions $\delta \in \widehat{X}$ with the bound $m_{1}$ and the Lipschitz modulus $l_{1}$. This set is convex and compact in $X$.
Let $\left(r_{\delta}, \eta_{\delta}, \pi_{\delta}, \chi_{\delta}\right)$ denote the Lagrange multipliers associated with $\left(w_{\delta}, z_{\delta}\right)$. It is shown in Malanowski (1994) (cf. Propositions 6.6 and 6.7) that, after possible shrinking of $\Delta, \mathbf{x}_{\delta}:=\left(w_{\delta}, \dot{z}_{\delta}, \dot{r}_{\delta}, \pi_{\delta}, \dot{\chi}_{\delta}\right)$ are bounded and Lipschitz continuous uniformly with respect to $\delta \in \Delta$. Let $m_{2}$ and $l_{2}$ be the bound and the Lipschitz modulus, respectively for all these functions.

The set $\widetilde{\Gamma}$ needed in (IV 7) is chosen as the set of functions $\mathbf{x} \in \widehat{X}$ with the above bound and Lipschitz mudulus. Certainly $\widetilde{\Gamma}$ is convex and by the Arzeli-Ascola theorem it is compact in $X$, i.e., condition (IV 7) is satisfied. It is shown in Section 7 of Malanowski (1994), that also condition (IV 8) holds. Hence all assumptions of Theorem 5.2 and Corollary 5.3 are satisfied and by that corollary we obtain:

Theorem 6.1 If assumptions (V 1)-(V 6) are satisfied then there exist a neighborhood $G_{\bar{h}} \subset G$ of $\bar{h}$ and a neighborhood $\mathcal{Z}_{\bar{h}} \subset Z$ of $\left(u_{\bar{h}}, x_{\bar{h}}\right)$, such that for each $h \in G_{\bar{h}}$ there is a unique in $\mathcal{Z}_{\bar{h}}$ solution $\left(u_{h}, z_{h}\right)$ of $\left(O_{h}\right)$ and a unique associated Lagrange multipliers ( $q_{h}, \rho_{h}, \kappa_{h}, \nu_{h}$ ).

Moreover, there exists a constant $c>0$ such that
$\left\|u_{h_{1}}-u_{h_{2}}\right\|_{2},\left\|x_{h_{1}}-x_{h_{2}}\right\|_{1,2} \leq c \llbracket h_{1}-h_{2} \rrbracket_{H}$,
$\left\|q_{h_{1}}-q_{h_{2}}\right\|_{1,2},\left|\rho_{h_{1}}-\rho_{h_{2}}\right|,\left\|\kappa_{h_{1}}-\kappa_{h_{2}}\right\|_{2},\left\|\nu_{h_{1}}-\nu_{h_{2}}\right\|_{1,2} \leq c \rrbracket h_{1}-h_{2} \|_{H}$,
for all $h_{1}, h_{2} \in G_{\bar{h}}$.

## References

Alt W. (1990A) Parametric programming with applications to optimal control and sequential quadratic programming, Bayreuther Mathematische Schriften 35, 1-37.
Alt W. (1990B) Stability of Solutions and the Lagrange-Newton Method for Nonlinear Optimization and Optimal Control Problems, (Habilitationsschrift), Universität Bayreuth, Bayreuth.
Auslender A. and Cominetti R. (1990) First and second order sensitivity analysis of nonlinear programs under directional constraint qualification condition, Optimization 21, 351-363.
Brezis H. (1983) Analyse Functionelle: Théorie et Applications, Masson, Paris.
Bonnans J. F. (1992) Directional derivatives of optimal solutions in smooth nonlinear programming, J. Optim. Theory Appl. 73, 27-46.
Dontchev A. L. and Hager W. W. (1993) Lipschitz stability in nonlinear control and optimization, SIAM J. Control and Optimization 31, 569-603.
Dontchev A.L., Hager W.W., Poore A.B., Yang B. (1994) Optimality, stability and convergence in nonlinear control (to be published).
Dunn J. C. and Tian T. (1992) Variants of the Kuhn-Tucker sufficient conditions in cones of nonnegative functions, SIAM J. Control and Optimization 30, 1361-1384.
Gauvin J. and Janin R. (1988) Directional behaviour of optimal solutions in nonlinear mathematical programming, Math. Oper. Res. 13, 629-649.
Gollan B. (1981) Perturbation theory for abstract optimization problems, J. Optim. Theory Appl. 35, 417-441.
HAGER W. W. (1979) Lipschitz continuity for constrained processes, SIAM J. Control and Optimization 17, 321-337.

Ioffe A. (1979) Necessary and sufficient condition for a local minimum 3: Second order conditions and augmented. duality, SIAM J. Control and Optimization 17, 266-288.
Ioffe A. (1991) On sensitivity analysis of nonlinear programs in Banach spaces: the approach via composite unconstrained optimization, to appear in SIAM J. Optimization.

Ito K. and Kunisch K. (1992) Sensitivity analysis of solutions to optimization problems in Hilbert spaces with applications to optimal control and estimation, J. Diff. Equations 99, 1-40.
Jittorntrum K. (1984) Solution point differentiability without strict complementarity in nonlinear programming, Mathematical Programming 21, 127-138.
Kojima M. (1980) Strongly stable stationary solutions in nonlinear programs, in: S.M. Robinson (ed.): Analysis and computation of fixed points, Academic Press, New York.
Malanowski K. (1992) Second order conditions and constraint qualifications in stability and sensitivity analysis of solutions to optimization problems in Hilbert spaces, Appl. Math. Opt. 25, 51-79.
Malanowski K. (1993) Two-norm approach in stability and sensitivity analysis of optimization and optimal control problems, Advances in Math. Sc. Appl. 2, 397-443.
Malanowski K. (1994) Stability and sensitivity of solutions to nonlinear optimal control problems, to appear in Appl. Math. Opt..
Maurer H. (1981) First and second order sufficient optimality conditions in mathematical programming and optimal control, Mathematical Programming Study 14, 163-177.
Robinson S. M. (1980) Strongly regular generalized equations, Mathematics of Operations Research 5, 43-62.
Robinson S. M. (1987) Local structure of feasible sets in nonlinear programming, Part III: Stability and sensitivity, Mathematical Programming Study 30, 45-66.
Shapiro A. (1992) Perturbation analysis of optimization problems in Banach spaces, Numerical Funct. Anal. Optim. 13, 97-116.
Shapiro A. (1994) Sensitivity analysis of parametrized programs via generalized equations, (to be published).
Shapiro A. and Bonnans J. F. (1992) Sensitivity analysis of parametrized programs under cone constraints, SIAM J. Control and Optimization 30, 1409-1421.
Tröltzsch F. (1991) Approximation of nonlinear parabolic boundary control problems by the Fourier method - convergence of optimal controls, Optimization 22, 83-98.


[^0]:    *) Supported by grant No. 3 P403 00205 from Komitet Badań Naukowych [State Committee for Scientific Research].

