# Necessary Conditions for the Optimal Control <br> of Differential-Algebraic Equations with Arbitrary Index <br> Peter Kunkel and Volker Mehrmann 

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## Necessary Conditions for the Optimal Control of Differential-Algebraic Equations with Arbitrary Index

1 Preliminaries
2 DAE theory
3 Necessary conditions
4 Numerical treatment
5 Conclusions

## 1 Preliminaries

### 1.1 Problem formulation

We consider optimal control problems of the form

$$
\mathcal{J}(x, u)=\mathcal{M}(x(\bar{t}))+\int_{\underline{t}}^{\bar{t}} \mathcal{K}(t, x(t), u(t)) d t=\min !
$$

subject to the differential-algebraic equation (DAE)

$$
F(t, x, u, \dot{x})=0, \quad x(\underline{t})=\underline{x}
$$

where

$$
\mathcal{M} \in C\left(\mathbb{D}_{x}, \mathbb{R}\right), \quad \mathcal{K} \in C\left(\mathbb{I} \times \mathbb{D}_{x} \times \mathbb{D}_{u}, \mathbb{R}\right), \quad F \in C\left(\mathbb{I} \times \mathbb{D}_{x} \times \mathbb{D}_{\dot{x}} \times \mathbb{D}_{u}, \mathbb{R}^{n}\right)
$$

with

$$
\mathbb{I}=[\underline{t}, \bar{t}] \quad \text { and } \quad \mathbb{D}_{x}, \mathbb{D}_{\dot{x}} \subseteq \mathbb{R}^{n}, \mathbb{D}_{u} \subseteq \mathbb{R}^{l} \text { open. }
$$

We assume that all functions are sufficiently smooth.

The corresponding linearized problem reads (omitting the argument $t$ )

$$
\mathcal{J}(x, u)=\frac{1}{2} x(\bar{t})^{T} M x(\bar{t})+\frac{1}{2} \int_{\underline{t}}^{\bar{t}}\left(x^{T} W x+2 x^{T} S u+u^{T} R u\right) d t=\min !
$$

subject to

$$
E \dot{x}=A x+B u+f, \quad x(\underline{t})=\underline{x},
$$

where

$$
M \in \mathbb{R}^{n, n}, \quad W \in C\left(\mathbb{I}, \mathbb{R}^{n, n}\right), \quad S \in C\left(\mathbb{I}, \mathbb{R}^{n, l}\right), \quad R \in C\left(\mathbb{I}, \mathbb{R}^{l, l}\right),
$$

and

$$
E, A \in C\left(\mathbb{I}, \mathbb{R}^{n, n}\right), \quad B \in C\left(\mathbb{I}, \mathbb{R}^{n, l}\right), \quad f \in C\left(\mathbb{I}, \mathbb{R}^{n}\right)
$$

are sufficiently smooth.
We are allowed to assume that $M$ is symmetric and $W$ und $R$ are pointwise symmetric.

### 1.3 Abstract optimization

For the abstract problem

$$
\mathcal{J}(z)=\min !
$$

subject to

$$
\mathcal{F}(z)=0
$$

where

$$
\mathcal{J} \in C(\mathbb{D}, \mathbb{R}), \quad \mathcal{F} \in C(\mathbb{D}, \mathbb{Y}), \quad \mathbb{D} \subseteq \mathbb{Z} \text { open }
$$

with real Banach spaces $\mathbb{Y}, \mathbb{Z}$, we have the following (classical) result due to Ljusternik (1934):

Theorem
Let $z^{*} \in \mathbb{Z}$ be a local minimum of the above problem and assume that

- $\mathcal{J}$ is Fréchet differentiable in $z^{*}$,
- $\mathcal{F}$ is a submersion in $z^{*}$, i.e., $\mathcal{F}$ is Fréchet differentiable in a neighborhood of $z^{*}$ with Fréchet derivative $D \mathcal{F}\left(z^{*}\right): \mathbb{Z} \rightarrow \mathbb{Y}$ surjective and kernel $D \mathcal{F}\left(z^{*}\right)$ continuously projectable.

Then there exists a functional $\wedge$ in the dual space $\mathbb{Y}^{*}$ of $\mathbb{Y}$ with $D \mathcal{J}\left(z^{*}\right) \wedge z+\wedge\left(D \mathcal{F}\left(z^{*}\right) \wedge z\right)=0$ for all $\wedge z \in \mathbb{Z}$.

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D \mathcal{J}\left(z^{*}\right) \Delta z+\wedge\left(D \mathcal{F}\left(z^{*}\right) \Delta z\right)=0 \text { for all } \Delta z \in \mathbb{Z}
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## Properties

- The Lagrange multiplier $\wedge$ is unique.
- The above necessary conditions transform covariantly with diffeomorphisms $\phi: \mathbb{Z} \rightarrow \mathbb{Z}$.

Questions

- What is a suitable abstract formulation of a DAE?
- What is a suitable representation of $\wedge$ ?
- How do the necessary conditions in the case of DAEs look like?


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## 2 DAE theory

### 2.1 Illustrative example

The system

$$
\begin{aligned}
& \dot{x}_{1}=x_{4}, \quad \dot{x}_{4}=2 x_{1} x_{7} \\
& \dot{x}_{2}=x_{5}, \quad \dot{x}_{5}=2 x_{2} x_{7} \\
& \dot{x}_{3}=x_{6}, \quad \dot{x}_{6}=-1-x_{7}, \\
& 0=x_{3}-x_{1}^{2}-x_{2}^{2}
\end{aligned}
$$

describes the movement of a mass point on a paraboloid under the influence of gravity.

Differentiating the constraint twice and eliminating the arising derivatives of the unknowns yields

$$
\begin{aligned}
& 0=x_{6}-2 x_{1} x_{4}-2 x_{2} x_{5} \\
& 0=-1-x_{7}-2 x_{4}^{2}-4 x_{1}^{2} x_{7}-2 x_{5}^{2}-4 x_{2}^{2} x_{7}
\end{aligned}
$$

Hence, we may replace the original problem by

```
\mp@subsup{\dot{x}}{1}{}}=\mp@subsup{x}{4}{}
x}=\mp@subsup{x}{5}{\prime
    0}=\mp@subsup{x}{3}{}-\mp@subsup{x}{1}{2}-\mp@subsup{x}{2}{2}
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\end{aligned}
$$

Hence, we may replace the original problem by

$$
\begin{aligned}
\dot{x}_{1} & =x_{4} \\
\dot{x}_{2} & =x_{5} \\
0 & =x_{3}-x_{1}^{2}-x_{2}^{2} \\
\dot{x}_{4} & =2 x_{1} x_{7} \\
\dot{x}_{5} & =2 x_{2} x_{7} \\
0 & =x_{6}-2 x_{1} x_{4}-2 x_{2} x_{5}, \\
0 & =-1-x_{7}-2 x_{4}^{2}-4 x_{1}^{2} x_{7}-2 x_{5}^{2}-4 x_{2}^{2} x_{7}
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### 2.2 Observations

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- In order to solve a DAE we may be forced to differentiate parts of the given DAE.
- If differentiation is necessary, we get additional (so-called hidden) constraints for the states.
- The second formulation allows for weaker smoothness requirements for the solution.

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- The second formulation allows for weaker smoothness requirements for the solution.

For the following, it is convenient not to distinguish between states and controls, so-called behavior approach. Introducing

$$
z=\left[\begin{array}{l}
x \\
u
\end{array}\right]
$$

we write the DAE as

$$
F(t, z, \dot{z})=0
$$

or

$$
\mathcal{E} \dot{z}=\mathcal{A} z+f, \quad \mathcal{E}=\left[\begin{array}{ll}
E & 0
\end{array}\right], \mathcal{A}=\left[\begin{array}{ll}
A & B
\end{array}\right]
$$

in the linear case.
We first study linear DAEs.

Since we may have to differentiate the given DAE in order to determine its solutions, we consider so-called derivative array equations

$$
M_{\ell} \dot{z}_{\ell}=N_{\ell} z_{\ell}+g_{\ell},
$$

where

$$
\begin{aligned}
\left(M_{\ell}\right)_{i, j} & =\binom{i}{j} \mathcal{E}^{(i-j)}-\binom{i}{j+1} \mathcal{A}^{(i-j-1)}, i, j=0, \ldots, \ell, \\
\left(N_{\ell}\right)_{i, j} & = \begin{cases}\mathcal{A}^{(i)} & \text { for } i=0, \ldots, \ell, j=0, \\
0 & \text { otherwise, },\end{cases} \\
\left(z_{\ell}\right)_{j} & =z^{(j)}, j=0, \ldots, \ell, \\
\left(g_{\ell}\right)_{i} & =f^{(i)}, i=0, \ldots, \ell,
\end{aligned}
$$

consisting of

$$
\begin{array}{rlr}
\mathcal{E} \dot{z} & =\mathcal{A} z+f, \\
(\dot{\mathcal{E}}-\mathcal{A}) \dot{z}+ & \mathcal{E} \ddot{z} & =\dot{\mathcal{A}} z+\dot{f}, \\
(\ddot{\mathcal{E}}-2 \dot{\mathcal{A}}) \dot{z}+(2 \dot{\mathcal{E}}-\mathcal{A}) \ddot{z}+\mathcal{E} \dddot{z} & =\ddot{\mathcal{A}} z+\ddot{f},
\end{array}
$$

etc.

### 2.5 Hypothesis

Hypothesis
There are integers $\mu, a, d \in \mathbb{N}_{0}$ with $a+d=n$ such that

1. $\operatorname{rank} M_{\mu}=(\mu+1) n-a$ on $\mathbb{I}$ and thus the existence of
$Z_{2}$ smooth matrix function, max. rank $a$, orth. columns, $Z_{2}^{T} M_{\mu}=0$ on $\mathbb{I}$,
2. $\operatorname{rank} Z_{2}^{T} N_{\mu}=a$ on $\mathbb{I}$ and thus the existence of
$T_{2}$ smooth matrix function, max. rank $d+l$, orth. columns, $Z_{2}^{T} N_{\mu}\left[\begin{array}{llll}I_{n+l} & 0 & \cdots & 0\end{array}\right]^{T} T_{2}=0$ on $\mathbb{I}$,
3. $\operatorname{rank} \mathcal{E} T_{2}=d$ on $\mathbb{I}$ and thus the existence of
$Z_{1}$ smooth matrix function, max. rank $d$, orth. columns, $\operatorname{rank} Z_{1}^{T} \mathcal{E} T_{2}=d$ on $\mathbb{I}$.

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There are integers $\mu, a, d \in \mathbb{N}_{\mathrm{O}}$ with $a+d=n$ such that

```
1. rank }\mp@subsup{M}{\mu}{}=(\mu+1)n-a on \mathbb{I}\mathrm{ and thus the existence of
    Z}\mp@subsup{Z}{2}{}\mathrm{ smooth matrix function, max. rank a, orth. columns,
    Z}\mp@subsup{2}{2}{T}\mp@subsup{M}{\mu}{}=0\mathrm{ on }\mathbb{I}\mathrm{ ,
2. rank Z}\mp@subsup{Z}{2}{T}\mp@subsup{N}{\mu}{}=a\mathrm{ on }\mathbb{I}\mathrm{ and thus the existence of
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    Z2
3. rankEET2}=d\mathrm{ on }\mathbb{I}\mathrm{ and thus the existence of
    Z1 smooth matrix function, max. rank d, orth. columns,
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### 2.6 Remarks

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- In the case $l=0$, the hypothesis is equivalent with the assumption of a well-defined differentiation index.
- Under the hypothesis, we can extract a so-called reduced DAE

$$
\widehat{E} \dot{x}=\hat{A} x+\hat{B} u+\widehat{f}, \quad \widehat{E}=\left[\begin{array}{c}
\widehat{E}_{1} \\
0
\end{array}\right], \hat{A}=\left[\begin{array}{c}
\widehat{A}_{1} \\
\widehat{A}_{2}
\end{array}\right], \widehat{B}=\left[\begin{array}{c}
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$$

where

$$
\begin{aligned}
& \hat{E}_{1}=Z_{1}^{T} E, \quad \hat{A}_{1}=Z_{1}^{T} A, \quad \widehat{B}_{1}=Z_{1}^{T} B, \quad \widehat{f}_{1}=Z_{1}^{T} f, \\
& {\left[\widehat{A}_{2} \widehat{B}_{2}\right]=Z_{2}^{T} N_{\mu}\left[\begin{array}{llll}
I_{n+l} & 0 & \cdots & 0
\end{array}\right]^{T}, \quad \widehat{f}_{2}=Z_{2}^{T} g_{\mu}}
\end{aligned}
$$

out of the derivative array equations satifying the hypothesis with $\mu=0$.

- Original and reduced DAE have the same (smooth) solutions.
- For the reduced DAE there exists a linear feedback $u=K x+w$ such that the closed loop problem

$$
\widehat{E} \dot{x}=(\widehat{A}+\widehat{B} K) x+\widehat{B} w+\widehat{f}
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satisfies the above hypothesis with $l=0$ and $\mu=0$ for every continous $w$.

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- Original and reduced DAE have the same (smooth) solutions.
- For the reduced DAE there exists a linear feedback $u=K x+w$ such that the closed loop problem

$$
\widehat{E} \dot{x}=(\widehat{A}+\widehat{B} K) x+\widehat{B} w+\widehat{f}
$$

satisfies the above hypothesis with $l=0$ and $\mu=0$ for every continous $w$.

Consequences

- In the optimal control problem, one should replace the original DAE by the reduced DAE.
- Since a feedback defines a diffeomorphism with respect to $z$, we may assume that the reduced DAE satisfies the above hypothesis with $\mu=0$ for every continuous u.
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Using

$$
E \dot{x}=E E^{+} E \dot{x}=E \frac{d}{d t}\left(E^{+} E x\right)-E \frac{d}{d t}\left(E^{+} E\right) x
$$

we interpret

$$
E \dot{x}=A x+B u+f, \quad x(\underline{t})=\underline{x}, \quad \underline{x} \in \operatorname{range} E^{+} E(\underline{t})
$$

as

$$
E \frac{d}{d t}\left(E^{+} E x\right)=\left(A+E \frac{d}{d t}\left(E^{+} E\right)\right) x+B u+f, \quad\left(E^{+} E x\right)(\underline{t})=\underline{x} .
$$

Consider now

$$
\mathcal{F}: \mathbb{Z}=\mathbb{X} \times \mathbb{U} \rightarrow \mathbb{Y}
$$

with the Banach spaces

```
\mathbb{X}=\mp@subsup{C}{\mp@subsup{E}{}{+}E}{1}(\mathbb{I},\mp@subsup{\mathbb{R}}{}{n})={x\in\mp@subsup{C}{}{0}(\mathbb{I},\mp@subsup{\mathbb{R}}{}{n})|E+Ex\in\mp@subsup{C}{}{1}(\mathbb{I},\mp@subsup{\mathbb{R}}{}{n})},
U}=\mp@subsup{C}{}{0}(\mathbb{I},\mp@subsup{\mathbb{R}}{}{l})
    \mathbb{Y}=\mp@subsup{C}{}{0}(\mathbb{I},\mp@subsup{\mathbb{R}}{}{n})\times\operatorname{range}\mp@subsup{E}{}{+}E(\underline{t})
```

defined by

$$
\mathcal{F}(z)=\left(E \frac{d}{d t}\left(E^{+} E x\right)-\left(A+E \frac{d}{d t}\left(E^{+} E\right)\right) x-B u-f,\left(E^{+} E x\right)(\underline{t})-\underline{x}\right) .
$$

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$$

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$$
\mathcal{F}: \mathbb{Z}=\mathbb{X} \times \mathbb{U} \rightarrow \mathbb{Y}
$$

with the Banach spaces

$$
\begin{aligned}
& \mathbb{X}=C_{E^{+} E}^{1}\left(\mathbb{I}, \mathbb{R}^{n}\right)=\left\{x \in C^{0}\left(\mathbb{I}, \mathbb{R}^{n}\right) \mid E^{+} E x \in C^{1}\left(\mathbb{I}, \mathbb{R}^{n}\right)\right\}, \\
& \mathbb{U}=C^{0}\left(\mathbb{I}, \mathbb{R}^{l}\right), \\
& \mathbb{Y}=C^{0}\left(\mathbb{I}, \mathbb{R}^{n}\right) \times \text { range } E^{+} E(\underline{t})
\end{aligned}
$$

defined by

$$
\mathcal{F}(z)=\left(E \frac{d}{d t}\left(E^{+} E x\right)-\left(A+E \frac{d}{d t}\left(E^{+} E\right)\right) x-B u-f,\left(E^{+} E x\right)(\underline{t})-\underline{x}\right) .
$$

## Theorem

The operator $\mathcal{F}$ is Fréchet differentiable and the restriction $\mathcal{F}(\cdot, u): \mathbb{X} \rightarrow \mathbb{Y}$ is invertible for every $u \in \mathbb{U}$.

Corollary
The operator $\mathcal{F}$ is a submersion in every $z \in \mathbb{Z}$ with $\mathcal{F}(z)=0$.

## Theorem

The operator $\mathcal{F}$ is Fréchet differentiable and the restriction $\mathcal{F}(\cdot, u): \mathbb{X} \rightarrow \mathbb{Y}$ is invertible for every $u \in \mathbb{U}$.

## Corollary

The operator $\mathcal{F}$ is a submersion in every $z \in \mathbb{Z}$ with $\mathcal{F}(z)=0$.

In the nonlinear case, the derivative array equations have the form

$$
F_{\ell}\left(t, z, \dot{z}, \ldots, z^{(\ell+1)}\right)=0,
$$

where

$$
F_{\ell}\left(t, z, \dot{z}, \ldots, z^{(\ell+1)}\right)=\left[\begin{array}{c}
F(t, z, \dot{z}) \\
\frac{d}{d t} F(t, z, \dot{z}) \\
\left(\frac{d}{d t}\right)^{2} F(t, z, \dot{z}) \\
\vdots \\
\left(\frac{d}{d t}\right)^{\ell} F(t, z, \dot{z})
\end{array}\right],
$$

with

$$
\frac{d}{d t} F(t, z, \dot{z})=F_{t}(t, z, \dot{z})+F_{z}(t, z, \dot{z}) \dot{z}+F_{\dot{z}}(t, z, \dot{z}) \ddot{z}
$$

etc.

### 2.11 Hypothesis

Hypothesis

```
There are integers \(\mu, a, d \in \mathbb{N}\) owith \(a+d=n\) such that
\(\mathbb{L}_{\mu}=F_{\mu}^{-1}(\{0\}) \neq \emptyset\)
and (locally)
1. \(\operatorname{rank} F_{\mu ; \dot{z}, \ldots, z^{(\mu+1)}}=(\mu+1) n-a\) on \(\mathbb{L}_{\mu}\) and thus the existence of
\(Z_{2}\) smooth matrix function, max. rank \(a\), orth columns,
\(Z_{2}^{T} F_{\mu ; \dot{z}, \ldots, z^{(\mu+1)}}=0\) on \(\mathbb{L}_{\mu}\),
2. rank \(Z_{2}^{T} F_{\mu ; z}=a\) on \(\mathbb{L}_{\mu}\) and thus the existence of
\(T_{2}\) smooth matrix function, max. rank \(d+l\), orth. columns, \(Z_{2}^{T} F_{\mu ; z} T_{2}=0\) on \(\mathbb{L}_{\mu}\),
3. rank \(F_{\dot{z}} T_{2}=d\) on \(\mathbb{L}_{\mu}\) and thus the existence of
\(Z_{1}\) smooth matrix function, max rank \(d\), orth. columns, \(\operatorname{rank} Z_{1}^{T} F_{\dot{z}} T_{2}=d\) on \(\mathbb{L}_{\mu}\).
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There are integers $\mu, a, d \in \mathbb{N}_{0}$ with $a+d=n$ such that

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## Remarks

- Given a solution $z$ of the DAE in the form of a path

$$
(t, z, \mathcal{P}(t)) \in \mathbb{I}_{\mu},
$$

the matrix functions fixed by the above hypothesis can be defined globally.

- We can then (under some additional technical assumptions) extract a so-called reduced DAE

$$
\widehat{F}(t, x, u, \dot{x})=0, \quad \widehat{F}(t, x, u, \dot{x})=\left[\begin{array}{c}
\hat{F}_{1}(t, x, u, \dot{x}) \\
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\end{array}\right]
$$

where

$$
\widehat{F}_{1}(t, x, u, \dot{x})=Z_{1}^{T} F(t, x, u, \dot{x})
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- Locally original and reduced DAE have the same (smooth) solutions.
2.13 Remarks (cont)19


## Remarks (cont)

- With some further technical assumptions, the above reduced DAE is equivalent to a DAE of the form

$$
\dot{x}_{1}=\mathcal{L}\left(t, x_{1}, u\right), \quad x_{2}=\mathcal{R}\left(t, x_{1}, u\right),
$$

where $\left(x_{1}, x_{2}\right)=Q x$ with a pointwise orthogonal $Q \in C^{1}\left(\mathbb{I}, \mathbb{R}^{n, n}\right)$.

- In the optimal control problem, one should replace the original DAE by the reduced DAE.
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$$
\mathbb{Z}=C^{1}\left(\mathbb{I}, \mathbb{R}^{d}\right) \times C^{0}\left(\mathbb{I}, \mathbb{R}^{a}\right) \times C^{0}\left(\mathbb{I}, \mathbb{R}^{l}\right), \quad \mathbb{Y}=C^{1}\left(\mathbb{I}, \mathbb{R}^{d}\right) \times C^{0}\left(\mathbb{I}, \mathbb{R}^{a}\right) \times \mathbb{R}^{d}
$$

by

$$
\mathcal{F}(z)=\left(\dot{x}_{1}-\mathcal{L}\left(t, x_{1}, u\right), x_{2}-\mathcal{R}\left(t, x_{1}, u\right), x_{1}(\underline{t})\right)
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where $z=\left(x_{1}, x_{2}, u\right)$.

## Theorem

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## 3 Necessary conditions

## Theorem

The necessary condition for a local minimum of the linear-quadratic optimal control problem is given by the boundary value problem

$$
\begin{aligned}
E \frac{d}{d t}\left(E^{+} E x\right) & =\left(A+E \frac{d}{d t}\left(E^{+} E\right)\right) x+B u+f,\left(E^{+} E x\right)(\underline{t})=\underline{x} \\
E^{T} \frac{d}{d t}\left(E E^{+} \lambda\right) & =W x+S u-\left(A+E E^{+} \dot{E}\right)^{T} \lambda,\left(E E^{+} \lambda\right)(\bar{t})=-\left(E^{+T} M x\right)(\bar{t}) \\
0 & =S^{T} x+R u-B^{T} \lambda
\end{aligned}
$$

## Remark

The Lagrange multiplier $\wedge: \mathbb{Y} \rightarrow \mathbb{R}$ has the form

$$
\wedge(g, r)=\int_{\underline{t}}^{\bar{t}} \lambda^{T} g d t+\gamma^{T} r
$$

with
$\lambda \in C_{E E^{+}}^{1}\left(\mathbb{I}, \mathbb{R}^{n}\right), \quad \gamma \in$ cokernel $E(\underline{t})$
given by the above boundary value problem and by $\gamma=E(t)^{T} \lambda(\underline{t})$.

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### 3.2 Observations

23
## Observations

- The necessary condition is a boundary value problem for a DAE even when we start with an ODE.
- It may happen that this DAE satisfies the above hypothesis only for non-vanishing $\mu$.
- We can characterize the case $\mu=0$.
- We can achieve $\mu=0$ just by modifying the costs.
- The boundary value problem transforms covariantly with respect to feedback controls in the DAE.
- For every additional end condition on the state, we loose an initial condition on the Lagrangian function.


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### 3.3 Adjoint DAE

The DAEs

$$
E \frac{d}{d t}\left(E^{+} E x\right)=\left(A+E \frac{d}{d t}\left(E^{+} E\right)\right) x, \quad E^{T} \frac{d}{d t}\left(E E^{+} \lambda\right)=-\left(A+E E^{+} \dot{E}\right)^{T} \lambda
$$

are the correct formulations of the problems

$$
E \dot{x}=A x, \quad E^{T} \dot{\lambda}=-(A+\dot{E})^{T} \lambda
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The role of $\lambda$ suggests to call the DAE for $\lambda$ the adjoint DAE of the DAE for $x$ and to call $\left(-E^{T},(A+\dot{E})^{T}\right)$ the adjoint pair of the pair $(E, A)$.

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## Theorem

The necessary condition for a local minimum of the nonlinear optimal control problem is given by the boundary value problem

$$
\begin{aligned}
& \dot{x}_{1}=\mathcal{L}\left(t, x_{1}, u\right), x_{1}(\underline{t})=\underline{x}_{1} \\
& x_{2}=\mathcal{R}\left(t, x_{1}, u\right) \\
& \dot{\lambda}_{1}=\mathcal{K}_{x_{1}}\left(t, x_{1}, x_{2}, u\right)^{T}-\mathcal{L}_{x_{1}}\left(t, x_{1}, x_{2}, u\right)^{T} \lambda_{1}-\mathcal{R}_{x_{1}}\left(t, x_{1}, u\right)^{T} \lambda_{2} \\
& \quad \lambda_{1}(\bar{t})=-\mathcal{M}_{x_{1}}\left(x_{1}(\bar{t}), x_{2}(\bar{t})\right)^{T} \\
& 0=\mathcal{K}_{x_{2}}\left(t, x_{1}, x_{2}, u\right)^{T}+\lambda_{2}, \\
& 0=\mathcal{K}_{u}\left(t, x_{1}, x_{2}, u\right)^{T}-\mathcal{L}_{u}\left(t, x_{1}, u\right)^{T} \lambda_{1}-\mathcal{R}_{u}\left(t, x_{1}, u\right)^{T} \lambda_{2} \\
& \gamma=\lambda_{1}(\underline{t})
\end{aligned}
$$

4 Numerical treatment

Returning to the original data, we must deal in the linear case with the DAE

$$
\begin{aligned}
Z_{1}^{T} E \dot{x} & =Z_{1}^{T} A x+Z_{1}^{T} B u+Z_{1}^{T} f, \\
0 & =Z_{2}^{T} \widehat{N}_{\mu}\left[\begin{array}{ll}
I_{n} & 0
\end{array}\right]^{T} x+Z_{2}^{T} \widehat{N}_{\mu}\left[\begin{array}{ll}
0 & I_{l}
\end{array}\right]^{T} u+Z_{2}^{T} g_{\mu}, \\
\frac{d}{d t}\left(E^{T} Z_{1} \lambda_{1}\right) & =W x+S u-A^{T} Z_{1} \lambda_{1}-\left[\begin{array}{ll}
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\end{array}\right] \widehat{N}_{\mu}^{T} Z_{2} \lambda_{2}, \\
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\end{aligned}
$$

where

$$
\widehat{N}_{\mu}=N_{\mu}\left[\begin{array}{llll}
I_{n+l} & 0 & \cdots & 0
\end{array}\right]^{T}
$$

## Problem

We know that there exist smooth functions $Z_{1}$ and $Z_{2}$, but numerically we only can determine $Z_{1} U_{1}$ and $Z_{2} U_{2}$ with pointwise orthogonal but in general non-smooth $U_{1}$ and $U_{2}$.

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Observe that

$$
E^{T} Z_{1} \lambda_{1}=E^{T} Z_{1} Z_{1}^{T} Z_{1} \lambda_{1}=E^{T} Z_{1} Z_{1}^{T} \hat{\lambda}_{1},
$$

with $\hat{\lambda}_{1}=Z_{1} \lambda_{1} \in \operatorname{range} Z_{1}$.
Assuming that $\hat{Z}_{1}$ gives a pointwise orthogonal $\left[\begin{array}{lll}Z_{1} & \hat{Z}_{1}\end{array}\right]$, we can write the property $\widehat{\lambda}_{1} \in \operatorname{range} Z_{1}$ as $\widehat{Z}_{1}^{T} \widehat{\lambda}_{1}=0$.

The projector $Z_{1} Z_{1}^{T}$ is unique and thus smooth.
We can proceed in a similar way for $\hat{N}_{\mu}^{T} Z_{2} \lambda_{2}$ which gives $\hat{\lambda}_{2}=Z_{2} \lambda_{2} \in$ range $Z_{2}$.

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Therefore, we actually treat

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\end{aligned}
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## Remarks

- Numerical solutions do not depend on non-smooth scalings from the left.
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In the nonlinear case, a corresponding approach yields

$$
\begin{aligned}
Z_{1}^{T} F & =0 \\
F_{\mu} & =0 \\
\frac{d}{d t}\left(F_{\dot{x}}^{T} Z_{1} Z_{1}^{T} \hat{\lambda}_{1}\right) & =\mathcal{K}_{x}^{T}+F_{x}^{T} \widehat{\lambda}_{1}+F_{\mu ; x}^{T} \hat{\lambda}_{2} \\
0 & =\mathcal{K}_{u}^{T}+F_{u}^{T} \hat{\lambda}_{1}+F_{\mu ; u}^{T} \hat{\lambda}_{2} \\
0 & =\hat{Z}_{1}^{T} \hat{\lambda}_{1} \\
0 & =\widehat{Z}_{2}^{T} \widehat{\lambda}_{2}
\end{aligned}
$$

with the boundary conditions

$$
\left(\widehat{E}_{1}^{+} \hat{E}_{1} x\right)(\underline{t})=\underline{x}, \quad\left(Z_{1}^{T} \hat{\lambda}_{1}\right)(\bar{t})=-\widehat{E}_{1}^{+}(\bar{t})^{T} \mathcal{M}_{x}(x(\bar{t}))^{T}
$$

A model problem for a motor controlled pendulum to be driven into its equilibrium with minimal costs is given by

$$
\begin{aligned}
J(x, u) & =\int_{0}^{3} u(t)^{2} d t=\min ! & & \\
\text { s.t. } & & & x_{1}(0)=\frac{1}{2} \sqrt{2}, \\
\dot{x}_{1} & =x_{3}, & & x_{2}(0)=-\frac{1}{2} \sqrt{2}, \\
\dot{x}_{2} & =x_{4}, & & \\
\dot{x}_{3} & =-2 x_{1} x_{5}+x_{2} u, & & x_{3}(0)=0, \\
\dot{x}_{4} & =-g-2 x_{2} x_{5}-x_{1} u, & & x_{4}(0)=0, \\
0 & =x_{1}^{2}+x_{2}^{2}-1, & & x_{5}(0)=-\frac{1}{2} g x_{2}(0)
\end{aligned}
$$

It is known that the differential-algebraic equation in the constraint satisfies the above hypothesis with $\mu=2, a=3$, and $d=2$.

Hence, only two scalar initial values are sufficient to describe the initial state, e. g.

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x_{2}(0)=-\frac{1}{2} \sqrt{2}, \quad x_{3}(0)=0
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We used a simple symmetric discretization of order two.
As initial trajectory we took

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\begin{array}{ll}
x_{1}(t)=\frac{1}{2} \sqrt{2}-\frac{1}{6} \sqrt{2} t, & x_{3}(t)=0 \\
x_{2}(t)=-\sqrt{1-x_{1}(t)^{2}}, & x_{4}(t)=0, \quad x_{5}(t)=-\frac{1}{2} g x_{2}(t),
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with all other unknowns set to zero on an equidistant grid of 60 intervals. The required tolerance for the Gauß-Newton method was $10^{-7}$.

Denoting the Euclidian norm of the corresponding Gauß-Newton correction by $\left\|\Delta w_{k}\right\|_{2}$, the course of the iteration was as follows.


| $k$ | $\left\\|\triangle w_{k}\right\\|_{2}$ |
| :---: | :---: |
| 17 | $0.103 \mathrm{D}+01$ |
| 18 | $0.610 \mathrm{D}-02$ |
| 19 | $0.318 \mathrm{D}-06$ |
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The bad convergence behavior in the initial phase is due to the bad initial guess, especially for the Lagrange parameter.

In the final phase, we see quadratic convergence.
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| :---: | :---: |
| 1 | $0.140 \mathrm{D}+03$ |
| 2 | $0.223 \mathrm{D}+03$ |
| $\vdots$ | $\vdots$ |
| 16 | $0.561 \mathrm{D}+01$ |


| $k$ | $\left\\|\Delta w_{k}\right\\|_{2}$ |
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- How can we treat the case when we have $\mu \neq 0$ for the DAE of the boundary value problem?
- How can we utilize the underlying structure of self-adjointness?
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## 5 Conclusions

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- In order to derive necessary conditions for an optimal control, we must apply techniques of index reduction.
- We are allowed to regularize by feedback for simplification.
- The Lagrangian functional possesses an integral representation via a Lagrangian function.
- The necessary conditions have the form of a boundary value problem for a DAE in state, control and Lagrangian function involving the index reduction of the constraint DAE.
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