Friedrichs systems (with complex coefficients)

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Joint work with Krešimir Burazin, Ivana Crnjac, Marko Erceg and Marko Vrdoljak







Classical theory

What are Friedrichs systems? Examples Boundary conditions for Friedrichs systems Existence, uniqueness, well-posedness

Abstract formulation

Graph spaces
Cone formalism of Ern, Guermond and Caplain
Interdependence of different representations of boundary conditions
Kreĭn spaces
Equivalence of boundary conditions

What can we say for the Friedrichs operator now?

Some examples Two-field theory

Concluding remarks

Assumptions:

 $d, r \in \mathbf{N}$, $\Omega \subseteq \mathbf{R}^d$ open and bounded with Lipschitz boundary Γ ;

$$\mathbf{A}_k \in \mathrm{W}^{1,\infty}(\Omega;\mathrm{M}_r(\mathbf{C})), \ k \in 1..d, \ \mathrm{and} \ \mathbf{C} \in \mathrm{L}^\infty(\Omega;\mathrm{M}_r(\mathbf{C}))$$

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$$\mathcal{L}u=f$$

the symmetric positive system or the Friedrichs system.

Symmetric hyperbolic systems (KOF1954)

$$\sum_{k=1}^{d} \mathbf{A}^k \partial_k \mathbf{u} + \mathbf{B} \mathbf{u} = \mathbf{f} .$$

In divergence form:

$$\sum_{k=1}^d \partial_k (\mathbf{A}^k \mathsf{u}) + (\mathbf{B} - \partial_k \mathbf{A}^k) \mathsf{u} = \mathsf{f} \; .$$

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It is symmetric if all matrices \mathbf{A}^k are real and symmetric; and uniformly hyperbolic if there is a $\boldsymbol{\xi} \in \mathbf{R}^d$ such that for any $\mathbf{x} \in \mathsf{Cl}\,\Omega$ the matrix $\xi_k \mathbf{A}^k(\mathbf{x})$ is positive definite.

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Such systems can easily be transformed into Friedrichs' systems.

It is known that the wave equation and the Maxwell system can be written as an equivalent hyperbolic system.

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K. O. Friedrichs: Symmetric positive linear differential equations, Communications on Pure and Applied Mathematics 11 (1958), 333–418

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- treating the equations of mixed type, such as the Tricomi equation:

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The development of theory is nowadays mostly motivated by the needs in development of numerical methods.

An example - scalar elliptic equation

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which is a Friedrichs system with the choice of

$$\mathbf{A}_1 = \begin{bmatrix} 0 & 0 & 1 \\ 0 & 0 & 0 \\ 1 & 0 & 0 \end{bmatrix}, \quad \mathbf{A}_2 = \begin{bmatrix} 0 & 0 & 0 \\ 0 & 0 & 1 \\ 0 & 1 & 0 \end{bmatrix}, \quad \mathbf{C} = \begin{bmatrix} 1 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & \mu \end{bmatrix}.$$

... with zero initial and Dirichlet boundary condition:

$$\left\{ \begin{array}{l} \partial_t u - \mathrm{div}_{\,\mathbf{x}}(\mathbf{A} \nabla_{\mathbf{x}} u) + \mathbf{b} \cdot \nabla_{\mathbf{x}} u + c u = f \ \text{in} \ \Omega_T \\ \\ u = 0 \ \text{on} \ \partial \Omega \times \langle 0, T \rangle \\ \\ u(\cdot, 0) = 0 \ \text{on} \ \Omega \end{array} \right.$$

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The condition (F1) holds. The positivity condition $\mathbf{C} + \mathbf{C}^{\top} \geqslant 2\mu_0 \mathbf{I}$ is fulfilled if and only if $c - \frac{1}{4}\mathbf{A}^{-1}\mathbf{b} \cdot \mathbf{b}$ is uniformly positive.

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allows the treatment of different types of usual boundary conditions.

Assumptions on the boundary matrix ${f M}$

We assume (for ae $\mathbf{x} \in \Gamma$) [KOF1958] $(\mathsf{FM1}) \qquad (\forall \, \pmb{\xi} \in \mathbf{C}^r) \quad (\mathbf{M}(\mathbf{x}) + \mathbf{M}(\mathbf{x})^*) \pmb{\xi} \cdot \pmb{\xi} \geqslant 0 \, ,$

Assumptions on the boundary matrix M

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Such \mathbf{M} is called *the admissible boundary condition*.

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The boundary problem: for given $\mathbf{f} \in \mathrm{L}^2(\Omega; \mathbf{C}^r)$ find u such that

$$\left\{ \begin{aligned} \mathcal{L}u &= f \\ (\mathbf{A}_{\boldsymbol{\nu}} - \mathbf{M})u_{\big|_{\Gamma}} &= 0 \end{aligned} \right. .$$

$$\mathbf{M} \qquad \mathbf{A}_{\nu} - \mathbf{M} \qquad (\mathbf{A}_{\nu} - \mathbf{M}) \begin{bmatrix} \mathbf{p} \\ u \end{bmatrix}_{\mid \Gamma} = \mathbf{0}$$

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All above matrices M satisfy (FM).

Different ways to enforce boundary conditions

Instead of

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 on Γ ,

Lax proposed boundary conditions with

$$\mathbf{u}(\mathbf{x}) \in N(\mathbf{x}) \,, \quad \mathbf{x} \in \Gamma \,,$$

where $N=\{N(\mathbf{x}):\mathbf{x}\in\Gamma\}$ is a family of subspaces of $\mathbf{C}^r.$

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Boundary problem:

$$\begin{cases} \mathcal{L}\mathbf{u} = \mathbf{f} \\ \mathbf{u}(\mathbf{x}) \in N(\mathbf{x}), & \mathbf{x} \in \Gamma \end{cases}.$$

Assumptions on N

 $(\mathsf{FX1}) \begin{tabular}{ll} \textit{M}(\mathbf{x}) & \text{is non-negative with respect to } \mathbf{A}_{\nu}(\mathbf{x}): \\ & (\forall \, \pmb{\xi} \in N(\mathbf{x})) \quad \mathbf{A}_{\nu}(\mathbf{x}) \pmb{\xi} \cdot \pmb{\xi} \geqslant 0 \, ; \\ \\ & (\mathsf{FX2}) \begin{tabular}{ll} \textit{K}(\mathbf{x}) & \mathbf{A}_{\nu}(\mathbf{x}) & \mathbf{A}$

Assumptions on N

maximal boundary conditions: (for ae $x \in \Gamma$) [PDL] $N(\mathbf{x})$ is non-negative with respect to $\mathbf{A}_{\nu}(\mathbf{x})$: (FX1) $(\forall \boldsymbol{\xi} \in N(\mathbf{x})) \quad \mathbf{A}_{\boldsymbol{\nu}}(\mathbf{x})\boldsymbol{\xi} \cdot \boldsymbol{\xi} \geqslant 0$: there is no non-negative subspace with respect to (FX2) $\mathbf{A}_{\nu}(\mathbf{x})$, which contains $N(\mathbf{x})$; [RSP&LS1966] or Let $N(\mathbf{x})$ and $\tilde{N}(\mathbf{x}) := (\mathbf{A}_{\nu}(\mathbf{x})N(\mathbf{x}))^{\perp}$ satisfy (for ae $\mathbf{x} \in \Gamma$) $(\forall \boldsymbol{\xi} \in N(\mathbf{x})) \quad \mathbf{A}_{\nu}(\mathbf{x})\boldsymbol{\xi} \cdot \boldsymbol{\xi} \geqslant 0$ (FV1) $(\forall \boldsymbol{\xi} \in \tilde{N}(\mathbf{x})) \quad \mathbf{A}_{\nu}(\mathbf{x})\boldsymbol{\xi} \cdot \boldsymbol{\xi} \leqslant 0$ $\tilde{N}(\mathbf{x}) = (\mathbf{A}_{\nu}(\mathbf{x})N(\mathbf{x}))^{\perp}$ and $N(\mathbf{x}) = (\mathbf{A}_{\nu}(\mathbf{x})\tilde{N}(\mathbf{x}))^{\perp}$. (FV2)

Equivalence of different descriptions of boundary conditions

Theorem. It holds

$$(FM1)$$
– $(FM2)$ \iff $(FX1)$ – $(FX2)$ \iff $(FV1)$ – $(FV2)$,

with

$$N(\mathbf{x}) := \ker \left(\mathbf{A}_{\nu}(\mathbf{x}) - \mathbf{M}(\mathbf{x}) \right).$$

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In fact, for a weak existence result some additional assumptions are needed [JR1994], [MJ2004].

Classical results on well-posedness

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Contributions:

- C. Morawetz, P. Lax, L. Sarason, R. S. Phillips, J. Rauch, ...
- the meaning of traces for functions in the graph space
- weak well-posedness results under additional assumptions (on ${f A}_
 u)$
- regularity of solution
- numerical treatment

Classical theory

What are Friedrichs systems? Examples Boundary conditions for Friedrichs systems Existence, uniqueness, well-posedness

Abstract formulation

Graph spaces
Cone formalism of Ern, Guermond and Caplain
Interdependence of different representations of boundary conditions
Kreĭn spaces
Equivalence of boundary conditions

What can we say for the Friedrichs operator now?

Some examples Two-field theory

Concluding remarks

A. Ern, J.-L. Guermond, G. Caplain: An intrinsic criterion for the bijectivity of Hilbert operators related to Friedrichs' systems, Comm. Partial Diff. Eq. 32 (2007) 317–341.

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...and new open questions.

They considered only the real case.

L — real (complex) Hilbert space (L' is (anti)dual of L), $\mathcal{D}\subseteq L$ — dense subspace,

```
\begin{array}{l} L \longrightarrow {\sf real} \ ({\sf complex}) \ {\sf Hilbert} \ {\sf space} \ (L' \ {\sf is} \ ({\sf anti}) {\sf dual} \ {\sf of} \ L), \\ \mathcal{D} \subseteq L \longrightarrow {\sf dense} \ {\sf subspace}, \\ T, \tilde{T} : \mathcal{D} \longrightarrow L \longrightarrow {\sf linear} \ {\sf unbounded} \ {\sf operators} \ {\sf satisfying} \end{array}
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The Friedrichs operator

Let
$$\mathcal{D}:=\mathrm{C}_c^\infty(\Omega;\mathbf{C}^r)$$
, $L=\mathrm{L}^2(\Omega;\mathbf{C}^r)$ and $T,\tilde{T}:\mathcal{D}\longrightarrow L$ be defined by
$$T\mathrm{u}:=\sum_{k=1}^d\partial_k(\mathbf{A}_k\mathrm{u})+\mathbf{C}\mathrm{u}\,,$$

$$\tilde{T}\mathrm{u}:=-\sum_{k=1}^d\partial_k(\mathbf{A}_k\mathrm{u})+(\mathbf{C}^*+\sum_{k=1}^d\partial_k\mathbf{A}_k)\mathrm{u}\,,$$

where A_k and C are as above (they satisfy (F1)–(F2)).

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... fits in this framework.

 $(\mathcal{D}, \langle \, \cdot \, | \, \cdot \, \rangle_T)$ is an inner product space, where

$$\langle \cdot | \cdot \rangle_T := \langle \cdot | \cdot \rangle_L + \langle T \cdot | T \cdot \rangle_L.$$

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Abusing notation: $T, \tilde{T} \in \mathcal{L}(L; W_0') \dots \text{(T1)-(T3)}$

Formulation of the problem

Lemma. The graph space

$$W := \{ u \in L : Tu \in L \} = \{ u \in L : \tilde{T}u \in L \},\$$

is a Hilbert space with respect to $\langle \cdot | \cdot \rangle_T$.

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Problem: for given $f \in L$ find $u \in W$ such that Tu = f.

Find sufficient conditions on $V\leqslant W$ such that $T_{\big|V}:V\longrightarrow L$ is an isomorphism.

Boundary operator

Boundary operator $D \in \mathcal{L}(W; W')$:

$${}_{W'}\!\langle\, Du,v\,\rangle_W := \langle\, Tu\mid v\,\rangle_L - \langle\, u\mid \tilde{T}v\,\rangle_L\,, \qquad u,v\in W\,.$$

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Lemma. *D* is selfadjoint

$$W'\langle Du, v \rangle_W = \overline{W'\langle Dv, u \rangle_W}$$

and satisfies

$$\ker D = W_0$$

$$\operatorname{im} D = W_0^0 := \left\{ g \in W' : (\forall u \in W_0) \quad {}_{W'} \langle g, u \rangle_W = 0 \right\}.$$

In particular, im D is closed in W'.

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In particular, im D is closed in W'.

If T is the Friedrichs operator \mathcal{L} , then for $\mathbf{u},\mathbf{v}\in \mathrm{C}^\infty_c(\mathbf{R}^d;\mathbf{C}^r)$ we have

$$W'\langle D\mathbf{u}, \mathbf{v} \rangle_W = \int_{\Gamma} \mathbf{A}_{\nu}(\mathbf{x}) \mathbf{u}_{\mid \Gamma}(\mathbf{x}) \cdot \mathbf{v}_{\mid \Gamma}(\mathbf{x}) dS(\mathbf{x}).$$

Well-posedness theorem

Let V and \tilde{V} be subspaces of W that satisfy

$$(\forall u \in V) \qquad \qquad _{W'} \langle Du, u \rangle_W \geqslant 0$$

$$(\forall v \in \tilde{V}) \qquad \qquad _{W'} \langle Dv, v \rangle_W \leqslant 0$$

$$(\forall v) \qquad \qquad V = D(\tilde{V})^0 \,, \qquad \tilde{V} = D(V)^0 \,.$$

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(cone formalism)

Theorem. Under assumptions (T1)-(T3) and (V1)-(V2), the operators $T_{|_{V}}:V\longrightarrow L$ and $\tilde{T}_{|_{\tilde{V}}}:\tilde{V}\longrightarrow L$ are isomorphisms.

In the real case [AE&JLG&GC2007].

Correspondence with classical assumptions

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Correspondence with *classical* assumptions

$$(\forall 1) \qquad (\forall u \in V) \qquad {}_{W'}\langle Du, u \rangle_W \geqslant 0 \, , \\ (\forall v \in \tilde{V}) \qquad {}_{W'}\langle Dv, v \rangle_W \leqslant 0 \, , \\ (\forall 2) \qquad V = D(\tilde{V})^0 \, , \qquad \tilde{V} = D(V)^0 \, , \\ (\forall \xi \in N(\mathbf{x})) \qquad \mathbf{A}_{\nu}(\mathbf{x}) \boldsymbol{\xi} \cdot \boldsymbol{\xi} \geqslant 0 \, , \\ (\forall \xi \in \tilde{N}(\mathbf{x})) \qquad \mathbf{A}_{\nu}(\mathbf{x}) \boldsymbol{\xi} \cdot \boldsymbol{\xi} \leqslant 0 \, , \\ (\forall \nabla \mathbf{x}) \qquad \tilde{N}(\mathbf{x}) = (\mathbf{A}_{\nu}(\mathbf{x})N(\mathbf{x}))^{\perp} \quad \text{and} \quad N(\mathbf{x}) = (\mathbf{A}_{\nu}(\mathbf{x})\tilde{N}(\mathbf{x}))^{\perp} \, , \\ (\text{for ae } \mathbf{x} \in \Gamma)$$

Other sets of conditions in the classical setting (recall)

maximal boundary conditions: (for ae $\mathbf{x} \in \Gamma$)

(FX1)
$$(\forall \boldsymbol{\xi} \in N(\mathbf{x})) \quad \mathbf{A}_{\nu}(\mathbf{x})\boldsymbol{\xi} \cdot \boldsymbol{\xi} \geqslant 0 \,,$$

(FX2) there is no non-negative subspace with respect to $\mathbf{A}_{\nu}(\mathbf{x}), \text{ which contains } N(\mathbf{x})\,,$

admissible boundary conditions: there exists a matrix function $\mathbf{M}:\Gamma\longrightarrow \mathrm{M}_r(\mathbf{C})$ such that (for ae $\mathbf{x}\in\Gamma$)

(FM1)
$$(\forall \boldsymbol{\xi} \in \mathbf{C}^r) \quad (\mathbf{M}(\mathbf{x}) + \mathbf{M}(\mathbf{x})^*) \boldsymbol{\xi} \cdot \boldsymbol{\xi} \geqslant 0,$$

(FM2)
$$\mathbf{C}^r = \ker \left(\mathbf{A}_{\nu}(\mathbf{x}) - \mathbf{M}(\mathbf{x}) \right) + \ker \left(\mathbf{A}_{\nu}(\mathbf{x}) + \mathbf{M}(\mathbf{x}) \right).$$

Correspondence — maximal b.c.

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$$\mathbf{A}_{\nu}(\mathbf{x})$$
, which contains $N(\mathbf{x})$,

subspace V is maximal non-negative with respect to D:

- $\text{(X1)} \quad V \text{ is non-negative with respect to } D \colon \quad (\forall \, v \in V) \quad {}_{W'}\!\langle \, Dv, v \, \rangle_{W} \geqslant 0 \,,$
- (X2) there is no non-negative subspace with respect to ${\cal D}$ that contains ${\cal V}$.

Correspondence — admissible b.c.

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admissible boundary condition: there exist $M \in \mathcal{L}(W;W')$ that satisfy

(M1)
$$(\forall u \in W) \quad _{W'} \langle (M+M^*)u, u \rangle_W \geqslant 0,$$

(M2)
$$W = \ker(D - M) + \ker(D + M).$$

Equivalence of different descriptions of b.c.

Theorem. (classical) It holds
$$(\textit{FM1})\text{-}(\textit{FM2}) \iff (\textit{FV1})\text{-}(\textit{FV2}) \iff (\textit{FX1})\text{-}(\textit{FX2})\,,$$
 with
$$N(\mathbf{x}) := \ker \Big(\mathbf{A}_{\nu}(\mathbf{x}) - \mathbf{M}(\mathbf{x})\Big)\,.$$

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Theorem. (A. Ern, J.-L. Guermond, G. Caplain) It holds
$$(M1)\text{-}(M2) \ \stackrel{\Longrightarrow}{\longleftarrow} \ (V1)\text{-}(V2) \ \Longrightarrow \ (X1)\text{-}(X2)\,,$$
 with
$$V:=\ker(D-M)\,.$$

This was obtained in the real case only.

$$(M1)$$
– $(M2)$ \leftarrow $(V1)$ – $(V2)$

Theorem. Let V and \tilde{V} satisfy (V1)–(V2), and suppose that there exist operators $P \in \mathcal{L}(W;V)$ and $Q \in \mathcal{L}(W;\tilde{V})$ such that

$$\begin{split} (\forall\,v\in V) \quad D(v-Pv) &= 0\,,\\ (\forall\,v\in \tilde{V}) \quad D(v-Qv) &= 0\,,\\ DPQ &= DQP\,. \end{split}$$

Let us define $M \in \mathcal{L}(W; W')$ (for $u, v \in W$) with

$$W'\langle Mu, v \rangle_W = W'\langle DPu, Pv \rangle_W - W'\langle DQu, Qv \rangle_W + W'\langle D(P+Q-PQ)u, v \rangle_W - W'\langle Du, (P+Q-PQ)v \rangle_W.$$

Then $V := \ker(D - M)$, $\tilde{V} := \ker(D + M^*)$, and M satisfies (M1)–(M2).

New notation

$$[\,u\mid v\,] := {}_{W'}\!\langle\, Du,v\,\rangle_W = \langle\, Tu\mid v\,\rangle_L - \langle\, u\mid \tilde{T}v\,\rangle_L\,, \qquad u,v\in W$$
 is an indefinite inner product on $W.$

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$$(\forall v \in V) \quad \begin{bmatrix} v \mid v \end{bmatrix} \geqslant 0 \, ,$$

$$(\forall v \in \tilde{V}) \quad \begin{bmatrix} v \mid v \end{bmatrix} \leqslant 0 \, ;$$

$$({\rm V2}) \hspace{1cm} V = \tilde{V}^{[\perp]} \,, \qquad \tilde{V} = V^{[\perp]} \,. \label{eq:V2}$$

($^{[\bot]}$ stands for $[\,\cdot\,\,|\,\,\cdot\,] \text{-orthogonal complement})$

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$$(\forall v \in V) \quad [v \mid v] \geqslant 0,$$

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($^{[\perp]}$ stands for $[\cdot | \cdot]$ -orthogonal complement) subspace V is maximal non-negative in $(W, [\cdot | \cdot])$:

- $\text{(X1)} \qquad V \text{ is non-negative in } (W,[\,\cdot\,\,|\,\,\cdot\,]) \colon \quad (\forall\,v\in V) \quad [\,v\,\,|\,\,v\,] \geqslant 0\,,$
- (X2) there is no non-negative subspace in $(W,[\,\cdot\,\,|\,\,\cdot\,\,])$ containing V .

 $(W, \lceil \cdot \mid \cdot \rceil)$ is not a Kreı̆n space – it is a degenerate space, because its Gramm operator $G := j \circ D$ $(j: W' \longrightarrow W \text{ is the canonical isomorphism})$ has large kernel:

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Theorem. If G is the Gramm operator of the space W, then the quotient space $\hat{W} := W/\ker G$ is a Kreĭn space if and only if $\operatorname{im} G$ is closed.

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 $\hat{W} := W/W_0$ is the Kreĭn space, with

$$\left[\, \hat{u} \mid \hat{v} \, \right] \, := \left[\, u \mid v \, \right], \qquad u,v \in W \, .$$

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.

Theorem. If G is the Gramm operator of the space W, then the quotient space $\hat{W}:=W/\ker G$ is a Kreĭn space if and only if $\operatorname{im} G$ is closed.

 $\hat{W} := W/W_0$ is the Kreĭn space, with

$$\left[\, \hat{u} \mid \hat{v} \, \right] \hat{} := \left[\, u \mid v \, \right], \qquad u,v \in W \, .$$

Important: im D is closed and $\ker D = W_0$.

Quotient Kreĭn space

Lemma. Let $U \supseteq W_0$ and Y be subspaces of W. Then

- a) U is closed if and only if $\hat{U}:=\{\hat{v}:v\in U\}$ is closed in \hat{W} ;
- b) $\widehat{(U+Y)} = \{u+v+W_0 : u \in U, v \in Y\} = \hat{U} + \hat{Y};$
- c) U+Y is closed if and only if $\hat{U}+\hat{Y}$ is closed;
- d) $(\hat{Y})^{[\perp]} = \widehat{Y^{[\perp]}}$.
- e) if Y is maximal non-negative (non-positive) in W, than \hat{Y} is maximal non-negative (non-positive) in \hat{W} ;
- f) if \hat{U} is maximal non-negative (non-positive) in \hat{W} , then U is maximal non-negative (non-positive) in W.

$$(V1)$$
- $(V2)$ \iff $(X1)$ - $(X2)$

Theorem. a) If subspaces V and \tilde{V} satisfy (V1)–(V2), then V is maximal non-negative in W (satisfies (X1)–(X2)) and \tilde{V} is maximal non-positive in W.

b) If V is maximal non-negative in W, then V and $\tilde{V}:=V^{[\perp]}$ satisfy (V1)–(V2).

$$(M1)–(M2) \implies (V1)–(V2) \qquad (recall)$$

Theorem. [EGC] (T1)–(T3) and
$$M \in \mathcal{L}(W;W')$$
 satisfy (M) imply
$$V := \ker(D-M) \quad \text{and} \quad \tilde{V} := \ker(D+M^*) \quad \text{satisfy (V)}.$$

Corollary. Under above assumptions

$$T_{\big|\ker(D-M)}:\ker(D-M)\longrightarrow L \qquad \text{i}\qquad \tilde{T}_{\big|\ker(D+M^*)}:\ker(D+M^*)\longrightarrow L$$
 are isomorphisms.

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$$(M1)-(M2) \quad \longleftarrow \quad (V1)-(V2) \qquad (recall)$$

Theorem. Let V and \tilde{V} satisfy (V1)–(V2), and suppose that there exist operators $P \in \mathcal{L}(W;V)$ and $Q \in \mathcal{L}(W;\tilde{V})$ such that

$$(\forall v \in V) \quad D(v - Pv) = 0,$$
$$(\forall v \in \tilde{V}) \quad D(v - Qv) = 0,$$
$$DPQ = DQP.$$

Let us define $M \in \mathcal{L}(W; W')$ (for $u, v \in W$) with

$$W'\langle Mu, v \rangle_W = W'\langle DPu, Pv \rangle_W - W'\langle DQu, Qv \rangle_W + W'\langle D(P+Q-PQ)u, v \rangle_W - W'\langle Du, (P+Q-PQ)v \rangle_W.$$

Then $V := \ker(D - M)$, $\tilde{V} := \ker(D + M^*)$, and M satisfies (M1)–(M2).

$$(M1)$$
– $(M2)$ \iff $(V1)$ – $(V2)$ (direct proof)

Theorem. If V, \tilde{V} are two closed subspaces of W that satisfy $W_0 \subseteq V \cap \tilde{V}$, then the following statements are equivalent:

a) There exist operators $P \in \mathcal{L}(W; V)$ and $Q \in \mathcal{L}(W; \tilde{V})$, such that

$$\begin{split} (\forall\,v\in V) \quad &D(v-Pv)=0\,,\\ (\forall\,v\in \tilde{V}) \quad &D(v-Qv)=0\,,\\ &DPQ=DQP\,. \end{split}$$

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b) There exist projectors $P',Q'\in\mathcal{L}(W;W)$, such that

$$P^{\prime ^2}=P^\prime$$
 and ${Q^\prime}^2=Q^\prime\,,$
$$\operatorname{im} P^\prime=V \quad \text{and} \quad \operatorname{im} Q^\prime=\tilde{V}\,,$$

$$P^\prime Q^\prime=Q^\prime P^\prime\,.$$

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$$P'^2=P'$$
 and ${Q'}^2=Q'$,
 $\operatorname{im} P'=V$ and $\operatorname{im} Q'=\tilde{V}$,
 $P'Q'=Q'P'$.

(b) is equivalent to closedness of $V+\tilde{V}.$

$$(M1)$$
- $(M2) \iff (V1)$ - $(V2) \pmod{.}$

Theorem.

a) $V, \ddot{V} \leqslant W$ satisfy (V), and exists a closed subspace $W_2 \subseteq C^-$ of W, $V \dot{+} W_2 = W$, then there exist an operator $M \in \mathcal{L}(W;W')$ satisfying (M) and $V = \ker(D-M)$.

If we define W_1 as orthogonal complement of W_0 in V, so that $W=W_1\dot{+}W_0\dot{+}W_2$, and denote by R_1,R_0,R_2 projectors that correspond to above direct sum, then one such operator is given with $M=D(R_1-R_2)$.

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b) $M \in \mathcal{L}(W; W')$ an operator satisfying (M1)–(M2), $V := \ker(D-M)$. For W_2 , the orthogonal complement of W_0 in $\ker(D+M)$, $W_2 \subseteq C^-$ is closed, $V \dot{+} W_2 = W$, and M coincide with the operator in (a).

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Lemma. Let $W_2'' \leqslant W$ satisfies $W_2'' \subseteq C^-$ and $W_2'' + V = W$. Then there is a closed subspace W_2 of W, such that $W_2 \subseteq C^-$ and $W_2 \dot{+} V = W$.

$$(M1)$$
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Lemma. If $U_1 + U_2 = W$ for some subspaces $U_1 \subseteq C^+$ and $U_2 \subseteq C^-$ of W, then $U_1 \cap U_2 \subseteq W_0$.

If additionally U_1 is maximal nonnegative and U_2 maximal nonpositive, then $U_1 \cap U_2 = W_0$.

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Theorem. For a maximal nonnegative subspace V of W, it is equivalent:

- a) There is a maximal nonpositive subspace W_2 of W, such that $W_2 + V = W$;
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Corollary. The conditions (V) and (M) are equivalent.

Classical theory

What are Friedrichs systems? Examples Boundary conditions for Friedrichs systems Existence, uniqueness, well-posedness

Abstract formulation

Graph spaces
Cone formalism of Ern, Guermond and Caplain
Interdependence of different representations of boundary conditions
Kreĭn spaces
Equivalence of boundary conditions

What can we say for the Friedrichs operator now?

Some examples Two-field theory

Concluding remarks

Scalar elliptic equation

Consider

$$-\mathsf{div}\left(\mathbf{A}\nabla u\right) + cu = f$$

in $\Omega \subseteq \mathbf{R}^d$,

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$$\mathbf{A}\in\mathcal{M}_d(\alpha',\beta';\Omega):=\left\{\mathbf{A}\in\mathrm{L}^\infty(\Omega;\mathrm{M}_d(\mathbf{R})): \left(\forall\,\pmb{\xi}\in\mathbf{R}^d\right)\,\mathbf{A}\pmb{\xi}\cdot\pmb{\xi}\geq\alpha'|\pmb{\xi}|^2\,\&\,\mathbf{A}\pmb{\xi}\cdot\pmb{\xi}\geq\frac{1}{\beta'}|\mathbf{A}\pmb{\xi}|^2\right\}$$

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$$\mathbf{A} \in \mathcal{M}_d(\alpha', \beta'; \Omega) := \left\{ \mathbf{A} \in L^{\infty}(\Omega; M_d(\mathbf{R})) : \right.$$
$$(\forall \boldsymbol{\xi} \in \mathbf{R}^d) \ \mathbf{A} \boldsymbol{\xi} \cdot \boldsymbol{\xi} \ge \alpha' |\boldsymbol{\xi}|^2 \ \& \ \mathbf{A} \boldsymbol{\xi} \cdot \boldsymbol{\xi} \ge \frac{1}{\beta'} |\mathbf{A} \boldsymbol{\xi}|^2 \right\}$$

New unknown vector function taking values in \mathbf{R}^{d+1} :

$$\mathbf{u} = \begin{bmatrix} \mathbf{u}_d \\ u_{d+1} \end{bmatrix} = \begin{bmatrix} -\mathbf{A} \nabla_{\mathbf{x}} u \\ u \end{bmatrix} .$$

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Then the starting equation can be written as a first-order system

$$\begin{cases} \nabla_{\mathbf{x}} u_{d+1} + \mathbf{A}^{-1} \mathbf{u}_d = \mathbf{0} \\ \operatorname{div} \mathbf{u}_d + c u_{d+1} = f \end{cases},$$

Scalar elliptic equation (cont.)

which is a Friedrichs system with the choice of

$$\mathbf{A}_k = \mathbf{e}_k \otimes \mathbf{e}_{d+1} + \mathbf{e}_{d+1} \otimes \mathbf{e}_k \in \mathcal{M}_{d+1}(\mathbf{R}), \qquad \mathbf{C} = \begin{bmatrix} \mathbf{A}^{-1} & \mathbf{0} \\ \mathbf{0} & c \end{bmatrix}.$$

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The graph space: $W = L^2_{\mathrm{div}}(\Omega) \times \mathrm{H}^1(\Omega)$.

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The graph space: $W = L^2_{div}(\Omega) \times H^1(\Omega)$.

Dirichlet, Neumann and Robin boundary conditions are imposed by the following choice of V and $\widetilde{V}\colon$

$$\begin{split} V_D &= \widetilde{V}_D := & \mathbf{L}_{\mathrm{div}}^2(\Omega) \times \mathbf{H}_0^1(\Omega) \,, \\ V_N &= \widetilde{V}_N := & \left\{ (\mathbf{u}_d, u_{d+1})^\top \in W : \boldsymbol{\nu} \cdot \mathbf{u}_d = 0 \right\}, \\ V_R &:= & \left\{ (\mathbf{u}_d, u_{d+1})^\top \in W : \boldsymbol{\nu} \cdot \mathbf{u}_d = a u_{d+1}_{\left| \Gamma \right|} \right\}, \\ \widetilde{V}_R &:= & \left\{ (\mathbf{u}_d, u_{d+1})^\top \in W : \boldsymbol{\nu} \cdot \mathbf{u}_d = -a u_{d+1}_{\left| \Gamma \right|} \right\}. \end{split}$$

Heat equation

... with zero initial and Dirichlet boundary condition:

$$\left\{ \begin{array}{l} \partial_t u - \mathrm{div}_{\,\mathbf{x}} (\mathbf{A} \nabla_{\mathbf{x}} u) + \mathbf{b} \cdot \nabla_{\mathbf{x}} u + c u = f \ \text{in} \ \Omega_T \\ \\ u = 0 \ \text{on} \ \partial \Omega \times \langle 0, T \rangle \\ \\ u(\cdot, 0) = 0 \ \text{on} \ \Omega \end{array} \right.$$

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...as a Friedrichs system:

$$\left\{ \begin{aligned} &\nabla_{\mathbf{x}} u_{d+1} + \mathbf{A}^{-1} \mathbf{u}_d = \mathbf{0} \\ &\partial_t u_{d+1} + \operatorname{div}_{\mathbf{x}} \mathbf{u}_d + c u_{d+1} - \mathbf{A}^{-1} \mathbf{b} \cdot \mathbf{u}_d = f \end{aligned} \right. ,$$

(note that we use $\mathbf{u} = (\mathbf{u}_d, u_{d+1})^{\top}$).

Friedrichs operator and the graph space

The operator T is given by

$$T \begin{bmatrix} \mathbf{u}_d \\ u_{d+1} \end{bmatrix} = \begin{bmatrix} \nabla_{\mathbf{x}} u_{d+1} + \mathbf{A}^{-1} \mathbf{u}_d \\ \partial_t u_{d+1} + \operatorname{div}_{\mathbf{x}} \mathbf{u}_d + c u_{d+1} - \mathbf{A}^{-1} \mathbf{b} \cdot \mathbf{u}_d \end{bmatrix} \,,$$

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while the corresponding graph space is

$$\begin{split} W &= \left\{ \mathbf{u} \in \mathbf{L}^2(\Omega_T; \mathbf{R}^{d+1}) : \nabla_{\mathbf{x}} u_{d+1} \in \mathbf{L}^2(\Omega_T; \mathbf{R}^d) \right. \\ &\qquad \qquad \& \quad \partial_t u_{d+1} + \mathsf{div}_{\,\mathbf{x}} \mathbf{u}_d \in \mathbf{L}^2(\Omega_T) \right\} \\ &= \left\{ \mathbf{u} \in \mathbf{L}^2_{\mathrm{div}}(\Omega_T) : \nabla_{\mathbf{x}} u_{d+1} \in \mathbf{L}^2(\Omega_T; \mathbf{R}^d) \right\} \\ &= \left\{ \mathbf{u} \in \mathbf{L}^2_{\mathrm{div}}(\Omega_T) : u_{d+1} \in \mathbf{L}^2(0, T; \mathbf{H}^1(\Omega)) \right\}. \end{split}$$

Properties of the last component

Lemma. The projection $\mathbf{u}=(\mathbf{u}_d,u_{d+1})^{\top}\mapsto u_{d+1}$ is a continuous linear operator from W to W(0,T), which is continuously embedded to $\mathrm{C}([0,T];\mathrm{L}^2(\Omega))$.

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The space

$$W(0,T) = \left\{ u \in L^{2}(0,T; H^{1}(\Omega)) : \partial_{t}u \in L^{2}(0,T; H^{-1}(\Omega)) \right\},$$

is a Banach space when equipped by norm

$$\|\mathbf{u}\|_{W(0,T)} = \sqrt{\|u\|_{\mathrm{L}^2(0,T;\mathrm{H}^1(\Omega))}^2 + \|\partial_t u\|_{\mathrm{L}^2(0,T;\mathrm{H}^{-1}(\Omega))}^2}$$
.

Let

$$\begin{split} V &= \left\{ \mathbf{u} \in W : u_{d+1} \in \mathrm{L}^2(0,T;\mathrm{H}^1_0(\Omega)), \quad u_{d+1}(\cdot,0) = 0 \text{ a.e. on } \Omega \right\}, \\ \widetilde{V} &= \left\{ \mathbf{v} \in W : v_{d+1} \in \mathrm{L}^2(0,T;\mathrm{H}^1_0(\Omega)), \quad v_{d+1}(\cdot,T) = 0 \text{ a.e. on } \Omega \right\}. \end{split}$$

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Do they satisfy (V1)–(V2)?

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Do they satisfy (V1)–(V2)? Technical...

Theorem

The above V and \widetilde{V} satisfy (V1)–(V2), and therefore the operator $T_{|_{V}}:V\longrightarrow L$ is an isomorphism.

Heat equation with b = 0 and c = 0:

$$\left\{ \begin{array}{l} \partial_t u - \mathrm{div}_{\,\mathbf{x}}(\mathbf{A} \nabla_{\mathbf{x}} u) = f \text{ in } \Omega_T \\ \\ u = 0 \text{ on } \Gamma \times \langle 0, T \rangle \\ \\ u(\cdot, 0) = 0 \text{ on } \Omega \end{array} \right.$$

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Two field theory:

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matrices need to be of the form

$$\mathbf{A}^k = \begin{bmatrix} \mathbf{0} & \mathsf{B}^k \\ (\mathsf{B}^k)^\top & a^k \end{bmatrix} \quad \text{and} \quad \mathbf{C} = \begin{bmatrix} \mathbf{C}^d & \mathbf{0} \\ \mathbf{0}^\top & c^{d+1} \end{bmatrix} \;,$$

where $\mathsf{B}^k \in \mathbf{R}^d$ are constant vectors, $a^k \in \mathrm{W}^{1,\infty}(\Omega_T)$, $\mathbf{C}^d \in \mathrm{L}^\infty(\Omega_T; \mathrm{M}_d(\mathbf{R}))$ and $c^{d+1} \in \mathrm{L}^\infty(\Omega_T)$, $k \in 1..(d+1)$.

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For the heat equation matrices have this form!

Instead of coercivity (positivity) condition (F2), the following is required:

$$\begin{split} (\exists\, \mu_1 > 0) (\forall\, \pmb{\xi} &= (\pmb{\xi}_d, \xi_{d+1}) \in \mathbf{R}^{d+1}) \\ \Big(\mathbf{C} + \mathbf{C}^\top + \sum_{k=1}^{d+1} \partial_k \mathbf{A}_k \Big) \pmb{\xi} \cdot \pmb{\xi} \geqslant 2\mu_1 |\pmb{\xi}_d|^2 \qquad \text{ (a.e. on } \Omega)\,, \end{split}$$

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where $Bu_{d+1} := \sum_{k=1}^{d+1} B^k \partial_k u_{d+1} = \nabla_{\mathbf{x}} u_{d+1}$.

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For our system both conditions are trivially fulfilled.

Therefore, we have the well-posedness result.

Classical theory

What are Friedrichs systems? Examples Boundary conditions for Friedrichs systems Existence, uniqueness, well-posedness

Abstract formulation

Graph spaces
Cone formalism of Ern, Guermond and Caplain
Interdependence of different representations of boundary conditions
Kreĭn spaces
Equivalence of boundary conditions

What can we say for the Friedrichs operator now?

Some examples Two-field theory

Concluding remarks

Some further applications . . .

Already known:

- non-stationary theory the Maxwell system, non-stationary div-grad problem, the wave equation
- homogenisation

Some further applications . . .

Already known:

- non-stationary theory
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- homogenisation

Work in progress:

- Dirac system
- time-harmonic Maxwell system
- what can be done for the Schrödinger equation?

Open problems ...

- Find all representations of a particular equation in the form of a Friedrichs system.
- Application to other equations of practical importance (mixed-type problems).
- Compare the results to those already known in the classical setting.

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Some used properties

Theorem. a) $[\cdot | \cdot]$ -orthogonal complement of a maximal non-negative (non-positive) subspace is non-positive (non-negative).

- b) Each maximal semi-definite subspace contains all isotropic vectors in W.
- c) If L is a non-negative (non-positive) subspace of a Krein space, such that $L^{[\perp]}$ is non-positive (non-negative), then $\operatorname{Cl} L$ is maximal non-negative (non-positive).
- d) Each maximal semi-definite subspace of a Krein space is closed.
- e) A subspace L of a Krein space is closed if and only if $L = L^{[\perp][\perp]}$.
- f) For a subspace L of a Krein space W it holds

$$L\cap L^{[\perp]}=\{0\}\qquad\Longleftrightarrow\qquad \operatorname{Cl}\,(L+L^{[\perp]})=W\,.$$

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