### Friedrichs systems with complex coefficients

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SISSA, Trieste, 27th October 2016

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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  $\Gamma$ ;

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

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The operator  $\mathcal{L}: L^2(\Omega; \mathbf{C}^r) \longrightarrow \mathcal{D}'(\Omega; \mathbf{C}^r)$ 

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

the symmetric positive system or the Friedrichs system.

# Symmetric hyperbolic systems (KOF1954)

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

In divergence form:

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

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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  $\boldsymbol{\xi}_k \mathbf{A}^k(\mathbf{x})$  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.

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

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Goals:

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

$$y\frac{\partial^2 u}{\partial x^2} + \frac{\partial^2 u}{\partial y^2} = 0;$$

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- unified treatment of equations and systems of different type.

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

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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:

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

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

(note that we use  $u = (u_d, u_{d+1})^{\top}$ , where  $u_d = -\mathbf{A}\nabla u$ , and  $u_{d+1} = u$ ).

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$$\begin{bmatrix} \mathbf{0} & \mathbf{0} \\ \mathbf{0}^{\top} & \mathbf{1} \end{bmatrix} \partial_t \begin{bmatrix} \mathbf{u}_d \\ u \end{bmatrix} + \sum_{i=1}^d \begin{bmatrix} 0 & \cdots & 0 & \cdots & 0 \\ \vdots & \ddots & 0 & \cdots & 0 \\ 0 & \cdots & 0 & \cdots & \mathbf{1} \\ \vdots & \cdots & 0 & \cdots & 0 \\ 0 & \cdots & \mathbf{1} & \cdots & \mathbf{1} \end{bmatrix} \partial_{x^i} \begin{bmatrix} \mathbf{u}_d \\ u \end{bmatrix} + \begin{bmatrix} \mathbf{A}^{-1} & \mathbf{0} \\ -(\mathbf{A}^{-1}\mathbf{b})^{\top} & c \end{bmatrix} \begin{bmatrix} \mathbf{u}_d \\ u \end{bmatrix} = \begin{bmatrix} \mathbf{0} \\ f \end{bmatrix}$$

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

Boundary conditions are enforced via matrix valued boundary field:

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Boundary condition

$$(\mathbf{A}_{\boldsymbol{\nu}} - \mathbf{M})\mathbf{u}_{|_{\Gamma}} = \mathbf{0}$$

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is sufficient for treatment of different types of usual boundary conditions.

## Assumptions on boundary matrix ${\bf M}$

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

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

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Such M is called *the admissible boundary condition*.

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

$$\begin{cases} \mathcal{L} u = f \\ (\mathbf{A}_{\boldsymbol{\nu}} - \mathbf{M}) u_{|_{\Gamma}} = 0 \end{cases}$$

# Elliptic equation - different boundary conditions

$$\mathbf{M} \qquad \mathbf{A}_{\nu} - \mathbf{M} \qquad (\mathbf{A}_{\nu} - \mathbf{M}) \begin{bmatrix} \mathbf{p} \\ u \end{bmatrix}_{|\Gamma} = \mathbf{0} \\ \begin{bmatrix} 0 & 0 & -\nu_1 \\ 0 & 0 & -\nu_2 \\ \nu_1 & \nu_2 & 0 \end{bmatrix} \qquad \begin{bmatrix} 0 & 0 & 2\nu_1 \\ 0 & 0 & 2\nu_2 \\ 0 & 0 & 0 \end{bmatrix} \qquad u_{|\Gamma} = \mathbf{0}$$

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All above matrices M satisfy (FM).

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

# Different ways to enforce boundary conditions

Instead of

$$(\mathbf{A}_{\boldsymbol{\nu}} - \mathbf{M})\mathbf{u} = \mathbf{0} \quad \text{on } \Gamma,$$

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:

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

#### Assumptions on N

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

(FX1)

(FX2)

$$\begin{split} N(\mathbf{x}) \text{ is non-negative with respect to } \mathbf{A}_{\boldsymbol{\nu}}(\mathbf{x}) \text{:} \\ (\forall \, \boldsymbol{\xi} \in N(\mathbf{x})) \quad \mathbf{A}_{\boldsymbol{\nu}}(\mathbf{x}) \boldsymbol{\xi} \cdot \boldsymbol{\xi} \geqslant 0 \, ; \end{split}$$

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

#### Assumptions on N

or

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

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

[RSP&LS1966]

Let 
$$N(\mathbf{x})$$
 and  $\tilde{N}(\mathbf{x}) := (\mathbf{A}_{\boldsymbol{\nu}}(\mathbf{x})N(\mathbf{x}))^{\perp}$  satisfy (for ae  $\mathbf{x} \in \Gamma$ )  
(FV1)  
 $(\forall \boldsymbol{\xi} \in N(\mathbf{x})) \quad \mathbf{A}_{\boldsymbol{\nu}}(\mathbf{x})\boldsymbol{\xi} \cdot \boldsymbol{\xi} \ge 0$  $(\forall \boldsymbol{\xi} \in \tilde{N}(\mathbf{x})) \quad \mathbf{A}_{\boldsymbol{\nu}}(\mathbf{x})\boldsymbol{\xi} \cdot \boldsymbol{\xi} \le 0$ 

(FV2)  $\tilde{N}(\mathbf{x}) = (\mathbf{A}_{\nu}(\mathbf{x})N(\mathbf{x}))^{\perp}$  and  $N(\mathbf{x}) = (\mathbf{A}_{\nu}(\mathbf{x})\tilde{N}(\mathbf{x}))^{\perp}$ .

Equivalence of different descriptions of boundary conditions

#### Theorem. It holds

 $(FM1)-(FM2) \iff (FX1)-(FX2) \iff (FV1)-(FV2),$ with  $N(\mathbf{x}) := \ker \left(\mathbf{A}_{\boldsymbol{\nu}}(\mathbf{x}) - \mathbf{M}(\mathbf{x})\right).$ 

Equivalence of different descriptions of boundary conditions

#### Theorem. It holds

 $\begin{array}{ll} (FM1)-(FM2) & \iff & (FX1)-(FX2) & \iff & (FV1)-(FV2) \,, \\ \mbox{with} & & \\ & & N(\mathbf{x}) := \ker \Bigl( \mathbf{A}_{\boldsymbol{\nu}}(\mathbf{x}) - \mathbf{M}(\mathbf{x}) \Bigr) \,. \end{array}$ 

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

# Classical results on well-posedness

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- uniqueness of the classical solution
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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  $A_{
  u}$ )
- 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,

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#### The Friedrichs operator

Let  $\mathcal{D}:=\mathrm{C}^\infty_c(\Omega;\mathbf{C}^r)$ ,  $L=\mathrm{L}^2(\Omega;\mathbf{C}^r)$  and  $T,\tilde{T}:\mathcal{D}\longrightarrow L$  be defined by

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

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... fits in this framework.

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

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Therefore  $T = \tilde{T}^*_{|_{W_0}}$ , and analogously  $\tilde{T} = T^*_{|_{W_0}}$ . Abusing notation:  $T, \tilde{T} \in \mathcal{L}(L; W'_0) \dots (T1)$ –(T3)

## Formulation of the problem

Lemma. The graph space

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

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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

$$\begin{split} &\ker D = W_0 \\ &\operatorname{im} D = W_0^0 := \{g \in W' : (\forall \, u \in W_0) \mid_{W'} \langle \, g, u \, \rangle_W = 0\} \end{split}$$

In particular, im D is closed in W'.

## For classical Friedrichs operator

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

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

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With the assumptions:

(FV1) 
$$\begin{array}{l} (\forall \boldsymbol{\xi} \in N(\mathbf{x})) \quad \mathbf{A}_{\boldsymbol{\nu}}(\mathbf{x})\boldsymbol{\xi} \cdot \boldsymbol{\xi} \ge 0, \\ (\forall \boldsymbol{\xi} \in \tilde{N}(\mathbf{x})) \quad \mathbf{A}_{\boldsymbol{\nu}}(\mathbf{x})\boldsymbol{\xi} \cdot \boldsymbol{\xi} \leqslant 0, \end{array}$$

(FV2) 
$$\tilde{N}(\mathbf{x}) = (\mathbf{A}_{\boldsymbol{\nu}}(\mathbf{x})N(\mathbf{x}))^{\perp}$$
 and  $N(\mathbf{x}) = (\mathbf{A}_{\boldsymbol{\nu}}(\mathbf{x})\tilde{N}(\mathbf{x}))^{\perp}$ ,

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we are lead to consider subspaces V and  $\tilde{V}$  in the functional framework:

(V1) 
$$\begin{array}{l} (\forall u \in V) & _{W'} \langle \, Du, u \, \rangle_W \geqslant 0 \,, \\ (\forall v \in \tilde{V}) & _{W'} \langle \, Dv, v \, \rangle_W \leqslant 0 \,, \end{array}$$

(V2) 
$$V = D(\tilde{V})^0, \qquad \tilde{V} = D(V)^0.$$

# Well-posedness theorem

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

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$$(\forall 1) \qquad (\forall v \in V) \quad [v \mid v] \ge 0, \\ (\forall v \in \tilde{V}) \quad [v \mid v] \le 0;$$

(V2) 
$$V = \tilde{V}^{[\perp]}, \qquad \tilde{V} = V^{[\perp]}.$$

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$$\begin{array}{ll} (\forall v \in V) & [v \mid v] \geqslant 0, \\ (\forall v \in \tilde{V}) & [v \mid v] \leqslant 0; \end{array} \end{array}$$

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$$V = \tilde{V}^{[\perp]}, \qquad \tilde{V} = V^{[\perp]}.$$

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

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

Correspondence — maximal b.c.

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

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

(FX2)

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

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(FX2) there is no non-negative subspace with respect to  ${f A}_{m 
u}({f x}),$  which contains  $N({f x}),$ 

subspace V is maximal non-negative in  $(W, [\cdot | \cdot])$ :

(X1) V is non-negative in  $(W, [\cdot | \cdot])$ :  $(\forall v \in V) [v | v] \ge 0$ ,

(X2) there is no non-negative subspace in  $(W, [\cdot | \cdot])$  containing V.

#### Correspondence — admissible b.c.

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

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

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

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abstract admissible boundary condition: there exists  $M \in \mathcal{L}(W; W')$  such that (M1)  $(\forall u \in W) \quad _{W'} \langle (M + M^*)u, u \rangle_W \ge 0$ ,

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

# Equivalence of different descriptions of b.c.

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

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Theorem. (A. Ern, J.-L. Guermond, G. Caplain) It holds  

$$(M1)-(M2) \stackrel{\Longrightarrow}{\leftarrow} (V1)-(V2) \implies (X1)-(X2),$$

with

$$V := \ker(D - M).$$

This was obtained in the real case only.

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

**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

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

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).

Graph space  $(W, \langle \cdot | \cdot \rangle_T)$  is a Hilbert space, where another (indefinite) inner product is defined:

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This inner product is dominated by the graph norm, which insures the existence of a linear operator  $G \in \mathcal{L}(W; W)$  such that

$$[ \, u \mid v \,] = \langle \, Gu \mid v \, \rangle_T \qquad \text{and} \qquad \langle \, Gu \mid v \, \rangle_T = \langle \, u \mid Gv \, \rangle_T \; .$$

Such an operator is called the Gramm operator.

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Inner product space is a *Kreĭn space* if it admits an orthogonal decomposition to its nonnegative and nonpositive parts, which are complete.

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Such an operator is called the Gramm operator.

Inner product space is a *Kreĭn space* if it admits an orthogonal decomposition to its nonnegative and nonpositive parts, which are complete.

Equivalently, a Hilbert space is a Kreĭn space if its Gramm operator is invertible.

 $(W,[\,\cdot\,|\,\cdot\,])$  is not a Kreı́n space – it is a degenerate space, because its Gramm operator  $G:=j\circ D \quad (j:W'\longrightarrow W$  is the canonical isomorphism) has large kernel:

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

# Quotient Krein 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.

# $(V1)-(V2) \qquad \Longleftrightarrow \qquad (X1)-(X2)$

**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).

 $(M1)-(M2) \implies (V1)-(V2) \quad (recall)$ 

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

#### Corollary. Under above assumptions

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

are isomorphisms.

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

**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

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

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:

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

$$P'^2 = P'$$
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im  $P' = V$  and im  $Q' = \tilde{V}$ ,  
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(b) is equivalent to closedness of  $V + \tilde{V}$ .

#### Theorem.

a)  $V, \tilde{V} \leq W$  satisfy (V), and exists a closed subspace  $W_2 \subseteq C^-$  of W,  $V + W_2 = W$ , 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  $W = W_1 + W_0 + W_2$ , 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,  $V \dotplus W_2 = W$ , and M coincide with the operator in (a).

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**Lemma.** Let  $W_2'' \leq W$  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  $W_2 + V = W$ .

**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  $W_2 + V = W$ ; b) There is a nonpositive subspace  $W_2$  of  $\hat{W}$ , such that  $W_2 + \hat{V} = \hat{W}$ .

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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
Consider

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

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

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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

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

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

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

# Heat equation

... with zero initial and Dirichlet boundary condition:

$$\begin{cases} \partial_t u - \operatorname{div}_{\mathbf{x}} (\mathbf{A} \nabla_{\mathbf{x}} u) + \mathbf{b} \cdot \nabla_{\mathbf{x}} u + cu = f \text{ in } \Omega_T \\ u = 0 \text{ on } \langle 0, T \rangle \times \Gamma \\ u(0, \cdot) = 0 \text{ on } \Omega \end{cases}$$

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

(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

$$W = \left\{ \mathbf{u} \in \mathbf{L}^{2}(\Omega_{T}; \mathbf{R}^{d+1}) : \nabla_{\mathbf{x}} u_{d+1} \in \mathbf{L}^{2}(\Omega_{T}; \mathbf{R}^{d}) \\ \& \quad \partial_{t} u_{d+1} + \operatorname{div}_{\mathbf{x}} \mathbf{u}_{d} \in \mathbf{L}^{2}(\Omega_{T}) \right\} \\ = \left\{ \mathbf{u} \in \mathbf{L}^{2}_{\operatorname{div}}(\Omega_{T}) : \nabla_{\mathbf{x}} u_{d+1} \in \mathbf{L}^{2}(\Omega_{T}; \mathbf{R}^{d}) \right\} \\ = \left\{ \mathbf{u} \in \mathbf{L}^{2}_{\operatorname{div}}(\Omega_{T}) : u_{d+1} \in \mathbf{L}^{2}(0, T; \mathbf{H}^{1}(\Omega)) \right\}.$$

## 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  $C([0, T]; L^2(\Omega))$ .

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The space

$$W(0,T) = \left\{ u \in L^2(0,T; \mathrm{H}^1(\Omega)) : \partial_t u \in \mathrm{L}^2(0,T; \mathrm{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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#### 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:

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

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

Heat equation with b = 0 and c = 0:

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

Two field theory:

developed by Ern and Guermond for elliptic problems

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

For the heat equation matrices have this form!

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

$$\begin{aligned} (\exists \,\mu_1 > 0)(\forall \,\boldsymbol{\xi} &= (\boldsymbol{\xi}_d, \boldsymbol{\xi}_{d+1}) \in \mathbf{R}^{d+1}) \\ & \left( \mathbf{C} + \mathbf{C}^\top + \sum_{k=1}^{d+1} \partial_k \mathbf{A}_k \right) \boldsymbol{\xi} \cdot \boldsymbol{\xi} \ge 2\mu_1 |\boldsymbol{\xi}_d|^2 \qquad (\text{a.e. on } \Omega) \,, \end{aligned}$$

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

$$(\exists \mu_1 > 0)(\forall \boldsymbol{\xi} = (\boldsymbol{\xi}_d, \boldsymbol{\xi}_{d+1}) \in \mathbf{R}^{d+1}) \left(\mathbf{C} + \mathbf{C}^\top + \sum_{k=1}^{d+1} \partial_k \mathbf{A}_k\right) \boldsymbol{\xi} \cdot \boldsymbol{\xi} \ge 2\mu_1 |\boldsymbol{\xi}_d|^2 \qquad (\text{a.e. on } \Omega), (\exists \mu_2 > 0)(\forall \mathbf{u} \in V \cup \widetilde{V}) \sqrt{\langle \mathcal{L}\mathbf{u} \mid \mathbf{u} \rangle_{\mathrm{L}^2(\Omega_T; \mathbf{R}^{d+1})}} + \|\mathbf{B}u_{d+1}\|_{\mathrm{L}^2(\Omega_T; \mathbf{R}^d)} \ge \mu_2 \|u_{d+1}\|_{\mathrm{L}^2(\Omega_T)},$$

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

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

$$(\exists \mu_1 > 0)(\forall \boldsymbol{\xi} = (\boldsymbol{\xi}_d, \boldsymbol{\xi}_{d+1}) \in \mathbf{R}^{d+1}) \left(\mathbf{C} + \mathbf{C}^\top + \sum_{k=1}^{d+1} \partial_k \mathbf{A}_k\right) \boldsymbol{\xi} \cdot \boldsymbol{\xi} \ge 2\mu_1 |\boldsymbol{\xi}_d|^2 \qquad (\text{a.e. on } \Omega), (\exists \mu_2 > 0)(\forall \mathbf{u} \in V \cup \widetilde{V}) \sqrt{\langle \mathcal{L}\mathbf{u} \mid \mathbf{u} \rangle_{\mathrm{L}^2(\Omega_T; \mathbf{R}^{d+1})}} + \|\mathbf{B}u_{d+1}\|_{\mathrm{L}^2(\Omega_T; \mathbf{R}^d)} \ge \mu_2 \|u_{d+1}\|_{\mathrm{L}^2(\Omega_T)},$$

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

For our system both conditions are trivially fulfilled.

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

$$\begin{split} (\exists \, \mu_1 > 0)(\forall \, \boldsymbol{\xi} &= (\boldsymbol{\xi}_d, \boldsymbol{\xi}_{d+1}) \in \mathbf{R}^{d+1}) \\ & \left( \mathbf{C} + \mathbf{C}^\top + \sum_{k=1}^{d+1} \partial_k \mathbf{A}_k \right) \boldsymbol{\xi} \cdot \boldsymbol{\xi} \geqslant 2\mu_1 |\boldsymbol{\xi}_d|^2 \qquad (\text{a.e. on } \Omega) \,, \\ (\exists \, \mu_2 > 0)(\forall \, \mathbf{u} \in V \cup \widetilde{V}) \\ & \sqrt{\langle \, \mathcal{L}\mathbf{u} \mid \mathbf{u} \, \rangle_{\mathbf{L}^2(\Omega_T; \mathbf{R}^{d+1})}} + \| \mathbf{B} u_{d+1} \|_{\mathbf{L}^2(\Omega_T; \mathbf{R}^d)} \geqslant \mu_2 \| u_{d+1} \|_{\mathbf{L}^2(\Omega_T)} \,, \end{split}$$

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

For our system both conditions are trivially fulfilled.

Therefore, we have the well-posedness result.

Some further applications ....

Dirac system Maxwell system

# 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 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

$$L \cap L^{[\perp]} = \{0\} \qquad \Longleftrightarrow \qquad \mathsf{Cl}\left(L + L^{[\perp]}\right) = W$$