# WELL-POSEDNESS OF LINEAR FIRST ORDER PORT-HAMILTONIAN SYSTEMS ON MULTIDIMENSIONAL SPATIAL DOMAINS 

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#### Abstract

We consider a port-Hamiltonian system on an open spatial domain $\Omega \subseteq \mathbb{R}^{n}$ with bounded Lipschitz boundary. We show that there is a boundary triple associated to this system. Hence, we can characterize all boundary conditions that provide unique solutions that are non-increasing in the Hamiltonian. As a by-product we develop the theory of quasi Gelfand triples. Adding "natural" boundary controls and boundary observations yields scattering/impedance passive boundary control systems. This framework will be applied to the wave equation, Maxwell's equations and Mindlin plate model. Probably, there are even more applications.


1. Introduction. The aim of this paper is to develop a port-Hamiltonian framework on multidimensional spatial domains that justifies existence and uniqueness of solutions. Those systems can be described by the following equations

$$
\begin{aligned}
\frac{\partial}{\partial t} x(t, \zeta) & =\sum_{i=1}^{n} \frac{\partial}{\partial \zeta_{i}} P_{i}(\mathcal{H}(\zeta) x(t, \zeta))+P_{0}(\mathcal{H}(\zeta) x(t, \zeta)), & & \zeta \in \Omega, t \geq 0 \\
x(0, \zeta) & =x_{0}(\zeta), & & \zeta \in \Omega
\end{aligned}
$$

where $x$ is the state, $P_{i}$ and $P_{0}$ are matrices, $\mathcal{H}$ is the Hamiltonian density, and $\Omega$ is an open subset of $\mathbb{R}^{n}$ with bounded Lipschitz boundary. We will restrict ourselves to the case, where the matrices $P_{i}$ have the block shape $\left[\begin{array}{cc}0 & L_{i} \\ L_{i}^{H} & 0\end{array}\right]$ for $i \in$ $\{1, \ldots, n\}$. We also introduce "natural" boundary controls and observations which make the system a scattering passive (energy preserving) or impedance passive (energy preserving) boundary control system. This PDE perfectly matches the description of port-Hamiltonian systems in one spatial dimension in [8], if we set $n=1$. The additional restriction $P_{1}=\left[\begin{array}{cc}0 & L_{1} \\ L_{1}^{H} & 0\end{array}\right]$ is not needed in [8], since the boundary of a line automatically satisfies certain symmetry properties. We decided to not demand an analogous symmetry from $\Omega$ in the multidimensional case, because

[^0]it did not seem very restrictive to ask for $P_{i}=\left[\begin{array}{cc}0 & L_{i} \\ L_{i}^{\mathrm{H}} & 0\end{array}\right]$ as all the examples satisfy this anyway. However, it is probably possible to drop this restriction and ask instead for certain a symmetry of the boundary.

The port-Hamiltonian formulation has proven to be a powerful tool for the modeling and control of complex multiphysics systems. An introductory overview can be found in [14]. For one-dimensional spatial domains concerns about existence and uniqueness of solutions are covered in [8].

Chapter 8 of the Ph.D. thesis [15] also regards such port-Hamiltonian systems that have multidimensional spatial domains, but the results demand very strong assumptions on the boundary operators (they have to map into $H^{1 / 2}(\partial \Omega)^{k}$ and its dual respectively), which are in case of Maxwell's equations and the Mindlin plate model not satisfied, as Example A. 4 shows for Maxwell's equations. With the following approach we will overcome these limits.

The strategy is to find a boundary triple associated to the differential operator. The multidimensional integration by parts formula already suggests possible operators for a boundary triple, but unfortunately these operators cannot be extended to the entire domain of the differential operator. Hence, we need to adapt the codomain of these boundary operators, which will lead to the construction of suitable boundary spaces for this problem. These boundary spaces behave like a Gelfand triple with the original codomain as pivot space, but lack of a chain inclusion.

Up to the author's best knowledge there is no earlier theory about this setting. So we will develop the notion of quasi Gelfand triples in section 5, which equips us with the tools to state the boundary condition in terms of the pivot space instead of the artificially constructed boundary spaces (Theorem 7.6). Section 5 can be read isolated from the rest.

One can think of using a quasi boundary triple ( $\mathcal{G}, \Gamma_{0}, \Gamma_{1}$ ) (see [1]) to overcome the extension problem of the boundary mappings, but unfortunately the condition ker $\Gamma_{0}$ is self-adjoint (or in this setting skew-adjoint) is in general not satisfied for our purpose.

The approach to the wave equation in [9] perfectly fits the framework presented in this paper. In fact, many ideas from [9] are generalized in this work. Also Maxwell's equations can be formulated as such a port-Hamiltonian system and the results in [16] can also be derived with the tools of this paper. Moreover, this theory can be applied on the model of the Mindlin plate in [2, 10]. In section 8 we give examples of how this framework can be applied to these three PDEs.

## Symbols.

| SymboL | Meaning | PAGE |
| :--- | :--- | :---: |
| $B_{r}\left(\zeta_{0}\right)$ | $\left\{\zeta \in X:\left\\|\zeta-\zeta_{0}\right\\|_{X}<r\right\}$ ball with radius $r$ and center $\zeta_{0}$ in |  |
|  | a normed space $X$ |  |
| $\mathcal{D}(\Omega)$ | set of $C^{\infty}(\Omega)$ functions with compact support | 969 |
| $\mathcal{D}^{\prime}(\Omega)$ | $($ anti)dual space of $\mathcal{D}(\Omega)$ | 969 |
| $\left.\mathcal{D}\left(\mathbb{R}^{n}\right)\right\|_{\Omega}$ | $\left\{\left.f\right\|_{\Omega}: f \in \mathcal{D}\left(\mathbb{R}^{n}\right)\right\}$ | 969 |
| $\nu$ | outward pointing normed normal vector on $\partial \Omega$ | 969 |
| $\gamma_{0}$ | $H^{1}(\Omega, X) \rightarrow L^{2}(\partial \Omega, X) ;$ extension of $\left.f \mapsto f\right\|_{\partial \Omega}$ | 969 |
| $L_{\partial}$ | $\sum_{i=1}^{n} \partial_{i} L_{i}$ a differential operator from $L^{2}(\Omega)^{m_{2}}$ to $L^{2}(\Omega)^{m_{1}}$ | 969 |


| $L_{\partial}^{\mathrm{H}}$ | $\sum_{i=1}^{n} \partial_{i} L_{i}^{\mathrm{H}}$ | 969 |
| :---: | :---: | :---: |
| $H\left(L_{\partial}, \Omega\right)$ | $\left\{f \in L^{2}(\Omega)^{m_{2}}: L_{\partial} f \in L^{2}(\Omega)^{m_{1}}\right\}$ maximal domain of $L_{\partial}$ | 969 |
| $H_{0}\left(L_{\partial}, \Omega\right)$ | The closure of $\mathcal{D}(\Omega)^{m_{2}}$ in $H\left(L_{\partial}, \Omega\right)$ | 969 |
| $H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$ | $\operatorname{ker} \bar{\pi}_{L}^{\Gamma_{0}}$ | 989 |
| $L_{\nu}$ | $\sum_{i=1}^{n} \nu_{i} L_{i}: L^{2}(\partial \Omega)^{m_{2}} \rightarrow L^{2}(\partial \Omega)^{m_{1}}$ | 969 |
| $\bar{L}_{\nu}^{\Gamma}$ | $H\left(L_{\partial}, \Omega\right) \rightarrow \mathcal{V}_{L, \Gamma}^{\prime} ;$ extension of $\mathbb{1}_{\Gamma} L_{\nu} \gamma_{0}$ on $H\left(L_{\partial}, \Omega\right)$ | 989 |
| $\bar{L}_{\nu}$ | $\bar{L}_{\nu}^{\partial \Omega}: H\left(L_{\partial}, \Omega\right) \rightarrow \mathcal{V}_{L}^{\prime} ;$ extension of $L_{\nu} \gamma_{0}$ on $H\left(L_{\partial}, \Omega\right)$ | 990 |
| $L_{\pi}^{2}(\Gamma)$ | $\overline{\operatorname{ran} \mathbb{1}_{\Gamma} L_{\nu} \gamma_{0}} \subseteq L^{2}(\Gamma)^{m_{1}}$ | 987 |
| $\pi_{L}^{\Gamma}$ | $H^{1}(\Omega)^{m_{1}} \rightarrow L_{\pi}^{2}(\Gamma)$; projection on $L_{\pi}^{2}(\Gamma)$ composed with $\gamma_{0}$ | 987 |
| $\pi_{L}$ | $\pi_{L}^{\partial \Omega}: H^{1}(\Omega)^{m_{1}} \rightarrow L^{2}(\partial \Omega)^{m_{1}}$ | 987 |
| $\bar{\pi}_{L}^{\Gamma}$ | $H\left(L_{\partial}^{\mathrm{H}}, \Omega\right) \rightarrow \mathcal{V}_{L, \Gamma}$; extension of $\pi_{L}^{\Gamma}$ on $H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$ | 988 |
| $\bar{\pi}_{L}$ | $\bar{\pi}_{L}^{\partial \Omega}: H\left(L_{\partial}^{\mathrm{H}}, \Omega\right) \rightarrow \mathcal{V}_{L}$ | 988 |
| $M_{\Gamma}$ | $\operatorname{ran} \pi_{L}^{\Gamma} \subseteq L_{\pi}^{2}(\Gamma)$ | 988 |
| $\mathcal{V}_{L, \Gamma_{1}}$ | $\left.\operatorname{ran} \bar{\pi}_{L}\right\|_{H_{\Gamma_{0}}\left(L_{\partial}^{H}, \Omega\right)}$ | 989 |
| $\mathcal{V}_{L}$ | $\mathcal{V}_{L, \partial \Omega}$ | 989 |
| $\mathcal{H}$ | Hamiltonian density | 976 |
| $\mathcal{X}_{\mathcal{H}}$ | $L^{2}(\Omega)^{m}$ equipped with $\left.\langle\mathcal{H} .,\rangle.\right\rangle_{L^{2}(\Omega)^{m}}$; the state space | 976 |
| $\left[\begin{array}{c} (\mathcal{H} x)_{L_{L}} \\ (\mathcal{H} x)_{L} \end{array}\right]$ | splitting of $\mathcal{H} x$ w.r.t. the dimensions of $L$ | 992 |
| $\mathcal{X}_{0}$ | Hilbert space; pivot space of a quasi Gelfand triple | 977 |
| $\tilde{D}_{+}$ | dense subspace of $\mathcal{X}_{0}$ with an alternative inner product | 977 |
| $D_{-}$ | $\left\{g \in \mathcal{X}_{0}: \sup _{g \in \tilde{D}_{+} \backslash\{0\}} \frac{\left\|\langle g, f\rangle_{\mathcal{X}_{0}}\right\|}{\\|f\\|_{\mathcal{X}_{+}}}<+\infty\right\}$ | 978 |

2. Boundary triple. In this section we state the most important properties of boundary triples for skew-symmetric operators for this work. More details can be found in [6, chapter 3] and [9].

A linear relation $T$ from a vector space $X$ to a vector space $Y$ is a linear subspace of $X \times Y$. Clearly, every linear operator is also a linear relation (we do not distinguish between a function and its graph). We will use the following notation

$$
\begin{aligned}
\operatorname{ker} T & :=\{x \in X:(x, 0) \in T\}, & \operatorname{ran} T:=\{y \in Y: \exists x:(x, y) \in T\} \\
\operatorname{mul} T & :=\{y \in Y:(0, y) \in T\}, & \operatorname{dom} T:=\{x \in X: \exists y:(x, y) \in T\}
\end{aligned}
$$

Thus, $T$ is single-valued, if mul $T=\{0\}$. The closure $\bar{T}$ of a linear relation $T$ is the closure in $X \times Y$. Note that every linear relation is closable. Also every operator has a closure as a linear relation, but its closure can be multi-valued. Therefore, showing mul $\bar{T}=\{0\}$ is necessary, even if $\operatorname{mul} T=\{0\}$. For an additional linear relation $S$ from $Y$ to another vector space $Z$ we define the composition $S T$ as

$$
S T:=\{(x, z) \in X \times Z: \exists y \in Y \text { such that }(x, y) \in T \text { and }(y, z) \in S\}
$$

For a linear relation $T$ from a Hilbert space $X$ to a Hilbert space $Y$ the adjoint relation is defined by

$$
T^{*}:=\left\{(u, v) \in Y \times X:\langle u, y\rangle_{Y}=\langle v, x\rangle_{X} \text { for all }(x, y) \in T\right\}
$$

and the following holds true

$$
\operatorname{ker} T^{*}=(\operatorname{ran} T)^{\perp}, \quad \operatorname{mul} T^{*}=(\operatorname{dom} T)^{\perp} \quad \text { and } \quad T^{*}=\left[\begin{array}{cc}
0 & \mathrm{I}_{Y} \\
-\mathrm{I}_{X} & 0
\end{array}\right] T^{\perp}
$$

where $\left[\begin{array}{cc}0 & \mathrm{I}_{Y} \\ -\mathrm{I}_{X} & 0\end{array}\right] T:=\{(y,-x):(x, y) \in T\}$ and $T^{\perp}$ is the orthogonal complement in $X \times Y$. A linear relation $T$ on a Hilbert space $H$ (from $H$ to $H$ ) is dissipative, if $\operatorname{Re}\langle x, y\rangle_{H} \leq 0$ for every $(x, y) \in T$ and maximal dissipative, if additionally there is no proper dissipative extension of $T$. The linear relation $T$ is (maximal) accretive, if $-T:=\{(x,-y):(x, y) \in T\}$ is (maximal) dissipative. More details can be found in [4].

Definition 2.1. Let $A_{0}$ be a densely defined, skew-symmetric, and closed operator on a Hilbert space $X$. By a boundary triple for $A_{0}^{*}$ we mean a triple ( $\mathcal{B}, B_{1}, B_{2}$ ) consisting of a Hilbert space $\mathcal{B}$, and two linear operators $B_{1}, B_{2}$ : $\operatorname{dom} A_{0}^{*} \rightarrow \mathcal{B}$ such that
(i) the mapping $\left[\begin{array}{l}B_{1} \\ B_{2}\end{array}\right]: \operatorname{dom} A_{0}^{*} \rightarrow \mathcal{B} \times \mathcal{B}, x \mapsto\left[\begin{array}{l}B_{1} x \\ B_{2} x\end{array}\right]$ is surjective, and
(ii) for $x, y \in \operatorname{dom} A_{0}^{*}$ there holds

$$
\begin{equation*}
\left\langle A_{0}^{*} x, y\right\rangle_{X}+\left\langle x, A_{0}^{*} y\right\rangle_{X}=\left\langle B_{1} x, B_{2} y\right\rangle_{\mathcal{B}}+\left\langle B_{2} x, B_{1} y\right\rangle_{\mathcal{B}} . \tag{2.1}
\end{equation*}
$$

The operator $A_{0}$ can be recovered by restricting $-A_{0}^{*}$ to $\operatorname{ker} B_{1} \cap \operatorname{ker} B_{2}$ as the next lemma will show. However, if $A_{0}^{*}$ satisfied item (i) and item (ii) but wasn't the adjoint of a skew-symmetric operator, then the next lemma would not hold as Example A. 1 demonstrates. Consequently, Proposition 2.3 would also not hold. This should highlight the importance of $A_{0}^{*}$ being the adjoint of a skew-symmetric operator in the definition of a boundary triple.

Lemma 2.2. Let $A_{0}$ be a densely defined, skew-symmetric, and closed operator on a Hilbert space $X$ and $\left(\mathcal{B}, B_{1}, B_{2}\right)$ be a boundary triple for $A_{0}^{*}$. Then $A_{0}=$ $-\left.A_{0}^{*}\right|_{\text {ker } B_{1} \cap \text { ker } B_{2}}$.

A proof can be found in [6, p. 155]. The following result is Theorem 2.2 from [9].
Proposition 2.3. Let $A_{0}$ be a skew-symmetric operator and $\left(\mathcal{B}, B_{1}, B_{2}\right)$ be a boundary triple for $A_{0}^{*}$. Consider the restriction $A$ of $A_{0}^{*}$ to a subspace $\mathcal{D}$ containing $\operatorname{ker} B_{1} \cap \operatorname{ker} B_{2}$. Define a subspace of $\mathcal{B} \times \mathcal{B}$ by $\mathcal{C}:=\left[\begin{array}{l}B_{1} \\ B_{2}\end{array}\right] \mathcal{D}$. Then the following claims are true:
(i) The domain of $A$ can be written as

$$
\operatorname{dom} A=\mathcal{D}=\left\{d \in \operatorname{dom} A_{0}^{*}:\left[\begin{array}{l}
B_{1} \\
B_{2}
\end{array}\right] d \in \mathcal{C}\right\}
$$

(ii) The operator closure of $A$ is $A_{0}^{*}$ restricted to

$$
\tilde{\mathcal{D}}:=\left\{d \in \operatorname{dom} A_{0}^{*}:\left[\begin{array}{l}
B_{1} \\
B_{2}
\end{array}\right] d \in \overline{\mathcal{C}}\right\}
$$

where $\overline{\mathcal{C}}$ is the closure in $\mathcal{B}^{2}$. Therefore, $A$ is closed if and only if $\mathcal{C}$ is closed.
(iii) The adjoint $A^{*}$ is the restriction of $-A_{0}^{*}$ to $\mathcal{D}^{\prime}$, where

$$
\mathcal{D}^{\prime}:=\{d^{\prime} \in \operatorname{dom} A_{0}^{*}:\left[\begin{array}{l}
B_{1} \\
B_{2}
\end{array}\right] d^{\prime} \in \underbrace{\left[\begin{array}{ll}
0 & \mathrm{I} \\
\mathrm{I} & 0
\end{array}\right] \mathcal{C}^{\perp}}_{=-\mathcal{C}^{*}}\}
$$

(iv) The operator $A$ is (maximal) dissipative if and only if $\mathcal{C}$ is a (maximal) dissipative relation. It also holds that $A$ is (maximal) accretive, if and only if $\mathcal{C}$ is (maximal) accretive.
3. Differential operators. Before we start analyzing port-Hamiltonian systems we will make some observation about the differential operators that will appear in the PDE. In this section we take care of all the technical details of these differential operators. Since it doesn't really make a difference whether we use the scalar field $\mathbb{R}$ or $\mathbb{C}$ we will use $\mathbb{K} \in\{\mathbb{R}, \mathbb{C}\}$ for the scalar field. The following assumption will be made for the rest of this work.

Assumption 3.1. Let $m_{1}, m_{2}, n \in \mathbb{N}, \Omega \subseteq \mathbb{R}^{n}$ be open with a bounded Lipschitz boundary, and $L=\left(L_{i}\right)_{i=1}^{n}$ such that $L_{i} \in \mathbb{K}^{m_{1} \times m_{2}}$ for all $i \in\{1, \ldots, n\}$. Corresponding to $L$ we also have $L^{\mathrm{H}}:=\left(L_{i}^{\mathrm{H}}\right)_{i=1}^{n}$, where $L_{i}^{\mathrm{H}}$ denotes the complex conjugated transposed (Hermitian transposed) matrix.

We will write $\mathcal{D}(\Omega)$ for the set of all $C^{\infty}(\Omega)$ functions with compact support in $\Omega$. Its dual space, the space of distributions, will be denoted by $\mathcal{D}^{\prime}(\Omega)$ (details on distributions can be found in [7]). Moreover, we will write $\left.\mathcal{D}\left(\mathbb{R}^{n}\right)\right|_{\Omega}$ for $\left\{\left.f\right|_{\Omega}: f \in\right.$ $\left.\mathcal{D}\left(\mathbb{R}^{n}\right)\right\}$. We will use $\partial_{i}$ as a short notation for $\frac{\partial}{\partial \zeta_{i}}$. We denote the boundary trace operator by $\gamma_{0}: H^{1}(\Omega, X) \rightarrow L^{2}(\partial \Omega, X)$ for a Banach space $X$.

Sometimes it can be confusing to pay attention to the antilinear structure of an inner product of a Hilbert space, when switching between the inner product and the dual pairing. Thus, for the sake of clarity we will always consider the antidual space instead of the dual space, which is the space of all continuous antilinear mappings from the topological vector space into its scalar field. Hence, both the inner product and the (anti)dual pairing is linear in one component and antilinear in the other. So also $\mathcal{D}^{\prime}(\Omega)$ is actually the antidual space of $\mathcal{D}(\Omega)$.

Sometimes we will write $\langle\psi, \phi\rangle_{\mathcal{D}^{\prime}, \mathcal{D}}$ instead of $\langle\psi, \phi\rangle_{\mathcal{D}^{\prime}(\Omega)^{k}, \mathcal{D}(\Omega)^{k}}$, if $\Omega$ and $k \in \mathbb{N}$ are clear or $\langle\psi, \phi\rangle_{\mathcal{D}^{\prime}(\Omega), \mathcal{D}(\Omega)}$, if only $k \in \mathbb{N}$ is clear.

Definition 3.2. Let $L$ be as in Assumption 3.1. Then we define

$$
L_{\partial}:=\sum_{i=1}^{n} \partial_{i} L_{i} \quad \text { and } \quad L_{\partial}^{\mathrm{H}}:=\left(L^{\mathrm{H}}\right)_{\partial}=\sum_{i=1}^{n} \partial_{i} L_{i}^{\mathrm{H}}
$$

as operators from $\mathcal{D}^{\prime}(\Omega)^{m_{2}}$ to $\mathcal{D}^{\prime}(\Omega)^{m_{1}}$ and from $\mathcal{D}^{\prime}(\Omega)^{m_{1}}$ to $\mathcal{D}^{\prime}(\Omega)^{m_{2}}$, respectively. Furthermore, we define the space

$$
H\left(L_{\partial}, \Omega\right):=\left\{f \in L^{2}\left(\Omega, \mathbb{K}^{m_{2}}\right): L_{\partial} f \in L^{2}\left(\Omega, \mathbb{K}^{m_{1}}\right)\right\}
$$

This space is endowed with the inner product

$$
\langle f, g\rangle_{H\left(L_{\partial}, \Omega\right)}:=\langle f, g\rangle_{L^{2}\left(\Omega, \mathbb{K}^{m_{2}}\right)}+\left\langle L_{\partial} f, L_{\partial} g\right\rangle_{L^{2}\left(\Omega, \mathbb{K}^{m_{1}}\right)}
$$

The space $H_{0}\left(L_{\partial}, \Omega\right)$ is defined as $\overline{\mathcal{D}(\Omega)^{m_{2}}}\|\cdot\|_{H\left(L_{\partial}, \Omega\right)}$. We denote the outward pointing normed normal vector on $\partial \Omega$ by $\nu$ and its $i$-th component by $\nu_{i}$. Moreover, we define

$$
L_{\nu}:=\sum_{i=1}^{n} \nu_{i} L_{i}:\left\{\begin{array}{rll}
L^{2}\left(\partial \Omega, \mathbb{K}^{m_{2}}\right) & \rightarrow & L^{2}\left(\partial \Omega, \mathbb{K}^{m_{1}}\right) \\
f & \mapsto & \sum_{i=1}^{n} \nu_{i} L_{i} f
\end{array}\right.
$$

and $L_{\nu}^{\mathrm{H}}:=\left(L^{\mathrm{H}}\right)_{\nu}$.

The operator $L_{\partial}$ can also be regarded as a linear unbounded operator from $L^{2}\left(\Omega, \mathbb{K}^{m_{2}}\right)$ to $L^{2}\left(\Omega, \mathbb{K}^{m_{1}}\right)$ with domain $H\left(L_{\partial}, \Omega\right)$. In fact this is what we will do most of the time. The same goes for $L_{\partial}^{\mathrm{H}}$ with domain $H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$. Since $\nu \in$ $L^{\infty}\left(\partial \Omega, \mathbb{R}^{n}\right)$ the mappings $L_{\nu}$ and $L_{\nu}^{\mathrm{H}}$ are well-defined and bounded.

For convenience we will write $H^{1}(\Omega)^{k}$ instead of $H^{1}\left(\Omega, \mathbb{K}^{k}\right)$ and $L^{2}(\Omega)^{k}$ instead of $L^{2}\left(\Omega, \mathbb{K}^{k}\right)$ for $k \in \mathbb{N}$.

Clearly, $\left.\mathcal{D}\left(\mathbb{R}^{n}\right)^{m_{2}}\right|_{\Omega} \subseteq H^{1}(\Omega)^{m_{2}} \subseteq H\left(L_{\partial}, \Omega\right)$ and $\left.\mathcal{D}\left(\mathbb{R}^{n}\right)^{m_{1}}\right|_{\Omega} \subseteq H^{1}(\Omega)^{m_{1}} \subseteq$ $H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$.

Example 3.3. Let us regard the following matrices

$$
L_{1}=\left[\begin{array}{lll}
1 & 0 & 0
\end{array}\right], \quad L_{2}=\left[\begin{array}{lll}
0 & 1 & 0
\end{array}\right], \quad \text { and } \quad L_{3}=\left[\begin{array}{lll}
0 & 0 & 1
\end{array}\right] .
$$

Then we obtain the corresponding differential operators

$$
L_{\partial}=\left[\begin{array}{lll}
\partial_{1} & \partial_{2} & \partial_{3}
\end{array}\right]=\operatorname{div} \quad \text { and } \quad L_{\partial}^{\mathrm{H}}=\left[\begin{array}{c}
\partial_{1} \\
\partial_{2} \\
\partial_{3}
\end{array}\right]=\operatorname{grad}
$$

The corresponding operator $L_{\nu}$ that acts on $L^{2}(\partial \Omega)$ can be written as an inner product

$$
L_{\nu} f=\left[\begin{array}{lll}
\nu_{1} & \nu_{2} & \nu_{3}
\end{array}\right]\left[\begin{array}{l}
f_{1} \\
f_{2} \\
f_{3}
\end{array}\right]=\nu \cdot f
$$

Clearly the previous example can be extended to any finite dimension.
Example 3.4. The following matrices will construct the rotation operator.

$$
L_{1}=\left[\begin{array}{ccc}
0 & 0 & 0 \\
0 & 0 & -1 \\
0 & 1 & 0
\end{array}\right], \quad L_{2}=\left[\begin{array}{ccc}
0 & 0 & 1 \\
0 & 0 & 0 \\
-1 & 0 & 0
\end{array}\right], \quad \text { and } \quad L_{3}=\left[\begin{array}{ccc}
0 & -1 & 0 \\
1 & 0 & 0 \\
0 & 0 & 0
\end{array}\right]
$$

In this example we have $L_{i}^{\mathrm{H}}=-L_{i}$. Furthermore, the corresponding differential operator is

$$
L_{\partial}=\left[\begin{array}{ccc}
0 & -\partial_{3} & \partial_{2} \\
\partial_{3} & 0 & -\partial_{1} \\
-\partial_{2} & \partial_{1} & 0
\end{array}\right]=\operatorname{rot}=-L_{\partial}^{\mathrm{H}} .
$$

The corresponding operator $L_{\nu}$ that acts on $L^{2}(\partial \Omega)$ can be written as a cross product

$$
L_{\nu} f=\left[\begin{array}{ccc}
0 & -\nu_{3} & \nu_{2} \\
\nu_{3} & 0 & -\nu_{1} \\
-\nu_{2} & \nu_{1} & 0
\end{array}\right]\left[\begin{array}{l}
f_{1} \\
f_{2} \\
f_{3}
\end{array}\right]=\nu \times f
$$

Lemma 3.5. The operator $L_{\partial}$ with $\operatorname{dom} L_{\partial}=H\left(L_{\partial}, \Omega\right)$ is a closed operator from $L^{2}(\Omega)^{m_{2}}$ to $L^{2}(\Omega)^{m_{1}}$ and $H\left(L_{\partial}, \Omega\right)$ endowed with the inner product $\langle., .\rangle_{H\left(L_{\partial}, \Omega\right)}$ is a Hilbert space.

Note that for $f \in \mathcal{D}^{\prime}(\Omega)^{m_{2}}$ and $\phi \in \mathcal{D}(\Omega)^{m_{1}}$ we have

$$
\begin{aligned}
\left\langle L_{\partial} f, \phi\right\rangle_{\mathcal{D}^{\prime}(\Omega)^{m_{1}}, \mathcal{D}(\Omega)^{m_{1}}} & =\sum_{i=1}^{n}\left\langle\partial_{i} L_{i} f, \phi\right\rangle_{\mathcal{D}^{\prime}(\Omega)^{m_{1}}, \mathcal{D}(\Omega)^{m_{1}}} \\
& =\sum_{i=1}^{n}\left\langle f,-\partial_{i} L_{i}^{\mathrm{H}} \phi\right\rangle_{\mathcal{D}^{\prime}(\Omega)^{m_{2}}, \mathcal{D}(\Omega)^{m_{2}}}=\left\langle f,-L_{\partial}^{\mathrm{H}} \phi\right\rangle_{\mathcal{D}^{\prime}(\Omega)^{m_{2}}, \mathcal{D}(\Omega)^{m_{2}}} .
\end{aligned}
$$

Proof. Let $\left(\left(f_{k}, L_{\partial} f_{k}\right)\right)_{k \in \mathbb{N}}$ be a sequence in $L_{\partial}$ that converges to a point $(f, g) \in$ $L^{2}(\Omega)^{m_{2}} \times L^{2}(\Omega)^{m_{1}}$. For an arbitrary $\phi \in \mathcal{D}(\Omega)^{m_{1}}$ we have

$$
\begin{aligned}
\langle g, \phi\rangle_{\mathcal{D}^{\prime}(\Omega)^{m_{1}}, \mathcal{D}(\Omega)^{m_{1}}} & =\lim _{k \rightarrow \infty}\left\langle L_{\partial} f_{k}, \phi\right\rangle_{\mathcal{D}^{\prime}(\Omega)^{m_{1}}, \mathcal{D}(\Omega)^{m_{1}}} \\
& =\lim _{k \rightarrow \infty}\left\langle f_{k},-L_{\partial}^{\mathrm{H}} \phi\right\rangle_{\mathcal{D}^{\prime}(\Omega)^{m_{1}}, \mathcal{D}(\Omega)^{m_{1}}} \\
& =\left\langle f,-L_{\partial}^{\mathrm{H}} \phi\right\rangle_{\mathcal{D}^{\prime}(\Omega)^{m_{2}}, \mathcal{D}(\Omega)^{m_{2}}} \\
& =\left\langle L_{\partial} f, \phi\right\rangle_{\mathcal{D}^{\prime}(\Omega)^{m_{1}}, \mathcal{D}(\Omega)^{m_{1}}}
\end{aligned}
$$

which implies $g=L_{\partial} f$. Since $g$ is also in $L^{2}(\Omega)^{m_{1}}$, we conclude that $L_{\partial}$ is closed. Hence, dom $L_{\partial}=H\left(L_{\partial}, \Omega\right)$ endowed with the graph norm of $L_{\partial}$, which is induced by $\langle., .\rangle_{H\left(L_{\partial}, \Omega\right)}$, is a Hilbert space.

Lemma 3.6. The adjoint of $L_{\partial}$ with dom $L_{\partial}=H\left(L_{\partial}, \Omega\right)$ (as an unbounded operator/linear relation from $L^{2}(\Omega)^{m_{2}}$ to $L^{2}(\Omega)^{m_{1}}$ ) is given by $L_{\partial}^{*} g=-L_{\partial}^{\mathrm{H}} g$ for $g \in \operatorname{dom} L_{\partial}^{*} \subseteq H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$, i.e. $L_{\partial}^{*} \subseteq-L_{\partial}^{\mathrm{H}}$.
Proof. For an arbitrary $g \in \operatorname{dom} L_{\partial}^{*}$ and an arbitrary $\phi \in \mathcal{D}(\Omega)^{m_{2}}$ we have

$$
\left\langle L_{\partial}^{*} g, \phi\right\rangle_{\mathcal{D}^{\prime}, \mathcal{D}}=\left\langle L_{\partial}^{*} g, \phi\right\rangle_{L^{2}}=\left\langle g, L_{\partial} \phi\right\rangle_{L^{2}}=\left\langle g, L_{\partial} \phi\right\rangle_{\mathcal{D}^{\prime}, \mathcal{D}}=\left\langle-L_{\partial}^{\mathrm{H}} g, \phi\right\rangle_{\mathcal{D}^{\prime}, \mathcal{D}}
$$

Therefore, $L_{\partial}^{*} g=-L_{\partial}^{\mathrm{H}} g$ and $L_{\partial}^{*} g \in L^{2}(\Omega)^{m_{2}}$ implies $L_{\partial}^{\mathrm{H}} g \in L^{2}(\Omega)^{m_{2}}$. Consequently, $\operatorname{dom} L_{\partial}^{*} \subseteq H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$.
Remark 3.7. If $L$ contains only Hermitian matrices $\left(L_{i}^{\mathrm{H}}=L_{i}\right)$, then $L_{\partial}^{\mathrm{H}}=L_{\partial}$ and $L_{\partial}^{*}$ is skew-symmetric by the previous lemma.

The next result is an integration by parts version for $L_{\partial}$. This will be helpful to construct a boundary triple for the differential operator in the port-Hamiltonian PDE.

Lemma 3.8. Let $f \in H^{1}(\Omega)^{m_{2}}$ and $g \in H^{1}(\Omega)^{m_{1}}$. Then we have

$$
\begin{align*}
\left\langle L_{\partial} f, g\right\rangle_{L^{2}(\Omega)^{m_{1}}}+\left\langle f, L_{\partial}^{\mathrm{H}} g\right\rangle_{L^{2}(\Omega)^{m_{2}}} & =\left\langle L_{\nu} \gamma_{0} f, \gamma_{0} g\right\rangle_{L^{2}(\partial \Omega)^{m_{1}}} \\
& =\left\langle\gamma_{0} f, L_{\nu}^{\mathrm{H}} \gamma_{0} g\right\rangle_{L^{2}(\partial \Omega)^{m_{2}}} \tag{3.1}
\end{align*}
$$

Proof. Let $\left.f \in \mathcal{D}\left(\mathbb{R}^{n}\right)^{m_{2}}\right|_{\Omega}$ and $\left.g \in \mathcal{D}\left(\mathbb{R}^{n}\right)^{m_{1}}\right|_{\Omega}$. By the definition of $L_{\partial}$ and $L_{\partial}^{\mathrm{H}}$, and the linearity of the scalar product we can write the left-hand-side of (3.1) as

$$
\int_{\Omega} \sum_{i=1}^{n}\left\langle\partial_{i} L_{i} f, g\right\rangle+\left\langle f, \partial_{i} L_{i}^{\mathrm{H}} g\right\rangle \mathrm{d} \lambda=\int_{\Omega} \sum_{i=1}^{n}\left\langle\partial_{i} L_{i} f, g\right\rangle+\left\langle L_{i} f, \partial_{i} g\right\rangle \mathrm{d} \lambda
$$

where $\lambda$ denotes the Lebesgue measure. By the product rule for derivatives and Gauß's theorem (divergence theorem) (see [7, eq. (3.1.6)] or [13, Remark 13.7.2]) this is equal to

$$
\int_{\Omega} \sum_{i=1}^{n} \partial_{i}\left\langle L_{i} f, g\right\rangle \mathrm{d} \lambda=\int_{\partial \Omega} \sum_{i=1}^{n} \nu_{i} \gamma_{0}\left\langle L_{i} f, g\right\rangle \mathrm{d} \mu=\int_{\partial \Omega}\left\langle L_{\nu} \gamma_{0} f, \gamma_{0} g\right\rangle \mathrm{d} \mu
$$

where $\nu$ denotes the outward pointing normed normal vector on $\partial \Omega$ and $\mu$ denotes the surface measure of $\partial \Omega$. By density we can extend this equality for $f \in H^{1}(\Omega)^{m_{2}}$ and $g \in H^{1}(\Omega)^{m_{1}}$.

Corollary 3.9. Let $f \in H^{1}(\Omega)^{m_{2}}$ and $g \in H^{1}(\Omega)^{m_{1}}$. Then we have

$$
\left|\left\langle L_{\nu} \gamma_{0} f, \gamma_{0} g\right\rangle_{L^{2}(\partial \Omega)^{m_{1}}}\right| \leq\|f\|_{H\left(L_{z}, \Omega\right)}\|g\|_{H\left(L_{\partial}^{H}, \Omega\right)} .
$$

Proof. Lemma 3.8, the triangular inequality and Cauchy Schwarz's inequality yield

$$
\begin{aligned}
\left|\left\langle L_{\nu} \gamma_{0} f, \gamma_{0} g\right\rangle_{L^{2}(\partial \Omega)^{m_{1}}}\right| & \leq\left|\left\langle L_{\partial} f, g\right\rangle_{L^{2}(\Omega)^{m_{1}}}\right|+\left|\left\langle f, L_{\partial}^{\mathrm{H}} g\right\rangle_{L^{2}(\Omega)^{m_{2}}}\right| \\
& \leq\left\|L_{\partial} f\right\|_{L^{2}(\Omega)^{m_{1}}}\|g\|_{L^{2}(\Omega)^{m_{1}}}+\|f\|_{L^{2}(\Omega)^{m_{2}}}\left\|L_{\partial}^{\mathrm{H}} g\right\|_{L^{2}(\Omega)^{m_{2}}} \\
& \leq \sqrt{\left\|L_{\partial} f\right\|_{L^{2}}^{2}+\|f\|_{L^{2}}^{2} \sqrt{\|g\|_{L^{2}}^{2}+\left\|L_{\partial}^{\mathrm{H}} g\right\|_{L^{2}}^{2}}} \\
& =\|f\|_{H\left(L_{\partial}, \Omega\right)}\|g\|_{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)} .
\end{aligned}
$$

Note that $\Omega=\mathbb{R}^{n}$ satisfies the assumptions in Assumption 3.1. Hence, all the previous results hold true for $\Omega=\mathbb{R}^{n}$.

Our next goal is to show that $\left.\mathcal{D}\left(\mathbb{R}^{n}\right)^{m_{2}}\right|_{\Omega}$ is dense in $H\left(L_{\partial}, \Omega\right)$; see Theorem 3.18. In order to archive this we will present some regularization and continuity results. In particular the density is needed to extend the integration by parts formula (Lemma 3.8) for $f \in H\left(L_{\partial}, \Omega\right)$ and $g \in H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$.
Lemma 3.10. The mapping $\iota: H\left(L_{\partial}, \mathbb{R}^{n}\right) \rightarrow H\left(L_{\partial}, \Omega\right),\left.f \mapsto f\right|_{\Omega}$ is well-defined and continuous for any open set $\Omega \subseteq \mathbb{R}^{n}$. In particular, $L_{\partial}\left(\left.f\right|_{\Omega}\right)=\left.\left(L_{\partial} f\right)\right|_{\Omega}$. Moreover, if $f_{k} \rightarrow f$ in $H\left(L_{\partial}, \mathbb{R}^{n}\right)$, then $f_{k} \rightarrow f$ in $H\left(L_{\partial}, \Omega\right)$.

Hence, we can always regard an $f \in H\left(L_{\partial}, \mathbb{R}^{n}\right)$ as an element of $H\left(L_{\partial}, \Omega\right)$, especially when supp $f \subseteq \bar{\Omega}$ - then it is also possible to recover $f$ from $\left.f\right|_{\Omega}$.

Proof. If $f \in H\left(L_{\partial}, \mathbb{R}^{n}\right)$, then $f \in L^{2}\left(\mathbb{R}^{n}\right)^{m_{2}}$ and $L_{\partial} f \in L^{2}\left(\mathbb{R}^{n}\right)^{m_{1}}$. Hence, it is easy to see that $\left\|\left.f\right|_{\Omega}\right\|_{L^{2}(\Omega)} \leq\|f\|_{L^{2}\left(\mathbb{R}^{n}\right)}$ and $\left\|\left.\left(L_{\partial} f\right)\right|_{\Omega}\right\|_{L^{2}(\Omega)} \leq\left\|L_{\partial} f\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}$. Note that $\mathcal{D}(\Omega) \subseteq \mathcal{D}\left(\mathbb{R}^{n}\right)$, and that for $g \in L^{2}\left(\mathbb{R}^{n}\right)$ and $\phi \in \mathcal{D}(\Omega)$

$$
\langle g, \phi\rangle_{\mathcal{D}^{\prime}\left(\mathbb{R}^{n}\right), \mathcal{D}\left(\mathbb{R}^{n}\right)}=\langle g, \phi\rangle_{L^{2}\left(\mathbb{R}^{n}\right)}=\left\langle\left. g\right|_{\Omega}, \phi\right\rangle_{L^{2}(\Omega)}=\left\langle\left. g\right|_{\Omega}, \phi\right\rangle_{\mathcal{D}^{\prime}(\Omega), \mathcal{D}(\Omega)}
$$

Hence, for $f \in H\left(L_{\partial}, \mathbb{R}^{n}\right)$ and $\phi \in \mathcal{D}(\Omega)$ we have

$$
\begin{aligned}
\left\langle L_{\partial}\left(\left.f\right|_{\Omega}\right), \phi\right\rangle_{\mathcal{D}^{\prime}(\Omega), \mathcal{D}(\Omega)} & =\left\langle\left. f\right|_{\Omega},-L_{\partial}^{\mathrm{H}} \phi\right\rangle_{\mathcal{D}^{\prime}(\Omega), \mathcal{D}(\Omega)}=\left\langle f,-L_{\partial}^{\mathrm{H}} \phi\right\rangle_{\mathcal{D}^{\prime}\left(\mathbb{R}^{n}\right), \mathcal{D}\left(\mathbb{R}^{n}\right)} \\
& =\left\langle L_{\partial} f, \phi\right\rangle_{\mathcal{D}^{\prime}\left(\mathbb{R}^{n}\right), \mathcal{D}\left(\mathbb{R}^{n}\right)} \\
& =\left\langle\left.\left(L_{\partial} f\right)\right|_{\Omega}, \phi\right\rangle_{\mathcal{D}^{\prime}(\Omega), \mathcal{D}(\Omega)}
\end{aligned}
$$

which implies $L_{\partial}\left(\left.f\right|_{\Omega}\right)=\left.\left(L_{\partial} f\right)\right|_{\Omega}$ in $\mathcal{D}^{\prime}(\Omega)$. Since the latter is in $L^{2}(\Omega)$, we conclude $\left.f\right|_{\Omega} \in H\left(L_{\partial}, \Omega\right)$. Consequently, $\iota$ is well-defined and $\|\iota f\|_{H\left(L_{\partial}, \Omega\right)} \leq\|f\|_{H\left(L_{\partial}, \mathbb{R}^{n}\right)}$ by the norm estimates from the beginning. Since $\iota$ is linear this implies the continuity of $\iota$ and in turn the last assertion of the lemma.

Lemma 3.11. Let $D_{\eta}: L^{2}\left(\mathbb{R}^{n}\right)^{k} \rightarrow L^{2}\left(\mathbb{R}^{n}\right)^{k}$ be the mapping defined by $\left(D_{\eta} f\right)(\zeta):=$ $f(\eta \zeta)$, where $\eta \in(0,+\infty)$ and $k \in \mathbb{N}$. Then $D_{\eta}$ converges in the strong operator topology to I for $\eta \rightarrow 1$.

Proof. For $\phi \in \mathcal{D}\left(\mathbb{R}^{n}\right)^{k}$ we will show that $\eta \mapsto D_{\eta} \phi$ from $(0,+\infty)$ to $L^{2}\left(\mathbb{R}^{n}\right)^{k}$ is continuous:

$$
\begin{aligned}
\left\|D_{\eta_{1}} \phi-D_{\eta_{2}} \phi\right\|_{L^{2}}^{2} & =\int_{\mathbb{R}^{n}}\left\|\phi\left(\eta_{1} \zeta\right)-\phi\left(\eta_{2} \zeta\right)\right\|_{\mathbb{K}^{k}}^{2} \mathrm{~d} \lambda(\zeta) \\
& =\frac{1}{\eta_{2}^{2 n}} \int_{\mathbb{R}^{n}}\left\|\phi\left(\frac{\eta_{1}}{\eta_{2}} \zeta\right)-\phi(\zeta)\right\|_{\mathbb{K}^{k}}^{2} \mathrm{~d} \lambda(\zeta) \rightarrow 0 \quad \text { for } \quad \eta_{2} \rightarrow \eta_{1}
\end{aligned}
$$

by Lebesgue's dominated convergence theorem, where $\lambda$ denotes the Lebesgue measure. For $f \in L^{2}\left(\mathbb{R}^{n}\right)^{k}$ there exists a sequence $\left(\phi_{m}\right)_{m \in \mathbb{N}}$ of $\mathcal{D}\left(\mathbb{R}^{n}\right)^{k}$ functions that converges to $f$ (w.r.t. $\|\cdot\|_{L^{2}}$ ). Hence,

$$
\left\|D_{\eta} \phi_{m}-D_{\eta} f\right\|_{L^{2}}=\frac{1}{\eta^{n}}\left\|\phi_{m}-f\right\|_{L^{2}}
$$

and $D_{\eta} \phi_{m}$ converges uniformly in $\eta \in(\epsilon,+\infty), \epsilon>0$ to $D_{\eta} f$ for $m \rightarrow \infty$. Consequently $\eta \mapsto D_{\eta} f$ is also continuous from $(\epsilon,+\infty)$ to $L^{2}(\Omega)^{k}$ and in particular $D_{\eta} f \rightarrow f$ for $\eta \rightarrow 1$.

Definition 3.12. A set $O \subseteq \mathbb{R}^{n}$ is strongly star-shaped with respect to $\zeta_{0}$, if for every $\zeta \in \bar{O}$ the half-open line segment $\left\{\theta\left(\zeta-\zeta_{0}\right)+\zeta_{0}: \theta \in[0,1)\right\}$ is contained in $O$. We call $O$ strongly star-shaped, if there is a $\zeta_{0}$ such that $O$ is strongly star-shaped with respect to $\zeta_{0}$.

Note that this is equivalent to

$$
\theta\left(\bar{O}-\zeta_{0}\right)+\zeta_{0} \subseteq O \quad \text { for all } \quad \theta \in[0,1)
$$

Lemma 3.13. Let $f \in H\left(L_{\partial}, \mathbb{R}^{n}\right)$ and $\zeta_{0} \in \mathbb{R}^{n}$. Furthermore, let $f_{\theta}(\zeta):=f\left(\frac{1}{\theta}(\zeta-\right.$ $\left.\left.\zeta_{0}\right)+\zeta_{0}\right)$ for $\theta \in(0,1)$ and a.e. $\zeta \in \mathbb{R}^{n}$. Then $f_{\theta} \in H\left(L_{\partial}, \mathbb{R}^{n}\right)$ and $f_{\theta} \rightarrow f$ in $H\left(L_{\partial}, \mathbb{R}^{n}\right)$ as $\theta \rightarrow 1$. If there exists a strongly star-shaped set $O$ with respect to the previous $\zeta_{0}$ such that supp $f \subseteq \bar{O}$, then $\operatorname{supp} f_{\theta} \subseteq O$ for $\theta \in(0,1)$.

Proof. Let $f \in H\left(L_{\partial}, \mathbb{R}^{n}\right)$ and $\alpha(\zeta):=\frac{1}{\theta}\left(\zeta-\zeta_{0}\right)+\zeta_{0}$. Then it is easy to see that $f_{\theta}=f \circ \alpha$ and $f_{\theta} \in L^{2}\left(\mathbb{R}^{n}\right)^{m_{2}}$. By change of variables we have

$$
\begin{aligned}
\left\langle L_{\partial}(f \circ \alpha), \phi\right\rangle_{\mathcal{D}^{\prime}\left(\mathbb{R}^{n}\right), \mathcal{D}\left(\mathbb{R}^{n}\right)} & =\left\langle f,-\left(L_{\partial}^{\mathrm{H}} \phi\right) \circ \alpha^{-1} \theta^{n}\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)} \\
& =\left\langle f,-\sum_{i=1}^{n} L_{i}^{\mathrm{H}} \partial_{i}\left(\phi \circ \alpha^{-1} \frac{1}{\theta}\right) \theta^{n}\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)} \\
& =\left\langle f,-L_{\partial}^{\mathrm{H}}\left(\frac{1}{\theta} \phi \circ \alpha^{-1}\right) \theta^{n}\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)}=\left\langle\frac{1}{\theta}\left(L_{\partial} f\right) \circ \alpha, \phi\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)} \\
& =\left\langle\frac{1}{\theta}\left(L_{\partial} f\right) \circ \alpha, \phi\right\rangle_{\mathcal{D}^{\prime}\left(\mathbb{R}^{n}\right), \mathcal{D}\left(\mathbb{R}^{n}\right)}
\end{aligned}
$$

Therefore, $L_{\partial} f_{\theta}=\frac{1}{\theta}\left(L_{\partial} f\right)_{\theta}$ and $f_{\theta} \in H\left(L_{\partial}, \mathbb{R}^{n}\right)$. We can also write $f_{\theta}$ as $T_{\zeta_{0}} D_{\frac{1}{\theta}} T_{-\zeta_{0}} f$, where $T_{\xi}: L^{2}\left(\mathbb{R}^{n}\right)^{m_{2}} \rightarrow L^{2}\left(\mathbb{R}^{n}\right)^{m_{2}}$ is the translation mapping $f \mapsto$ $f(.+\xi)$ and $D_{\eta}: L^{2}\left(\mathbb{R}^{n}\right)^{m_{2}} \rightarrow L^{2}\left(\mathbb{R}^{n}\right)^{m_{2}}$ is the mapping from Lemma 3.11. Since $T_{\xi}$ is bounded and $D_{\eta}$ converges strongly to I as $\eta \rightarrow 1$, we conclude $f_{\theta} \rightarrow f$ in $L^{2}\left(\mathbb{R}^{n}\right)^{m_{2}}$ as $\theta \rightarrow 1$ and $L_{\partial} f_{\theta}=\frac{1}{\theta}\left(L_{\partial} f\right)_{\theta} \rightarrow L_{\partial} f$ in $L^{2}\left(\mathbb{R}^{n}\right)^{m_{1}}$ as $\theta \rightarrow 1$. Hence, $f_{\theta} \rightarrow f$ in $H\left(L_{\partial}, \mathbb{R}^{n}\right)$.

Let $O$ be strongly star-shaped with respect to $\zeta_{0}$ and $\operatorname{supp} f \subseteq \bar{O}$. Then for $\theta \in(0,1)$

$$
\operatorname{supp} f_{\theta}=\theta\left(\operatorname{supp} f-\zeta_{0}\right)+\zeta_{0} \subseteq \theta\left(\bar{O}-\zeta_{0}\right)+\zeta_{0} \subseteq O
$$

Remark 3.14. If $f \in H\left(L_{\partial}, \Omega\right)$ and $\left.\psi \in \mathcal{D}\left(\mathbb{R}^{n}\right)\right|_{\Omega}$, then by the product rule for distributional derivatives also $\psi f \in H\left(L_{\partial}, \Omega\right)$ and $L_{\partial}(\psi f)=\psi L_{\partial} f+\sum_{i=1}^{n}\left(\partial_{i} \psi\right) L_{i} f$ (see [7, equation (3.1.1)']).

Lemma 3.15. For every $f \in H\left(L_{\partial}, \mathbb{R}^{n}\right)$ exists a sequence $\left(f_{k}\right)_{k \in \mathbb{N}}$ in $H\left(L_{\partial}, \mathbb{R}^{n}\right)$ with compact support $\operatorname{supp} f_{k} \subseteq \operatorname{supp} f$ that converges to $f$ in $H\left(L_{\partial}, \mathbb{R}^{n}\right)$.

Proof. Let $\psi \in C^{\infty}\left(\mathbb{R}^{n}, \mathbb{R}\right)$ be such that

$$
\psi(\zeta) \in \begin{cases}\{1\}, & \text { if }\|\zeta\| \leq 1 \\ {[0,1],} & \text { if } 1<\|\zeta\|<2 \\ \{0\}, & \text { if }\|\zeta\| \geq 2\end{cases}
$$

Then $f_{k}:=\psi(\dot{\bar{k}}) f \in L^{2}\left(\mathbb{R}^{n}\right)^{m_{2}}$ and $f_{k} \rightarrow f$ in $L^{2}$. By the previous remark we have $L_{\partial} f_{k}=\psi(\dot{\bar{k}}) L_{\partial} f+\frac{1}{k} \sum_{i=1}^{n}\left(\partial_{i} \psi\right)\left(\frac{1}{k}.\right) L_{i} f$ and therefore $f_{k} \in H\left(L_{\partial}, \mathbb{R}^{n}\right)$. Since $\left\|\partial_{i} \psi\right\|_{\infty}<\infty$ and $\left\|L_{i} f\right\|_{L^{2}} \leq\left\|L_{i}\right\|\|f\|_{L^{2}}<\infty$, we have $L_{\partial} f_{k} \rightarrow L_{\partial} f$ as $\psi(\dot{\bar{k}}) L_{\partial} f \rightarrow$ $L_{\partial} f$ in $L^{2}\left(\mathbb{R}^{n}\right)^{m_{2}}$ and consequently $f_{k} \rightarrow f$ in $H\left(L_{\partial}, \mathbb{R}^{n}\right)$.

The next result is essentially [3, Proposition 2.5.4, page 69], except that we allow $\Omega$ to be unbounded.

Lemma 3.16. For $\Omega \subseteq \mathbb{R}^{n}$ (open with bounded Lipschitz boundary) there exists an open covering $\left(O_{i}\right)_{i=0}^{k}$ of $\bar{\Omega}$ such that $O_{i} \cap \Omega$ is bounded and strongly star-shaped for $i \in\{1, \ldots, k\}$ and $\overline{\overline{O_{0}}} \subseteq \Omega$

Proof. Since $\Omega$ has a bounded Lipschitz boundary, there is an open ball $B_{r}(0)$ such that $\partial \Omega \subseteq B_{r}(0)$. Hence, $B_{r}(0) \cap \Omega$ is bounded and open with bounded Lipschitz boundary and we can apply [3, Proposition 2.5.4, page 69]. This gives an open covering $\left(O_{i}\right)_{i=1}^{k}$ of $B_{r}(0) \cap \Omega$ and in particular of $\partial \Omega$ such that $O_{i} \cap \Omega$ is strongly star-shaped. We define $O_{0}$ as $B_{\epsilon}\left(\Omega \backslash \bigcup_{i=1}^{k} O_{i}\right)$, where $\epsilon>0$ is small enough such that $\overline{O_{0}} \subseteq \Omega$.

The next lemma is similar to [5, Lemma 1, page 206], which proves the result for $L_{\partial}=$ rot. The main idea of the proof can be adopted.

Lemma 3.17. If $f \in H\left(L_{\partial}, \Omega\right)$ is such that

$$
\begin{equation*}
\left\langle L_{\partial} f, \phi\right\rangle_{L^{2}(\Omega)}+\left\langle f, L_{\partial}^{\mathrm{H}} \phi\right\rangle_{L^{2}(\Omega)}=0 \quad \text { for all } \quad \phi \in \mathcal{D}\left(\mathbb{R}^{n}\right)^{m_{1}} \tag{3.2}
\end{equation*}
$$

then $f \in H_{0}\left(L_{\partial}, \Omega\right)$.
Recall the definition of a positive mollifier: Let $\rho \in \mathcal{D}\left(\mathbb{R}^{n}\right)$. Then we define $\rho_{\epsilon}$ by $\rho_{\epsilon}(\zeta)=\epsilon^{-n} \rho\left(\frac{\zeta}{\epsilon}\right)$. We say that $\rho_{\epsilon}$ is a positive mollifier, if $\rho(\zeta) \geq 0, \int_{\mathbb{R}^{n}} \rho(\zeta) \mathrm{d} \zeta=1$ and $\lim _{\epsilon \rightarrow 0} \rho_{\epsilon}=\delta_{0}$ in the sense of distributions, where $\delta_{0}$ is the Dirac delta function $\left(\left\langle\delta_{0}, \phi\right\rangle_{\mathcal{D}^{\prime}, \mathcal{D}}=\phi(0)\right)$.

In particular, for every $f \in L^{2}\left(\mathbb{R}^{n}\right)$ holds

$$
\rho_{\epsilon} * f:=\int_{\mathbb{R}^{n}} \rho_{\epsilon}(\zeta) f(.-\zeta) \mathrm{d} \zeta \xrightarrow{\epsilon \rightarrow 0} f \quad \text { in } \quad L^{2}\left(\mathbb{R}^{n}\right) .
$$

Proof. Let $f \in H\left(L_{\partial}, \Omega\right)$ satisfy (3.2). Then we have to find a sequence $\left(f_{n}\right)_{n \in \mathbb{N}}$ in $\mathcal{D}(\Omega)^{m_{2}}$ that converges to $f$ with respect to $\|\cdot\|_{H\left(L_{\partial}, \Omega\right)}$.

We define $\tilde{f}$ and $\widetilde{L_{\partial} f}$ as the extension of $f$ and $L_{\partial} f$ respectively on $\mathbb{R}^{n}$ such that these functions are 0 outside of $\Omega$. By

$$
\begin{aligned}
\left\langle\widetilde{L_{\partial} f}, \phi\right\rangle_{\mathcal{D}^{\prime}\left(\mathbb{R}^{n}\right), \mathcal{D}\left(\mathbb{R}^{n}\right)} & =\left\langle\widetilde{L_{\partial} f}, \phi\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)}=\left\langle L_{\partial} f, \phi\right\rangle_{L^{2}(\Omega)} \stackrel{(3.2)}{=}\left\langle f,-L_{\partial}^{\mathrm{H}} \phi\right\rangle_{L^{2}(\Omega)} \\
& =\left\langle\tilde{f},-L_{\partial}^{\mathrm{H}} \phi\right\rangle_{L^{2}\left(\mathbb{R}^{n}\right)}=\left\langle\tilde{f},-L_{\partial}^{\mathrm{H}} \phi\right\rangle_{\mathcal{D}^{\prime}\left(\mathbb{R}^{n}\right), \mathcal{D}\left(\mathbb{R}^{n}\right)}
\end{aligned}
$$

for $\phi \in \mathcal{D}\left(\mathbb{R}^{n}\right)^{m_{1}}$, we see that $\widetilde{L_{\partial} f}=L_{\partial} \tilde{f}$ and $\tilde{f} \in H\left(L_{\partial}, \mathbb{R}^{n}\right)$ with $\operatorname{supp} \tilde{f} \subseteq \bar{\Omega}$.

By Lemma 3.16 there is a finite open covering $\left(O_{i}\right)_{i=0}^{k}$ of $\bar{\Omega}$ such that $O_{i} \cap \bar{\Omega}$ is strongly star-shaped for $i \in\{1, \ldots, k\}$ and $\overline{O_{0}} \subseteq \Omega$. We employ a partition of unity and obtain $\left(\alpha_{i}\right)_{i=0}^{k}$, subordinate to this covering, that is
$\alpha_{i} \in C^{\infty}\left(\mathbb{R}^{n}\right), \quad \operatorname{supp} \alpha_{i} \subseteq O_{i}, \quad \alpha_{i}(\zeta) \in[0,1], \quad$ and $\quad \sum_{i=0}^{k} \alpha_{i}(\zeta)=1 \quad$ for $\quad \zeta \in \Omega$.
Hence, $\tilde{f}=\sum_{i=0}^{k} \alpha_{i} \tilde{f}$ and we define $f_{i}:=\alpha_{i} \tilde{f}$. By construction $f_{i} \in H\left(L_{\partial}, \mathbb{R}^{n}\right)$ and $\operatorname{supp} f_{i} \subseteq \overline{O_{i} \cap \Omega}$.

- For $i \in\{1, \ldots, k\}$ we have $O_{i} \cap \Omega$ is strongly star-shaped. Lemma 3.13 ensures that $\operatorname{supp}\left(f_{i}\right)_{\theta} \subseteq O_{i} \cap \Omega$ for $\theta \in(0,1)$ and $\left(f_{i}\right)_{\theta} \rightarrow f_{i}$ in $H\left(L_{\partial}, \mathbb{R}^{n}\right)$ for $\theta \rightarrow 1$.

Let $\rho_{\epsilon}$ be a positive mollifier. Then $\rho_{\epsilon} * g \rightarrow g$ in $L^{2}\left(\mathbb{R}^{n}\right)$ for an arbitrary $g \in L^{2}\left(\mathbb{R}^{n}\right)$. Since $L_{\partial}\left(\rho_{\epsilon} * h\right)=\rho_{\epsilon} * L_{\partial} h$, we also have $\rho_{\epsilon} * h \rightarrow h$ in $H\left(L_{\partial}, \mathbb{R}^{n}\right)$ for $h \in H\left(L_{\partial}, \mathbb{R}^{n}\right)$ and since $\rho_{\epsilon} \in C^{\infty}\left(\mathbb{R}^{n}\right)$ we have $\rho_{\epsilon} * h \in C^{\infty}\left(\mathbb{R}^{n}\right)^{m_{2}}$.

For fixed $\theta \in(0,1)$ and $\epsilon$ sufficiently small, we can say $\operatorname{supp} \rho_{\epsilon} *\left(f_{i}\right)_{\theta} \subseteq O_{i} \cap \Omega$. Hence, by a diagonalization argument we find a sequence $\left(\rho_{\epsilon_{j}} *\left(f_{i}\right)_{\theta_{j}}\right)_{j \in \mathbb{N}}$ in $\mathcal{D}(\Omega)^{m_{2}}$ converging to $f_{i}$ in $H\left(L_{\partial}, \mathbb{R}^{n}\right)$. Doing this for every $i \in\{1, \ldots, k\}$ yields sequences $\left(f_{i, j}\right)_{j \in \mathbb{N}}$ in $\mathcal{D}(\Omega)^{m_{2}}$ converging to $f_{i}$ in $H\left(L_{\partial}, \mathbb{R}^{n}\right)$.

- For $f_{0}$ we have supp $f_{0} \subseteq \overline{O_{0}} \subseteq \Omega$ and by Lemma 3.15 there exists a sequence $\left(g_{l}\right)_{l \in \mathbb{N}}$ in $H\left(L_{\partial}, \mathbb{R}^{n}\right)$ that converges to $f_{0}$ in $H\left(L_{\partial}, \mathbb{R}^{n}\right)$ such that every $g_{l}$ has compact support in $\Omega$. Every $g_{l}$ can be approximated by $\rho_{\epsilon} * g_{l}$ for $\epsilon \rightarrow 0$ in $H\left(L_{\partial}, \mathbb{R}^{n}\right)$ and if $\epsilon$ is sufficiently small $\operatorname{supp} \rho_{\epsilon} * g_{l} \subseteq \Omega$. A diagonalization argument establishes a sequence $\left(f_{0, j}\right)_{j \in \mathbb{N}}$ in $\mathcal{D}(\Omega)^{m_{2}}$ that converges to $f_{0}$ in $H\left(L_{\partial}, \mathbb{R}^{n}\right)$.
Consequently, $\left(\sum_{i=0}^{k} f_{i, j}\right)_{j \in \mathbb{N}}$ is a sequence in $\mathcal{D}(\Omega)^{m_{2}}$ that converges to $\tilde{f}$ in $H\left(L_{\partial}, \mathbb{R}^{n}\right)$ and by Lemma 3.10 also in $H\left(L_{\partial}, \Omega\right)$.

Theorem 3.18. $\left.\mathcal{D}\left(\mathbb{R}^{n}\right)^{m_{2}}\right|_{\Omega}$ is dense in $H\left(L_{\partial}, \Omega\right)$.
Proof. Suppose $\left.\mathcal{D}\left(\mathbb{R}^{n}\right)^{m_{2}}\right|_{\Omega}$ is not dense in $H\left(L_{\partial}, \Omega\right)$. Then there exists a non zero $f \in H\left(L_{\partial}, \Omega\right)$ such that

$$
\begin{equation*}
\langle f, g\rangle_{H\left(L_{\partial}, \Omega\right)}=\langle f, g\rangle_{L^{2}}+\left\langle L_{\partial} f, L_{\partial} g\right\rangle_{L^{2}}=0 \quad \text { for all }\left.\quad g \in \mathcal{D}\left(\mathbb{R}^{n}\right)^{m_{2}}\right|_{\Omega} \tag{3.3}
\end{equation*}
$$

In particular, for an arbitrary $h \in \mathcal{D}(\Omega)^{m_{2}}$ we have

$$
\langle f, h\rangle_{\mathcal{D}^{\prime}, \mathcal{D}}=\langle f, h\rangle_{L^{2}}=-\left\langle L_{\partial} f, L_{\partial} h\right\rangle_{L^{2}}=-\left\langle L_{\partial} f, L_{\partial} h\right\rangle_{\mathcal{D}^{\prime}, \mathcal{D}}=\left\langle L_{\partial}^{\mathrm{H}} L_{\partial} f, h\right\rangle_{\mathcal{D}^{\prime}, \mathcal{D}}
$$

which implies that $f=L_{\partial}^{\mathrm{H}} L_{\partial} f \in L^{2}(\Omega)^{m_{2}}$ and $f_{0}:=L_{\partial} f \in H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$. Hence we can rewrite (3.3) as

$$
\langle L_{\partial}^{\mathrm{H}} \underbrace{L_{\partial} f}_{=f_{0}}, g\rangle_{L^{2}(\Omega)}+\langle\underbrace{L_{\partial} f}_{=f_{0}}, L_{\partial} g\rangle_{L^{2}(\Omega)}=0 \quad \text { for all }\left.g \in \mathcal{D}\left(\mathbb{R}^{n}\right)^{m_{2}}\right|_{\Omega}
$$

By Lemma 3.17 (switching the roles of $L_{\partial}$ and $L_{\partial}^{\mathrm{H}}$ ) we have $f_{0} \in H_{0}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$. Since $\mathcal{D}(\Omega)^{m_{1}}$ is dense in $H_{0}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$, there is a sequence $\left(f_{n}\right)_{n \in \mathbb{N}}$ in $\mathcal{D}(\Omega)^{m_{1}}$ converging to $f_{0}$ with respect to $\|\cdot\|_{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)}$. The fact $f=L_{\partial}^{\mathrm{H}} L_{\partial} f=L_{\partial}^{\mathrm{H}} f_{0}$ implies

$$
\begin{aligned}
\left\langle f_{0}, f_{n}\right\rangle_{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)} & =\left\langle f_{0}, f_{n}\right\rangle_{L^{2}}+\left\langle L_{\partial}^{\mathrm{H}} f_{0}, L_{\partial}^{\mathrm{H}} f_{n}\right\rangle_{L^{2}}=\left\langle L_{\partial} f, f_{n}\right\rangle_{L^{2}}+\left\langle f, L_{\partial}^{\mathrm{H}} f_{n}\right\rangle_{L^{2}} \\
& =\left\langle L_{\partial} f, f_{n}\right\rangle_{\mathcal{D}^{\prime}, \mathcal{D}}-\left\langle L_{\partial} f, f_{n}\right\rangle_{\mathcal{D}^{\prime}, \mathcal{D}} \\
& =0 .
\end{aligned}
$$

Since $\left\|f_{0}\right\|_{H\left(L_{\alpha}^{H}, \Omega\right)}^{2}=\lim _{n \rightarrow \infty}\left\langle f_{0}, f_{n}\right\rangle_{H\left(L_{,}^{H}, \Omega\right)}=0$, we have $f_{0}=0$, which implies $f=L_{\partial}^{\mathrm{H}} f_{0}=0$. Hence, $\left.\mathcal{D}\left(\mathbb{R}^{n}\right)^{m_{2}}\right|_{\Omega}$ is dense in $H\left(L_{\partial}, \Omega\right)$.
4. Port Hamiltonian systems. In this section we will introduce linear first order port-Hamiltonian systems on multidimensional spatial domains and illustrate the difficulties we want to overcome.
Definition 4.1. Let $m \in \mathbb{N}$ and $P=\left(P_{i}\right)_{i=1}^{n}$, where $P_{i}$ is a Hermitian $m \times m$ matrix. Moreover, let $\mathcal{H}: \Omega \rightarrow \mathbb{K}^{m \times m}$ be such that $\mathcal{H}(\zeta)^{\mathrm{H}}=\mathcal{H}(\zeta)$ and $c \mathrm{I} \leq \mathcal{H}(\zeta) \leq C \mathrm{I}$ for a.e. $\zeta \in \Omega$ and some constants $c, C>0$ independent of $\zeta$. Then we endow the space $\mathcal{X}_{\mathcal{H}}:=L^{2}(\Omega)^{m}$ with the scalar product

$$
\langle f, g\rangle_{\mathcal{X}_{\mathcal{H}}}:=\frac{1}{2}\langle\mathcal{H} f, g\rangle_{L^{2}(\Omega)^{m}}=\frac{1}{2} \int_{\Omega}\langle\mathcal{H}(\zeta) f(\zeta), g(\zeta)\rangle_{\mathbb{K}^{m}} \mathrm{~d} \lambda(\zeta) .
$$

We will refer to $\mathcal{X}_{\mathcal{H}}$ as the state space and to its elements as state variables or states. Furthermore, let $P_{0} \in \mathbb{K}^{m \times m}$ be such that $P_{0}^{\mathrm{H}}=-P_{0}$. Then we will call the differential equation

$$
\begin{align*}
\frac{\partial}{\partial t} x(t, \zeta) & =\sum_{i=1}^{n} \frac{\partial}{\partial \zeta_{i}} P_{i}(\mathcal{H}(\zeta) x(t, \zeta))+P_{0}(\mathcal{H}(\zeta) x(t, \zeta)), & & t \in \mathbb{R}_{+}, \zeta \in \Omega,  \tag{4.1}\\
x(0, \zeta) & =x_{0}(\zeta), & & \zeta \in \Omega
\end{align*}
$$

a linear, first order port-Hamiltonian system, where $x_{0} \in L^{2}(\Omega)^{m}$ is the initial state. The associated Hamiltonian $H: \mathcal{X}_{\mathcal{H}} \rightarrow \mathbb{R}_{+} \cup\{0\}$ is defined by

$$
H(x):=\langle x, x\rangle_{\mathcal{X}_{\mathcal{H}}}=\frac{1}{2} \int_{\Omega}\langle\mathcal{H}(\zeta) x(\zeta), x(\zeta)\rangle_{\mathbb{K}^{m}} \mathrm{~d} \lambda(\zeta)
$$

where $\mathcal{H}$ is called the Hamiltonian density.
In most applications the Hamiltonian describes the energy in the state space.
By the convention of regarding a function $x: \mathbb{R}_{+} \times \Omega \rightarrow \mathbb{K}^{m}$ as $x: \mathbb{R}_{+} \rightarrow$ $L^{2}\left(\Omega ; \mathbb{K}^{m}\right)$ by setting $x(t)=x(t,$.$) , we can rewrite the PDE (4.1) as$

$$
\dot{x}=\left(\sum_{i=1}^{n} \partial_{i} P_{i}+P_{0}\right) \mathcal{H} x=\left(P_{\partial}+P_{0}\right) \mathcal{H} x, \quad x(0)=x_{0},
$$

where $P_{\partial}$ is defined by Definition 3.2 replacing $L$ with $P$.
We want to add the following assumption on $P$.
Assumption 4.2. Let $m, m_{1}, m_{2} \in \mathbb{N}$ such that $m=m_{1}+m_{2}$ and let $L=\left(L_{i}\right)_{i=1}^{n}$ such that $L_{i} \in \mathbb{K}^{m_{1} \times m_{2}}$. Then we assume that $P=\left(P_{i}\right)_{i=1}^{n}$ has the block structure

$$
P_{i}=\left[\begin{array}{cc}
0 & L_{i} \\
L_{i}^{\mathrm{H}} & 0
\end{array}\right] .
$$

Assumption 4.2 implies that $P$ contains only Hermitian matrices. According to the block structure we split $x \in \mathbb{K}^{m}$ into $\left[\begin{array}{c}x_{x^{H}} \\ x_{L}\end{array}\right]$, where $x_{L^{H}}=\left(x_{i}\right)_{i=1}^{m_{1}}$ and $x_{L}=\left(x_{i}\right)_{i=m_{1}+1}^{m}$. We have the identities $H\left(P_{\partial}, \Omega\right)=H\left(L_{\partial}^{\mathrm{H}}, \Omega\right) \times H\left(L_{\partial}, \Omega\right)$,

$$
P_{\partial}=\left[\begin{array}{cc}
0 & L_{\partial} \\
L_{\partial}^{\mathrm{H}} & 0
\end{array}\right] \quad \text { and } \quad P_{\nu}=\left[\begin{array}{cc}
0 & L_{\nu} \\
L_{\nu}^{\mathrm{H}} & 0
\end{array}\right] .
$$

By Lemma 3.8 we have for $x, y \in H^{1}(\Omega)^{m}$

$$
\begin{align*}
\left\langle P_{\partial} x, y\right\rangle_{L^{2}(\Omega)}+\langle x, & \left.P_{\partial} y\right\rangle_{L^{2}(\Omega)} \\
& =\left\langle P_{\nu} \gamma_{0} x, \gamma_{0} y\right\rangle_{L^{2}(\partial \Omega)} \\
& =\left\langle\left[\begin{array}{cc}
0 & L_{\nu} \\
L_{\nu}^{\mathrm{H}} & 0
\end{array}\right] \gamma_{0}\left[\begin{array}{c}
x_{L^{\mathrm{H}}} \\
x_{L}
\end{array}\right], \gamma_{0}\left[\begin{array}{c}
y_{L^{\mathrm{H}}} \\
y_{L}
\end{array}\right]\right\rangle_{L^{2}(\partial \Omega)}  \tag{4.2}\\
& =\left\langle L_{\nu} \gamma_{0} x_{L}, \gamma_{0} y_{L^{\mathrm{H}}}\right\rangle_{L^{2}(\partial \Omega)}+\left\langle L_{\nu}^{\mathrm{H}} \gamma_{0} x_{L^{\mathrm{H}}}, \gamma_{0} y_{L}\right\rangle_{L^{2}(\partial \Omega)} \\
& =\left\langle L_{\nu} \gamma_{0} x_{L}, \gamma_{0} y_{L^{\mathrm{H}}}\right\rangle_{L^{2}(\partial \Omega)}+\left\langle\gamma_{0} x_{L^{\mathrm{H}}}, L_{\nu} \gamma_{0} y_{L}\right\rangle_{L^{2}(\partial \Omega)} .
\end{align*}
$$

Hence, $\mathcal{B}=L^{2}(\partial \Omega)^{m_{1}}, B_{1} x=L_{\nu} \gamma_{0} x_{L}$ and $B_{2} x=\gamma_{0} x_{L^{\text {H }}}$ is reminiscent of a boundary triple for $A_{0}^{*}=P_{\partial}\left(A_{0}=P_{\partial}^{*}\right.$ is skew-symmetric by Remark 3.7). However, we need to extend (4.2) for $x, y \in H\left(P_{\partial}, \Omega\right)$. In order to do this we have to introduce a new norm on $L^{2}(\partial \Omega)^{m_{1}}$, which will lead to the notion of quasi Gelfand triples.
5. Quasi Gelfand triples. Normally when we talk about Gelfand triples we have a Hilbert space $\mathcal{X}_{0}$ and another Hilbert space $\mathcal{X}_{+}$that can be continuously and densely embedded into $\mathcal{X}_{0}$. The third space $\mathcal{X}_{-}$is given by the completion of $\mathcal{X}_{0}$ with respect to

$$
\|g\|_{\mathcal{X}_{-}}:=\sup _{f \in \mathcal{X}_{+} \backslash\{0\}} \frac{|\langle g, f\rangle|_{\mathcal{X}_{0}}}{\|f\|_{\mathcal{X}_{+}}} .
$$

The duality between $\mathcal{X}_{+}$and $\mathcal{X}_{-}$is given by

$$
\langle g, f\rangle_{\mathcal{X}_{-}, \mathcal{X}_{+}}=\lim _{k \rightarrow \infty}\left\langle g_{k}, f\right\rangle_{\mathcal{X}_{0}}
$$

where $\left(g_{k}\right)_{k \in \mathbb{N}}$ is a sequence in $\mathcal{X}_{0}$ that converges to $g$ in $\mathcal{X}_{-}$. Details for "ordinary" Gelfand triple can be found in [6, ch. 2.1, p. 54] or in [13, ch. 2.9, p. 56].

We want to weaken the assumptions such that the norm of $\mathcal{X}_{+}$is not necessarily related to the norm of $\mathcal{X}_{0}$. This is in particular necessary for Maxwell's equations. In Example 8.10 we point out that is not possible to associate an "ordinary" Gelfand triple to the spatial differential operator of Maxwell's equations.

We will have the following setting: Let $\left(\mathcal{X}_{0},\langle., .\rangle_{\mathcal{X}_{0}}\right)$ be a Hilbert space and $\langle., .\rangle_{\mathcal{X}_{+}}$ another inner product (not necessarily related to $\langle.,.\rangle \mathcal{X}_{0}$ ) which is defined on a dense (w.r.t. $\|.\|_{\mathcal{X}_{0}}$ ) subspace $\tilde{D}_{+}$of $\mathcal{X}_{0}$. We denote the completion of $\tilde{D}_{+}$w.r.t. $\|.\|_{\mathcal{X}_{+}}=$ $\sqrt{\langle., .\rangle \mathcal{X}_{+}}$by $\mathcal{X}_{+}$. This completion is again a Hilbert space with the extension of $\langle., .\rangle_{\mathcal{X}_{+}}$, for which we use the same symbol. Now we have $\tilde{D}_{+}$is dense in $\mathcal{X}_{0}$ w.r.t. $\|.\|_{\mathcal{X}_{0}}$ and dense in $\mathcal{X}_{+}$w.r.t. $\|.\|_{\mathcal{X}_{+}}$.

Summarized:

- $\mathcal{X}_{0}$ Hilbert space endowed with $\left.\langle.,\rangle.\right\rangle_{\mathcal{X}_{0}}$.
- $\tilde{D}_{+}$dense subspace of $\mathcal{X}_{0}$ (w.r.t. $\|.\|_{\mathcal{X}_{0}}$ ).
- $\langle., .\rangle_{\mathcal{X}_{+}}$another inner product defined on $\tilde{D}_{+}$.
- $\mathcal{X}_{+}$completion of $\tilde{D}_{+}$with respect to $\|\cdot\|_{\mathcal{X}_{+}}$.

Example 5.1. Let $\mathcal{X}_{0}=\ell^{2}(\mathbb{Z} \backslash\{0\})$ with the standard inner product $\langle x, y\rangle_{\mathcal{X}_{0}}=$ $\sum_{n=1}^{\infty} x_{n} \overline{y_{n}}+x_{-n} \overline{y_{-n}}$. We define the inner product

$$
\langle x, y\rangle_{\mathcal{X}_{+}}:=\sum_{n=1}^{\infty} n^{2} x_{n} \overline{y_{n}}+\frac{1}{n^{2}} x_{-n} \overline{y_{-n}}
$$

and the set $\tilde{D}_{+}:=\left\{f \in \mathcal{X}_{0}:\|f\|_{\mathcal{X}_{+}}<+\infty\right\}$. Clearly, this inner product is welldefined on $\tilde{D}_{+}$. Let $e_{i}$ denote the sequence which is 1 on the $i$-th position and 0 elsewhere. Since $\left\{e_{i}: i \in \mathbb{Z} \backslash\{0\}\right\}$ is a orthonormal basis of $\mathcal{X}_{0}$ and contained in $\tilde{D}_{+}, \tilde{D}_{+}$is dense in $\mathcal{X}_{0}$. The sequence $\left(\sum_{i=1}^{n} e_{-i}\right)_{n \in \mathbb{N}}$ is a Cauchy sequence with respect to $\|\cdot\|_{\mathcal{X}_{+}}$but not with respect to $\|\cdot\|_{\mathcal{X}_{0}}$.
Definition 5.2. We define
$\|g\|_{\mathcal{X}_{-}}:=\sup _{f \in \tilde{D}_{+} \backslash\{0\}} \frac{\left|\langle g, f\rangle_{\mathcal{X}_{0}}\right|}{\|f\|_{\mathcal{X}_{+}}} \quad$ for $g \in \mathcal{X}_{0} \quad$ and $\quad D_{-}:=\left\{g \in \mathcal{X}_{0}:\|g\|_{\mathcal{X}_{-}}<+\infty\right\}$.
We denote the completion of $D_{-}$w.r.t. $\|.\|_{\mathcal{X}_{-}}$by $\mathcal{X}_{-}$. We will also denote the extension of $\|\cdot\|_{\mathcal{X}_{-}}$to $\mathcal{X}_{-}$by $\|\cdot\|_{\mathcal{X}_{-}}$.
Remark 5.3. By definition of $D_{-}$we can identify every $g \in D_{-}$with an element of $\mathcal{X}_{+}^{\prime}$ by the continuous extension of $f \in \tilde{D}_{+} \mapsto\langle g, f\rangle \mathcal{X}_{0}$ to $\mathcal{X}_{+}$. The completion $\mathcal{X}_{-}$ is isomorphic to the closure of $D_{-}$in $\mathcal{X}_{+}^{\prime}$ as $\|g\|_{\mathcal{X}_{+}^{\prime}}=\|g\|_{\mathcal{X}_{-}}$for $g \in D_{-}$.
Lemma 5.4. $D_{-}$is complete with respect to $\|g\|_{\mathcal{X}_{-} \cap \mathcal{X}_{0}}:=\sqrt{\|g\|_{\mathcal{X}_{0}}^{2}+\|g\|_{\mathcal{X}_{-}}^{2}}$.
Proof. Let $\left(g_{n}\right)_{n \in \mathbb{N}}$ be a Cauchy sequence in $D_{-}$with respect to $\|\cdot\|_{\mathcal{X}_{-} \cap \mathcal{X}_{0}}$. Then $\left(g_{n}\right)_{n \in \mathbb{N}}$ is a convergent sequence in $\mathcal{X}_{0}$ (w.r.t. $\|.\| \mathcal{X}_{0}$ ) and a Cauchy sequence in $D_{-}$ (w.r.t. $\|.\|_{\mathcal{X}_{-}}$). We denote the limit in $\mathcal{X}_{0}$ by $g_{0}$. By definition of $\|.\|_{\mathcal{X}_{-}}$we obtain for $f \in \tilde{D}_{+}$

$$
\left|\left\langle g_{0}, f\right\rangle_{\mathcal{X}_{0}}\right|=\lim _{n \rightarrow \infty}\left|\left\langle g_{n}, f\right\rangle_{\mathcal{X}_{0}}\right| \leq \lim _{n \rightarrow \infty}\left\|g_{n}\right\|_{\mathcal{X}_{-}}\|f\|_{\mathcal{X}_{+}} \leq C\|f\|_{\mathcal{X}_{+}}
$$

and consequently $g_{0} \in D_{-}$.
Let $\epsilon>0$ be arbitrary. Since $\left(g_{n}\right)_{n \in \mathbb{N}}$ is a Cauchy sequence with respect to $\|.\| \mathcal{X}_{-}$, there is an $n_{0} \in \mathbb{N}$ such that for all $f \in \tilde{D}_{+}$with $\|f\|_{\mathcal{X}_{+}}=1$

$$
\left|\left\langle g_{n}-g_{m}, f\right\rangle_{\mathcal{X}_{0}}\right| \leq \frac{\epsilon}{2}, \quad \text { if } \quad n, m \geq n_{0}
$$

holds true. Furthermore, for every $f \in \tilde{D}_{+}$there exists an $m_{f} \geq n_{0}$ such that $\left|\left\langle g_{0}-g_{m_{f}}, f\right\rangle_{\mathcal{X}_{0}}\right| \leq \frac{\epsilon\|f\|_{\mathcal{X}_{+}}}{2}$, because $g_{m} \rightarrow g_{0}$ w.r.t. $\|.\|_{\mathcal{X}_{0}}$. This yields

$$
\frac{\left|\left\langle g_{0}-g_{n}, f\right\rangle_{\mathcal{X}_{0}}\right|}{\|f\|_{\mathcal{X}_{+}}} \leq \frac{\left|\left\langle g_{0}-g_{m_{f}}, f\right\rangle_{\mathcal{X}_{0}}\right|}{\|f\|_{\mathcal{X}_{+}}}+\frac{\left|\left\langle g_{m_{f}}-g_{n}, f\right\rangle_{\mathcal{X}_{0}}\right|}{\|f\|_{\mathcal{X}_{+}}} \leq \epsilon, \quad \text { if } \quad n \geq n_{0} .
$$

Since the right-hand-side is independent of $f$, we obtain

$$
\left\|g_{0}-g_{n}\right\|_{\mathcal{X}_{-}}=\sup _{f \in \tilde{D}_{+} \backslash\{0\}} \frac{\left|\left\langle g_{0}-g_{n}, f\right\rangle_{\mathcal{X}_{0}}\right|}{\|f\|_{\mathcal{X}_{+}}} \leq \epsilon, \quad \text { if } \quad n \geq n_{0} .
$$

Hence, $g_{0}$ is also the limit of $\left(g_{n}\right)_{n \in \mathbb{N}}$ with respect to $\|\cdot\|_{\mathcal{X}_{-}}$and consequently the limit of $\left(g_{n}\right)_{n \in \mathbb{N}}$ with respect to $\|\cdot\|_{\mathcal{X}} \cap \mathcal{X}_{0}$.

Lemma 5.5. The embedding $\tilde{\iota}_{+}: \tilde{D}_{+} \subseteq \mathcal{X}_{+} \rightarrow \mathcal{X}_{0}, f \mapsto f$ is a densely defined operator with ran $\tilde{\iota}_{+}$is dense in $\mathcal{X}_{0}$ and $\operatorname{ker} \tilde{\tau}_{+}=\{0\}$. Furthermore, the embedding $\iota_{-}: D_{-} \subseteq \mathcal{X}_{-} \rightarrow \mathcal{X}_{0}, g \mapsto g$ is closed and $\operatorname{ker} \iota_{-}=\{0\}$.
Proof. By assumption on $\tilde{D}_{+}$and definition of $\mathcal{X}_{+}$the embedding $\tilde{\iota}_{+}$is densely defined and has a dense range. Clearly, $\operatorname{ker} \tilde{\iota}_{+}=\{0\}$ and $\operatorname{ker} \iota_{-}=\{0\}$. By Lemma 5.4 $\iota_{-}$is closed.

Lemma 5.6. Let $\tilde{\iota}_{+}^{*}$ denote the adjoint relation (from $\mathcal{X}_{0}$ to $\mathcal{X}_{+}$) of the embedding mapping $\tilde{\iota}_{+}$in the previous lemma. Then $\tilde{\iota}_{+}^{*}$ is an operator (single-valued, i.e. $\left.\operatorname{mul} \tilde{\iota}_{+}^{*}=\{0\}\right)$ and $\operatorname{ker} \tilde{\iota}_{+}^{*}=\{0\}$. Its domain coincides with $D_{-}$and $\tilde{\iota}_{+}^{*} \iota_{-}: D_{-} \subseteq$ $\mathcal{X}_{-} \rightarrow \mathcal{X}_{+}$is isometric.

If $\operatorname{ker} \tilde{\iota}_{+}=\{0\}$, then $\operatorname{ran} \tilde{\iota}_{+}^{*}$ is dense in $\mathcal{X}_{+}$.
Proof. The density of the domain of $\tilde{\iota}_{+}$yields mul $\tilde{\iota}_{+}^{*}=\left(\operatorname{dom} \tilde{\iota}_{+}\right)^{\perp}=\{0\}$, and $\overline{\operatorname{ran} \tilde{\iota}_{+}}{ }^{\mathcal{X}_{0}}=\mathcal{X}_{0}$ yields $\operatorname{ker} \tilde{\iota}_{+}^{*}=\{0\}$. The following equivalences show dom $\tilde{\iota}_{+}^{*}=D_{-}$:

$$
\begin{aligned}
g \in \operatorname{dom} \tilde{\iota}_{+}^{*} & \Leftrightarrow\left\langle g, \tilde{c}_{+} f\right\rangle_{\mathcal{X}_{0}} \text { is continuous in } f \in \tilde{D}_{+} \text {w.r.t. }\|\cdot\|_{\mathcal{X}_{+}} \\
& \Leftrightarrow \sup _{f \in \tilde{D}_{+} \backslash\{0\}} \frac{\left|\langle g, f\rangle_{\mathcal{X}_{0}}\right|}{\|f\|_{\mathcal{X}_{+}}}<+\infty \\
& \Leftrightarrow g \in D_{-}
\end{aligned}
$$

For $g \in D_{-} \subseteq \mathcal{X}_{-}$we have

$$
\|g\|_{\mathcal{X}_{-}}=\sup _{f \in \tilde{D}+\backslash\{0\}} \frac{\left|\left\langle\iota_{-} g, f\right\rangle_{\mathcal{X}_{0}}\right|}{\|f\|_{\mathcal{X}_{+}}}=\sup _{f \in \tilde{D}_{+} \backslash\{0\}} \frac{\left|\left\langle\tilde{\iota}_{+}^{*} \iota_{-} g, f\right\rangle_{\mathcal{X}_{+}}\right|}{\|f\|_{\mathcal{X}_{+}}}=\left\|\tilde{\tau}_{+}^{*} \iota_{-} g\right\|_{\mathcal{X}_{+}},
$$

which proves that $\tilde{\iota}_{+}^{*} \iota_{-}$is isometric.
If $\operatorname{ker} \overline{\tilde{\iota}_{+}}=\{0\}$, then the following equation implies the density of ran $\tilde{\iota}_{+}^{*}$ in $\mathcal{X}_{+}$

$$
\{0\}=\operatorname{ker} \overline{\tilde{\iota}_{+}}=\operatorname{ker} \tilde{\iota}_{+}^{* *}=\left(\operatorname{ran} \tilde{\iota}_{+}^{*}\right)^{\perp}
$$

Remark 5.7. As mentioned in Remark 5.3 every $g \in D_{-}$can be regarded as an element of $\mathcal{X}_{+}^{\prime}$ by the continuous extension of $\tilde{D}_{+} \ni f \mapsto\left\langle g, \tilde{\iota}_{+} f\right\rangle_{\mathcal{X}_{0}}$ on $\mathcal{X}_{+}$. Since $D_{-}=\operatorname{dom} \tilde{\iota}_{+}^{*}$, this extension equals $\left\langle\tilde{\iota}_{+}^{*} g, .\right\rangle_{\mathcal{X}_{+}}$.
Proposition 5.8. The following assertions are equivalent.
(i) There is a Hausdorff topological vector space $(Z, \mathcal{T})$ and two continuous embeddings $\phi_{\mathcal{X}_{+}}: \mathcal{X}_{+} \rightarrow Z$ and $\phi_{\mathcal{X}_{0}}: \mathcal{X}_{0} \rightarrow Z$ such that the diagram

commutes.
(ii) If $\tilde{D}_{+} \ni f_{n} \rightarrow 0$ w.r.t. $\|.\|_{\mathcal{X}_{+}}$and $\lim _{n \rightarrow \infty} f_{n}$ exists w.r.t. $\|.\|_{\mathcal{X}_{0}}$, then this limit is also 0 and if $\tilde{D}_{+} \ni f_{n} \rightarrow 0$ w.r.t. $\|.\|_{\mathcal{X}_{0}}$ and $\lim _{n \rightarrow \infty} f_{n}$ exists w.r.t. $\|\cdot\|_{\mathcal{X}_{+}}$, then this limit is also 0 .
(iii) $\tilde{\iota}_{+}: \tilde{D}_{+} \subseteq \mathcal{X}_{+} \rightarrow \mathcal{X}_{0}, f \mapsto f$ is closable (as an operator) and its closure is injective.
(iv) $D_{-}$is dense in $\mathcal{X}_{0}$ and dense in $\mathcal{X}_{+}^{\prime}$.

Proof. (i) $\Rightarrow$ (ii): Let $\left(f_{n}\right)_{n \in \mathbb{N}}$ be a sequence in $\tilde{D}_{+}$such that $f_{n} \rightarrow \hat{f}$ w.r.t. $\mathcal{X}_{+}$and $f_{n} \rightarrow f$ w.r.t. $\mathcal{X}_{0}$. Since $\mathcal{T}$ is coarser than both of the topologies induced by these norms, we also have


Since $\mathcal{T}$ is Hausdorff, we conclude $f=\hat{f}$. Hence, if either $\hat{f}$ or $f$ is 0 , then also the other is 0 .
(ii) $\Rightarrow$ (iii): If $\left(f_{n}, f_{n}\right)_{n \in \mathbb{N}}$ is a sequence in $\tilde{\iota}_{+}$that converges to $(0, f) \in \mathcal{X}_{+} \times \mathcal{X}_{0}$, then $f=0$ by (ii). Hence, mul $\tilde{\iota}_{+}=\{0\}$ and consequently $\tilde{\iota}_{+}$is closable. On the other hand, if $\left(f_{n}, f_{n}\right)_{n \in \mathbb{N}}$ is a sequence in $\tilde{\iota}_{+}$that converges to $(f, 0)$, then $f=0$ by (ii). Consequently, $\operatorname{ker} \overline{\tilde{\iota}_{+}}=\{0\}$ and $\overline{\tilde{\iota}_{+}}$is injective.
(iii) $\Rightarrow$ (iv): We have $\left(\operatorname{dom} \tilde{\iota}_{+}^{*}\right)^{\perp}=\operatorname{mul} \tilde{\iota}_{+}^{* *}=\operatorname{mul} \overline{\tilde{\iota}_{+}}$. Since $\tilde{\iota}_{+}$is closable, we have mul $\tilde{\iota}_{+}=\{0\}$, which yields that dom $\tilde{\iota}_{+}^{*}$ is dense in $\mathcal{X}_{0}$. By Lemma 5.6 dom $\tilde{\iota}_{+}^{*}$ coincides with $D_{-}$.

The second assertion of Lemma 5.6 yields that $\operatorname{ran} \tilde{\iota}_{+}^{*}=\tilde{\iota}_{+}^{*} D_{-}$is dense in $\mathcal{X}_{+}$. As mentioned in Remark 5.7 every element $g \in D_{-}$can be identified with $\left\langle\tilde{\imath}_{+}^{*} g,.\right\rangle \mathcal{X}_{+} \in$ $\mathcal{X}_{+}^{\prime}$. Therefore, the density of $\tilde{\iota}_{+}^{*} D_{-}$in $\mathcal{X}_{+}$implies the density of $D_{-}$in $\mathcal{X}_{+}^{\prime}$, because $f \mapsto\langle f, .\rangle_{\mathcal{X}_{+}}$is a unitary mapping between $\mathcal{X}_{+}$and $\mathcal{X}_{+}^{\prime}$.

$$
(\text { iv }) \Rightarrow(\text { i }): \text { Let } Y:=D_{-} \text {be equipped with }\|g\|_{Y}:=\|g\|_{\mathcal{X}_{-} \cap \mathcal{X}_{0}}=\sqrt{\|g\|_{\mathcal{X}_{-}}^{2}+\|g\|_{\mathcal{X}_{0}}^{2}} .
$$

We define $Z:=Y^{\prime}$ as the (anti)dual of $Y$. Then we have

$$
\begin{aligned}
& \left|\langle f, g\rangle_{\mathcal{X}_{0}}\right| \leq\|f\|_{\mathcal{X}_{0}}\|g\|_{\mathcal{X}_{0}} \leq\|f\|_{\mathcal{X}_{0}}\|g\|_{Y} \quad \text { for } \quad f \in \mathcal{X}_{0}, g \in Y \\
& \text { and } \quad\left|\left\langle f, \tilde{\iota}_{+}^{*} g\right\rangle_{\mathcal{X}_{+}}\right| \leq\|f\|_{\mathcal{X}_{+}} \underbrace{\left\|\tilde{\iota}_{+}^{*} g\right\|_{\mathcal{X}_{+}}}_{=\|g\|_{\mathcal{X}_{-}}} \leq\|f\|_{\mathcal{X}_{+}}\|g\|_{Y} \quad \text { for } \quad f \in \mathcal{X}_{+}, g \in Y \text {. }
\end{aligned}
$$

Hence, $\phi_{\mathcal{X}_{0}}: f \mapsto\langle f, .\rangle_{\mathcal{X}_{0}}$ and $\phi_{\mathcal{X}_{+}}: f \mapsto\left\langle f, \tilde{i}_{+}^{*} .\right\rangle_{\mathcal{X}_{+}}$are continuous mappings from $\mathcal{X}_{0}$ and $\mathcal{X}_{+}$, respectively, into $Z$. The injectivity of these mappings follows from the density of $D_{-}$in $\mathcal{X}_{0}$ and $D_{-}$in $\mathcal{X}_{+}^{\prime}\left(\tilde{\iota}_{+}^{*} D_{-}\right.$dense in $\left.\mathcal{X}_{+}\right)$, respectively. For $f \in \tilde{D}_{+}$ we have

$$
\phi_{\mathcal{X}_{+}} f=\left\langle f, \tilde{\iota}_{+}^{*} \cdot\right\rangle_{\mathcal{X}_{+}}=\left\langle\tilde{\iota}_{+} f, .\right\rangle_{\mathcal{X}_{0}}=\phi_{\mathcal{X}_{0}} \circ \tilde{\iota}_{+} f
$$

and consequently the diagram in (i) commutes.

If one and therefore all assertions in Proposition 5.8 are satisfied, then $\mathcal{X}_{+} \cap$ $\mathcal{X}_{0}$ is defined as the intersection in $Z$ and complete with the norm $\|\cdot\|_{\mathcal{X}_{+} \cap \mathcal{X}_{0}}:=$ $\sqrt{\|\cdot\|_{\mathcal{X}_{+}}^{2}+\|\cdot\|_{\mathcal{X}_{0}}^{2}}$. Moreover, we define $D_{+}$as the closure of $\tilde{D}_{+}$in $\mathcal{X}_{+} \cap \mathcal{X}_{0}$ (w.r.t. $\|.\|_{\mathcal{X}_{+} \cap \mathcal{X}_{0}}$ ). Note that although $\mathcal{X}_{+} \cap \mathcal{X}_{0}$ may depend on $Z, D_{+}$is independent of $Z$. We will denote the extension of $\tilde{\iota}_{+}$to $D_{+}$by $\iota_{+}$, which can be expressed by $\iota_{+}=\overline{\tilde{\iota}_{+}}$. The adjoint $\iota_{+}^{*}$ coincides with $\tilde{\iota}_{+}^{*}$. Also $D_{-}$does not change, if we replace $\tilde{D}_{+}$by $D_{+}$in Definition 5.2 and all previous results in this section also hold for $D_{+}$ and $\iota_{+}$instead of $\tilde{D}_{+}$and $\tilde{\iota}_{+}$, respectively. If $\tilde{\iota}_{+}$is already closed, then $D_{+}=\tilde{D}_{+}$.

Lemma 5.9. Let one assertion in Proposition 5.8 be satisfied. Let $Z=Y^{\prime}$, where $Y=D_{-}$endowed with $\|g\|_{Y}:=\|g\|_{\mathcal{X}_{-} \cap \mathcal{X}_{0}}=\sqrt{\|g\|_{\mathcal{X}_{-}}^{2}+\|g\|_{\mathcal{X}_{0}}^{2}}$ (from Proposition 5.8 (iv) $\Rightarrow(\mathrm{i})$ ). Then we have the following characterization for $D_{+}$:

- $D_{+}=\operatorname{dom} \iota_{-}^{*}$,
- $D_{+}=\mathcal{X}_{+} \cap \mathcal{X}_{0}$ in $Y^{\prime}$.

Proof. Note that for $g \in D_{-}$we have $g=\left(\iota_{+}^{*}\right)^{-1} \iota_{+}^{*} g$ and that $\iota_{+}^{*} \iota_{-}$is isometric from dom $\iota_{-}$onto $\operatorname{dom}\left(\iota_{+}^{*}\right)^{-1}$. The following equivalences show the first assertion:

$$
\begin{aligned}
f \in \operatorname{dom} \iota_{-}^{*} & \Leftrightarrow D_{-} \ni g \mapsto\left\langle f, \iota_{-} g\right\rangle_{\mathcal{X}_{0}} \text { is continious w.r.t. }\|\cdot\|_{\mathcal{X}_{-}} \\
& \Leftrightarrow D_{-} \ni g \mapsto\left\langle f,\left(\iota_{+}^{*}\right)^{-1} \iota_{+}^{*} \iota_{-} g\right\rangle_{\mathcal{X}_{0}} \text { is continious w.r.t. }\|\cdot\|_{\mathcal{X}_{-}} \\
& \Leftrightarrow f \in \operatorname{dom}\left(\left(\iota_{+}^{*}\right)^{-1}\right)^{*}=\operatorname{dom} \iota_{+}^{-1}=\operatorname{ran} \iota_{+}=D_{+}
\end{aligned}
$$

We define $P_{+}:=\mathcal{X}_{+} \cap \mathcal{X}_{0}$ and we define $P_{-}$analogously to $D_{-}$in Definition 5.2:

$$
\|g\|_{P_{-}}:=\sup _{f \in P_{+} \backslash\{0\}} \frac{\left|\langle g, f\rangle_{\mathcal{X}_{0}}\right|}{\|f\|_{\mathcal{X}_{+}}} \quad \text { and } \quad P_{-}:=\left\{g \in \mathcal{X}_{0}:\|g\|_{P_{-}}<+\infty\right\}
$$

Clearly, $\|g\|_{\mathcal{X}_{-}} \leq\|g\|_{P_{-}}$for $g \in P_{-}$and consequently $P_{-} \subseteq D_{-}$. Furthermore, we can define $\iota_{P_{+}}: P_{+} \subseteq \mathcal{X}_{+} \rightarrow \mathcal{X}_{0}, f \mapsto f$ analogously to $\tilde{\iota}_{+}$. Note that $\iota_{P_{+}}$is closed due the completeness of $\left(\mathcal{X}_{+} \cap \mathcal{X}_{0},\|\cdot\|_{\mathcal{X}_{+} \cap \mathcal{X}_{0}}\right)$. Then we have $\operatorname{dom} \iota_{P_{+}}^{*}=P_{-}$and $\tilde{\iota}_{+} \subseteq \iota_{P_{+}}$and therefore $\iota_{P_{+}}^{*} \subseteq \tilde{\iota}_{+}^{*}$. For $g \in D_{-}$and $f \in P_{+}$we have, by definition of $P_{+}=\mathcal{X}_{+} \cap \mathcal{X}_{0}$ in $Z$,

$$
\left|\langle g, f\rangle_{\mathcal{X}_{0}}\right|=\left|\left\langle\tilde{\imath}_{+} g, f\right\rangle_{\mathcal{X}_{+}}\right| \leq\left\|\tilde{\iota}_{+} g\right\|_{\mathcal{X}_{+}}\|f\|_{\mathcal{X}_{+}}=\|g\|_{\mathcal{X}_{-}}\|f\|_{\mathcal{X}_{+}},
$$

which yields $\|g\|_{P_{-}} \leq\|g\|_{\mathcal{X}_{-}}$. Hence, $P_{-}=D_{-}, \iota_{P_{+}}^{*}=\tilde{\iota}_{+}^{*}$ and $\iota_{P_{+}}=\overline{\tilde{\iota}_{+}}$, which is equivalent to $P_{+}=\mathcal{X}_{+} \cap \mathcal{X}_{0}={\overline{D_{+}}}^{\mathcal{X}_{+} \cap \mathcal{X}_{0}}=D_{+}$.

Theorem 5.10. Let one assertion in Proposition 5.8 be satisfied. Then the mapping $\iota_{+}^{*} \iota_{-}$can be uniquely extended to a isometric and surjective operator $\Psi: \mathcal{X}_{-} \rightarrow \mathcal{X}_{+}$. In particular $\mathcal{X}_{-}$is a Hilbert space whose (original) norm is induced by $\langle g, f\rangle_{\mathcal{X}_{-}}:=$ $\langle\Psi g, \Psi f\rangle_{\mathcal{X}_{+}}$and $\Psi$ is unitary.

Proof. By Lemma $5.6 \iota_{+}^{*} \iota_{-}: D_{-} \subseteq \mathcal{X}_{-} \rightarrow \mathcal{X}_{+}$is an isometry with dense range, since $\iota_{+}$is closed and injective by assumption. Since $D_{-}$is dense in $\mathcal{X}_{-}$by construction, we can extend $\iota_{+}^{*} \iota_{-}$by continuity to $\mathcal{X}_{-}$. We denote this extension by $\Psi$. For an arbitrary $g \in \mathcal{X}_{-}$there exists a sequence $\left(g_{n}\right)_{n \in \mathbb{N}}$ in $D_{-}$that converges to $g$ (w.r.t. $\left.\|\cdot\|_{\mathcal{X}_{-}}\right)$. Hence,

$$
\|\Psi g\|_{\mathcal{X}_{+}}=\lim _{n \rightarrow \infty}\left\|\Psi g_{n}\right\|_{\mathcal{X}_{+}}=\lim _{n \rightarrow \infty}\left\|\iota_{+}^{*} \iota_{-} g_{n}\right\|_{\mathcal{X}_{+}}=\lim _{n \rightarrow \infty}\left\|g_{n}\right\|_{\mathcal{X}_{-}}=\|g\|_{\mathcal{X}_{-}}
$$

This yields that $\Psi$ is isometric and $\operatorname{ran} \Psi$ is closed in $\mathcal{X}_{+}$. Since $\operatorname{ran} \Psi$ also contains the dense subspace $\operatorname{ran} \iota_{+}^{*}$, the mapping $\Psi$ is surjective.

Clearly, this implies that $\|.\|_{\mathcal{X}_{-}}$is induced by the inner product $\langle., .\rangle_{\mathcal{X}_{-}}=\langle\Psi ., \Psi .\rangle_{\mathcal{X}_{+}}$ and $\mathcal{X}_{-}$Hilbert space endowed with this inner product. Moreover, $\Psi$ is then unitary.

Corollary 5.11. If one assertion in Proposition 5.8 is satisfied, then $\mathcal{X}_{-}$can be identified with the (anti)dual space of $\mathcal{X}_{+}$by

$$
\Lambda:\left\{\begin{array}{rll}
\mathcal{X}_{-} & \rightarrow \mathcal{X}_{+}^{\prime} \\
g & \mapsto & \langle\Psi g, .\rangle_{\mathcal{X}_{+}}
\end{array}\right.
$$

where $\Psi$ is the mapping from Theorem 5.10.
If $\mathcal{X}_{1}$ and $\mathcal{X}_{2}$ are Hilbert spaces and $\mathcal{X}_{2}$ can be identified with the dual space of $\mathcal{X}_{1}$ by a unitary mapping $\Lambda: \mathcal{X}_{2} \rightarrow \mathcal{X}_{1}^{\prime}$, then we define

$$
\langle g, f\rangle_{\mathcal{X}_{2}, \mathcal{X}_{1}}:=\langle\Lambda g, f\rangle_{\mathcal{X}_{1}^{\prime}, \mathcal{X}_{1}}=(\Lambda g)(f) .
$$



Figure 1. Illustration of a quasi Gelfand triple, where $D_{+}=$ $\operatorname{dom} \iota_{+}$and $D_{-}=\operatorname{dom} \iota_{-}$.

Remark 5.12. For $f \in D_{+}$and $g \in D_{-}$we have

$$
\langle g, f\rangle_{\mathcal{X}_{-}, \mathcal{X}_{+}}=\langle\Psi g, f\rangle_{\mathcal{X}_{+}}=\left\langle\iota_{+}^{*} \iota_{-} g, f\right\rangle_{\mathcal{X}_{+}}=\left\langle\iota_{-} g, \iota_{+} f\right\rangle_{\mathcal{X}_{0}}=\langle g, f\rangle_{\mathcal{X}_{0}} .
$$

Since these two sets are dense in $\mathcal{X}_{+}$and $\mathcal{X}_{-}$respectively, we have for $f \in \mathcal{X}_{+}$and $g \in \mathcal{X}_{-}$

$$
\langle g, f\rangle_{\mathcal{X}_{-}, \mathcal{X}_{+}}=\lim _{(n, m) \rightarrow(\infty, \infty)}\left\langle g_{n}, f_{m}\right\rangle_{\mathcal{X}_{0}}
$$

where $\left(f_{m}\right)_{m \in \mathbb{N}}$ is a sequence in $D_{+}$that converges to $f$ in $\mathcal{X}_{+}$and $\left(g_{n}\right)_{n \in \mathbb{N}}$ is a sequence in $D_{-}$that converges to $g$ in $\mathcal{X}_{-}$.

Definition 5.13. Let $\mathcal{X}_{+}, \mathcal{X}_{0}$ and $\mathcal{X}_{-}$be Hilbert spaces, where $\mathcal{X}_{-}$can be identified with $\mathcal{X}_{+}^{\prime}$. Furthermore, let $\iota_{+}: \operatorname{dom} \iota_{+} \subseteq \mathcal{X}_{+} \rightarrow \mathcal{X}_{0}$ and $\iota_{-}: \operatorname{dom} \iota_{-} \subseteq \mathcal{X}_{-} \rightarrow \mathcal{X}_{0}$ be densely defined, closed, and injective linear mappings with dense range. We call $\left(\mathcal{X}_{+}, \mathcal{X}_{0}, \mathcal{X}_{-}\right)$a quasi Gelfand triple, if

$$
\begin{equation*}
\langle g, f\rangle_{\mathcal{X}_{-}, \mathcal{X}_{+}}=\left\langle\iota_{-} g, \iota_{+} f\right\rangle_{\mathcal{X}_{0}} \tag{5.1}
\end{equation*}
$$

for all $f \in \operatorname{dom} \iota_{+}$and $g \in \operatorname{dom} \iota_{-}$, and $\operatorname{dom} \iota_{+}^{*}=\operatorname{ran} \iota_{-}$. The space $\mathcal{X}_{0}$ will be referred as pivot space. We define $\mathcal{X}_{+} \cap \mathcal{X}_{0}:=\operatorname{ran} \iota_{+}$and $\mathcal{X}_{-} \cap \mathcal{X}_{0}:=\operatorname{ran} \iota_{-}$.

Figure 1 illustrates the setting of a quasi Gelfand triple. Since $\mathcal{X}_{-}$can be identified with $\mathcal{X}_{+}^{\prime}$ and $\mathcal{X}_{+}^{\prime}$ can be identified with $\mathcal{X}_{+}$, there exists a unitary operator $\Psi: \mathcal{X}_{-} \rightarrow \mathcal{X}_{+}$. In fact, by (5.1) this $\Psi$ is the extension of $\iota_{+}^{*} \iota_{-}$, which already appeared in Theorem 5.10. We will show this in detail in Proposition 5.16. We will call $\Psi$ the duality map of the quasi Gelfand triple.

In contrast to "ordinary" Gelfand triple, the setting for quasi Gelfand triple is somehow "symmetric", i.e. the roles of $\mathcal{X}_{+}$and $\mathcal{X}_{-}$are interchangeable, since neither $\iota_{+}$nor $\iota_{-}$have to be continuous, as indicated in the beginning of this section.
Lemma 5.14. Let $\left(\mathcal{X}_{+}, \mathcal{X}_{0}, \mathcal{X}_{-}\right)$with $\iota_{+}$and $\iota_{-}$satisfy all conditions of Definition 5.13 except dom $\iota_{+}^{*}=\operatorname{ran} \iota_{-}$. Then

$$
\operatorname{dom} \iota_{+}^{*}=\operatorname{ran} \iota_{-} \quad \Leftrightarrow \quad \operatorname{dom} \iota_{-}^{*}=\operatorname{ran} \iota_{+}
$$

In particular, if $\left(\mathcal{X}_{+}, \mathcal{X}_{0}, \mathcal{X}_{-}\right)$is a quasi Gelfand triple, then also $\operatorname{dom} \iota_{-}^{*}=\operatorname{ran} \iota_{+}$ holds true.

The proof of this is basically the first part of the proof of Lemma 5.9.

Proof. By (5.1), it is clear that dom $\iota_{+}^{*} \supseteq \operatorname{ran} \iota_{-}$and dom $\iota_{-}^{*} \supseteq \operatorname{ran} \iota_{+}$holds. Moreover, for $f \in \operatorname{dom} \iota_{+}, g \in \operatorname{dom} \iota_{-}$and the duality mapping $\Psi$ we have

$$
\left\langle g, \Psi^{*} f\right\rangle_{\mathcal{X}_{-}}=\langle\Psi g, f\rangle_{\mathcal{X}_{+}}=\langle g, f\rangle_{\mathcal{X}_{-}, \mathcal{X}_{+}}=\left\langle\iota_{-} g, \iota_{+} f\right\rangle_{\mathcal{X}_{0}}
$$

which implies $\iota_{+}^{*} \iota_{-} \subseteq \Psi$ and $\iota_{-}^{*} \iota_{+} \subseteq \Psi^{*}$. In particular, both $\iota_{+}^{*} \iota_{-}$and $\iota_{-}^{*} \iota_{+}$are isometric.

Let dom $\iota_{+}^{*}=\operatorname{ran} \iota_{-}$. Then $\iota_{+}^{*} \iota_{-}$is isometric from dom $\iota_{-}$onto $\operatorname{dom}\left(\iota_{+}^{*}\right)^{-1}$. The following equivalences

$$
\begin{aligned}
f \in \operatorname{dom} \iota_{-}^{*} & \Leftrightarrow \operatorname{dom} \iota_{-} \ni g \mapsto\left\langle f, \iota_{-} g\right\rangle_{\mathcal{X}_{0}} \text { is continuous w.r.t. }\|\cdot\|_{\mathcal{X}_{-}} \\
& \Leftrightarrow \operatorname{dom} \iota_{-} \ni g \mapsto\left\langle f,\left(\iota_{+}^{*}\right)^{-1} \iota_{+}^{*} \iota_{-} g\right\rangle_{\mathcal{X}_{0}} \text { is continuous w.r.t. }\|\cdot\|_{\mathcal{X}_{-}} \\
& \Leftrightarrow f \in \operatorname{dom}\left(\left(\iota_{+}^{*}\right)^{-1}\right)^{*}=\operatorname{dom} \iota_{+}^{-1}=\operatorname{ran} \iota_{+}
\end{aligned}
$$

imply $\operatorname{dom} \iota_{-}^{*}=\operatorname{ran} \iota_{+}$.
The other implication follows analogously.
Lemma 5.15. Let $\left(\mathcal{X}_{+}, \mathcal{X}_{0}, \mathcal{X}_{-}\right)$with $\iota_{+}$and $\iota_{-}$satisfy all conditions of Definition 5.13 except $\operatorname{dom} \iota_{+}^{*}=\operatorname{ran} \iota_{-}$. Then there exists an extension $\hat{\iota}_{-}$of $\iota_{-}$that respects (5.1) and satisfies $\operatorname{dom} \iota_{+}^{*}=\operatorname{ran} \hat{\iota}_{-}$. In particular, $\left(\mathcal{X}_{+}, \mathcal{X}_{0}, \mathcal{X}_{-}\right)$with $\iota_{+}$ and $\hat{\iota}_{-}$forms a quasi Gelfand triple.

Proof. Note that

$$
g \in \operatorname{dom} \iota_{+}^{*} \Leftrightarrow \operatorname{dom} \iota_{+} \ni f \mapsto\left\langle g, \iota_{+} f\right\rangle_{\mathcal{X}_{0}} \text { is continuous w.r.t. }\|.\|_{\mathcal{X}_{+}} .
$$

Hence, for $g \in \operatorname{dom} \iota_{+}^{*}$ there exists an $h \in \mathcal{X}_{-}$such that

$$
\begin{equation*}
\left\langle g, \iota_{+} f\right\rangle_{\mathcal{X}_{0}}=\langle h, f\rangle_{\mathcal{X}_{-}, \mathcal{X}_{+}} \quad \text { for all } \quad f \in \operatorname{dom} \iota_{+} . \tag{5.2}
\end{equation*}
$$

We define $\phi(g):=h$ for $g \in \operatorname{dom} \iota_{+}$. Clearly, $\phi(g)=\iota_{-}^{-1} g$ for $g \in \operatorname{ran} \iota_{-}$. Therefore, $\hat{\iota}_{-}:=\phi^{-1}$ is an extension of $\iota_{-}$that satisfies dom $\iota_{+}^{*}=\operatorname{ran} \hat{\iota}_{-}$. Moreover, by (5.2) we have $\left\langle\hat{\imath}_{-} g, \iota_{+} f\right\rangle_{\mathcal{X}_{0}}=\langle g, f\rangle_{\mathcal{X}_{-}, \mathcal{X}_{+}}$for $g \in \operatorname{dom} \hat{\iota}_{-}$and $f \in \operatorname{dom} \iota_{+}$.

Proposition 5.16. Let $\left(\mathcal{X}_{+}, \mathcal{X}_{0}, \mathcal{X}_{-}\right)$be a quasi Gelfand triple and $\Psi: \mathcal{X}_{-} \rightarrow \mathcal{X}_{+}$ be its duality map. Then

$$
\overline{\iota_{+}^{*} \iota_{-}}=\Psi, \quad \overline{\iota_{-}^{*} \iota_{+}}=\Psi^{*}, \quad \iota_{+}^{*}=\Psi \iota_{-}^{-1} \quad \text { and } \quad \iota_{-}^{*}=\Psi^{*} \iota_{+}^{-1}
$$

Proof. Let $f \in \operatorname{dom} \iota_{+}$and $g \in \operatorname{dom} \iota_{-}$. Then

$$
\langle\Psi g, f\rangle_{\mathcal{X}_{+}}=\langle g, f\rangle_{\mathcal{X}_{-}, \mathcal{X}_{+}}=\left\langle\iota_{-} g, \iota_{+} f\right\rangle_{\mathcal{X}_{0}}=\left\langle\iota_{+}^{*} \iota_{-} g, f\right\rangle_{\mathcal{X}_{+}} .
$$

Since dom $\iota_{+}$is dense in $\mathcal{X}_{+}$, we have $\Psi g=\iota_{+}^{*} \iota_{-} g$ for all $g \in \operatorname{dom} \iota_{-}$. Applying $\iota_{-}^{-1}$ on both sides gives $\Psi \iota_{-}^{-1}=\iota_{+}^{*}$. Moreover, the density of dom $\iota_{-}$in $\mathcal{X}_{-}$yields $\Psi=\overline{\iota_{+}^{*} \iota_{-}}$.

Analogously, we can show $\Psi^{*} \iota_{+}^{-1}=\iota_{-}^{*}$ and $\Psi^{*}=\overline{\iota_{-}^{*} \iota_{+}}$.
Theorem 5.17. Let $\mathcal{X}_{+}$and $\mathcal{X}_{0}$ be Hilbert spaces and $\iota_{+}: \operatorname{dom} \iota_{+} \subseteq \mathcal{X}_{+} \rightarrow \mathcal{X}_{0}$ be a densely defined, closed, and injective linear mapping with dense range. Then there exists a Hilbert space $\mathcal{X}_{-}$and a mapping $\iota_{-}$such that $\left(\mathcal{X}_{+}, \mathcal{X}_{0}, \mathcal{X}_{-}\right)$is a quasi Gelfand triple.

In particular, $\mathcal{X}_{-}$is given by Definition 5.2, where $D_{+}=\operatorname{ran} \iota_{+}$.
Proof. We will identify dom $\iota_{+}$with $\operatorname{ran} \iota_{+}$and denote it by $D_{+}$. Then item (iii) of Proposition 5.8 is satisfied. Hence, the corresponding $D_{-}$(Definition 5.2) is
dense in $\mathcal{X}_{0}$ and its completion $\mathcal{X}_{-}$(w.r.t. to $\|.\|_{\mathcal{X}_{-}}$) can be identified with $\mathcal{X}_{+}^{\prime}$ by Corollary 5.11. The linear mapping

$$
\iota_{-}:\left\{\begin{array}{rll}
D_{-} \subseteq \mathcal{X}_{-} & \rightarrow \mathcal{X}_{0} \\
g & \mapsto
\end{array}\right.
$$

is densely defined and injective by construction. By the already shown $\operatorname{ran} \iota_{-}=D_{-}$ is dense in $\mathcal{X}_{0}$. Finally, by Lemma $5.5 \iota_{-}$is closed and by Lemma $5.6 \operatorname{dom} \iota_{+}^{*}=$ $D_{-}=\operatorname{ran} \iota_{-}$.

Remark 5.18. By Theorem 5.17 the setting in the beginning of the section establishes a quasi Gelfand triple, if one assertion of Proposition 5.8 is satisfied.

Until the end of this section we will assume that ( $\left.\mathcal{X}_{+}, \mathcal{X}_{0}, \mathcal{X}_{-}\right)$is a quasi Gelfand triple and we will identify $\operatorname{dom} \iota_{+}$with $\operatorname{ran} \iota_{+}$and denote it by $D_{+}$. The set $D_{-}$ is defined by Definition 5.2 for $D_{+}$. This set coincides with ran $\iota_{-}$, which we will identify with dom $\iota_{-}$.

Proposition 5.19. The space $D_{+} \cap D_{-}$is complete with respect to $\|\cdot\|_{\mathcal{X}_{+} \cap \mathcal{X}_{-}}:=$ $\sqrt{\|\cdot\|_{\mathcal{X}_{+}}^{2}+\|\cdot\|_{\mathcal{X}_{-}}^{2}}$.
Proof. For $f \in D_{+} \cap D_{-}$we have

$$
\|f\|_{\mathcal{X}_{0}}^{2}=\left|\langle f, f\rangle_{\mathcal{X}_{0}}\right|=\left|\langle f, f\rangle_{\mathcal{X}_{-}, \mathcal{X}_{+}}\right| \leq\|f\|_{\mathcal{X}_{-}}\|f\|_{\mathcal{X}_{+}} \leq\|f\|_{\mathcal{X}_{+} \cap \mathcal{X}_{-}}^{2} .
$$

Hence, every Cauchy sequence in $D_{+} \cap D_{-}$with respect to $\|\cdot\|_{\mathcal{X}_{+} \cap \mathcal{X}_{-}}$is also a Cauchy sequence with respect to $\|\cdot\|_{\mathcal{X}_{0}},\|\cdot\|_{\mathcal{X}_{+}}$and $\|\cdot\|_{\mathcal{X}_{-}}$.

Let $\left(f_{n}\right)_{n \in \mathbb{N}}$ be a Cauchy sequence in $D_{+} \cap D_{-}$with respect to $\|\cdot\|_{\mathcal{X}_{+} \cap \mathcal{X}_{-}}$. By the closedness of $\iota_{+}$the limit with respect to $\|.\|_{\mathcal{X}_{0}}$ and the limit with respect to $\|\cdot\| \mathcal{X}_{+}$ coincide. The same argument for $\iota_{-}$yields that the limit with respect to $\|\cdot\|_{\mathcal{X}_{0}}$ and the limit with respect $\|\cdot\|_{\mathcal{X}_{-}}$also coincide. Therefore, all these limits have to coincide and $\left(f_{n}\right)_{n \in \mathbb{N}}$ converges to that limit in $D_{+} \cap D_{-}$w.r.t. $\|.\|_{\mathcal{X}_{+} \cap \mathcal{X}_{-}}$.

Lemma 5.20. The operator

$$
\left[\begin{array}{ll}
\iota_{+} & \iota_{-}
\end{array}\right]:\left\{\begin{array}{rll}
D_{+} \times D_{-} \subseteq \mathcal{X}_{+} \times \mathcal{X}_{-} & \rightarrow & \mathcal{X}_{0} \\
{\left[\begin{array}{l}
f \\
g
\end{array}\right]} & \mapsto & f+g
\end{array}\right.
$$

is closed.
Proof. Let $\left(\left(\left[\begin{array}{c}f_{n} \\ g_{n}\end{array}\right], z_{n}\right)\right)_{n \in \mathbb{N}}$ be a sequence in $\left[\iota_{+} \iota_{-}\right]$that converges to $\left(\left[\begin{array}{l}f \\ g\end{array}\right], z\right)$ in $\mathcal{X}_{+} \times \mathcal{X}_{-} \times \mathcal{X}_{0}$. Then we have

$$
\|z\|_{\mathcal{X}_{0}}^{2}=\lim _{n \rightarrow \infty}\left\|f_{n}+g_{n}\right\|_{\mathcal{X}_{0}}^{2}=\lim _{n \rightarrow \infty}\left(\left\|f_{n}\right\|_{\mathcal{X}_{0}}^{2}+\left\|g_{n}\right\|_{\mathcal{X}_{0}}^{2}+2 \operatorname{Re}\left\langle f_{n}, g_{n}\right\rangle \mathcal{X}_{0}\right)
$$

Since $2 \operatorname{Re}\left\langle f_{n}, g_{n}\right\rangle_{\mathcal{X}_{0}}$ converges to $2 \operatorname{Re}\langle f, g\rangle_{\mathcal{X}_{+}, \mathcal{X}_{-}}$, we conclude that $\left\|f_{n}\right\|_{\mathcal{X}_{0}}$ and $\left\|g_{n}\right\|_{\mathcal{X}_{0}}$ are bounded. Hence, there exists a subsequence of $\left(f_{n}\right)_{n \in \mathbb{N}}$ that converges weakly to an $\tilde{f} \in \mathcal{X}_{0}$. Moreover, by Lemma A. 3 we can pass on to a further subsequence $\left(f_{n(k)}\right)_{k \in \mathbb{N}}$ such that $\left(\frac{1}{j} \sum_{k=1}^{j} f_{n(k)}\right)_{j \in \mathbb{N}}$ converges to $\tilde{f}$ strongly (w.r.t. $\|.\|_{\mathcal{X}_{0}}$ ). The sequence $\left(\frac{1}{j} \sum_{k=1}^{j} f_{n(k)}\right)_{j \in \mathbb{N}}$ has still the limit $f$ in $\mathcal{X}_{+}$(w.r.t. $\|.\|_{\mathcal{X}_{+}}$) and because $\iota_{+}$is closed we conclude that $f=\tilde{f} \in D_{+}$. By linearity we also have $\frac{1}{j} \sum_{k=1}^{j} g_{n(k)} \rightarrow z-f$ in $\mathcal{X}_{0}$ for the same subsequence. Since $\frac{1}{j} \sum_{k=1}^{j} g_{n(k)}$ is a Cauchy sequence in both $\mathcal{X}_{-}$and $\mathcal{X}_{0}$, the closedness of $\iota_{-}$gives that $g=z-f \in D_{-}$. Hence, $z=\left[\begin{array}{ll}\iota_{+} & \iota_{-}\end{array}\right]\left[\begin{array}{c}f \\ g\end{array}\right]$ and the operator $\left[\iota_{+} \iota_{-}\right]$is closed.

Proposition 5.21. $D_{+} \cap D_{-}$is dense in $\mathcal{X}_{0}$ with respect to $\|\cdot\|_{\mathcal{X}_{0}}$.
Proof. By dom $\iota_{ \pm}^{*}=D_{\mp}$ (Lemma 5.14) we have

$$
\left.\mathcal{X}_{0}=\left(\operatorname{mul}\left[\iota_{+} \quad \iota_{-}\right]\right)^{\perp}=\overline{\operatorname{dom}\left[\iota_{+}\right.} \quad \iota_{-}\right]^{*}=\overline{\operatorname{dom} \iota_{+}^{*} \cap \operatorname{dom} \iota_{-}^{*}}=\overline{D_{-} \cap D_{+}} .
$$

The following theorem can be found in [17, Theorem 2 p. 200], we just changed that the operator maps into a different space, which does not change the proof.

Theorem 5.22 (J. von Neumann). Let $T$ be a closed linear operator from the Hilbert spaces $X$ to the Hilbert space $Y$. Then $T^{*} T$ and $T T^{*}$ are self-adjoint, and $\left(\mathrm{I}_{X}+T^{*} T\right)$ and $\left(\mathrm{I}_{Y}+T T^{*}\right)$ are boundedly invertible.

Corollary 5.23. The set $D_{+} \cap D_{-}$is dense in $\mathcal{X}_{+}$and $\mathcal{X}_{-}$with respect to their corresponding norms.

Proof. Applying Theorem 5.22 to $\iota_{+}$yields $\iota_{+}^{*} \iota_{+}$is self-adjoint. Hence, dom $\iota_{+}^{*} \iota_{+}$is dense in $\mathcal{X}_{+}$. By Lemma 5.14 dom $\iota_{+}^{*}=D_{-}$, consequently dom $\iota_{+}^{*} \iota_{+}=D_{+} \cap D_{-}$.

An analogous argument for $\iota_{-}$yields $D_{+} \cap D_{-}$is dense in $\mathcal{X}_{-}$.
Corollary 5.24. $D_{+}+D_{-}=\mathcal{X}_{0}$.
Proof. Applying Theorem 5.22 to $\iota_{+}$gives that ( $\mathrm{I}_{\mathcal{X}_{0}}+\iota_{+} \iota_{+}^{*}$ ) is onto. Hence, for every $x \in \mathcal{X}_{0}$ there exists a $g_{x} \in \operatorname{dom} \iota_{+} \iota_{+}^{*} \subseteq D_{-}$such that

$$
x=\underbrace{g_{x}}_{\in D_{-}}+\underbrace{\iota_{+} \iota_{+}^{*} g_{x}}_{\in D_{+}} .
$$

Since $g_{x} \in \operatorname{dom} \iota_{+} \iota_{+}^{*}$, we have $\iota_{+}^{*} g_{x} \in D_{+}$and consequently $x \in D_{+}+D_{-}$.
Proposition 5.25. Let $T$ be a bounded and boundedly invertible mapping from $\mathcal{X}_{0}$ to another Hilbert space $\mathcal{Y}_{0}$. Then $P_{+}:=T D_{+}$equipped with $\|f\|_{\mathcal{Y}_{+}}:=\left\|T^{-1} f\right\|_{\mathcal{X}_{+}}$ establishes a quasi Gelfand triple $\left(\mathcal{Y}_{+}, \mathcal{Y}_{0}, \mathcal{Y}_{-}\right)$, where $\mathcal{Y}_{+}$is the completion of $P_{+}$ and $\mathcal{Y}_{-}$is the completion of $P_{-}$defined as in Definition 5.2, where $D_{+}$is replaced by $P_{+}$. Moreover, $P_{-}=\left(T^{*}\right)^{-1} D_{-}$and $\|g\|_{\mathcal{Y}_{-}}=\left\|T^{*} g\right\|_{\mathcal{X}_{-}}$for $g \in P_{-}$.

Proof. The mapping $\left.T\right|_{D_{+}}: D_{+} \rightarrow P_{+}$is isometric and surjective, if we equip its domain with $\|\cdot\|_{\mathcal{X}_{+}}$and its codomain with $\|\cdot\|_{\mathcal{Y}_{+}}$. So the linear (single-valued) relation $\left[\begin{array}{cc}T & 0 \\ 0 & T\end{array}\right] \iota_{+}=\left\{(T f, T g):(f, g) \in \iota_{+}\right\} \subseteq \mathcal{Y}_{+} \times \mathcal{Y}_{0}$ is closed. Since this linear relation coincides with the embedding $\iota_{P_{+}}: P_{+} \subseteq \mathcal{Y}_{+} \rightarrow \mathcal{Y}_{0}, f \mapsto f$, Theorem 5.17 yields that $\left(\mathcal{Y}_{+}, \mathcal{Y}_{0}, \mathcal{Y}_{-}\right)$is a quasi Gelfand triple.

For $g \in P_{-}$we have

$$
\begin{aligned}
\|g\| \mathcal{Y}_{-}= & \sup _{h \in P_{+} \backslash\{0\}} \frac{\left|\langle g, h\rangle_{\mathcal{Y}_{0}}\right|}{\|h\|_{\mathcal{Y}_{+}}}=\sup _{f \in D_{+} \backslash\{0\}} \frac{\left|\langle g, T f\rangle_{\mathcal{Y}_{0}}\right|}{\|T f\|_{\mathcal{Y}_{+}}} \\
& =\sup _{f \in D_{+} \backslash\{0\}} \frac{\left|\left\langle T^{*} g, f\right\rangle_{\mathcal{X}_{0}}\right|}{\|f\|_{\mathcal{X}_{+}}}=\left\|T^{*} g\right\|_{\mathcal{X}_{-}}
\end{aligned}
$$

and consequently $P_{-}=\left(T^{*}\right)^{-1} D_{-}$.
Corollary 5.26. With the assumption from Proposition 5.25 the operators $\left.T\right|_{D_{+}}$ and $\left.\left(T^{*}\right)^{-1}\right|_{D_{-}}$can be continuously extended to unitary operators from $\mathcal{X}_{+}$and $\mathcal{X}_{-}$ to $\mathcal{Y}_{+}$and $\mathcal{Y}_{-}$respectively. These extension will be denoted by $T_{+}$and $\left(T^{*}\right)_{-}^{-1}$. Moreover, $\left\langle\left(T^{*}\right)_{-}^{-1} g, T_{+} f\right\rangle_{\mathcal{Y}_{-}, \mathcal{Y}_{+}}=\langle g, f\rangle_{\mathcal{X}_{-}, \mathcal{X}_{+}}$for $g \in \mathcal{X}_{-}$and $f \in \mathcal{X}_{+}$.

Corollary 5.27. Let $S, T$ be a bounded and boundedly invertible mappings on $\mathcal{X}_{0}$. Then $\left[\left.\left.S T\right|_{D_{+}} S\left(T^{*}\right)^{-1}\right|_{D_{-}}\right]$is a densely defined closed surjective linear operator from $\mathcal{X}_{+} \times \mathcal{X}_{-}$to $\mathcal{X}_{0}$. In particular $\operatorname{ran}\left[\left.\left.S T\right|_{D_{+}} S\left(T^{*}\right)^{-1}\right|_{D_{-}}\right]=\mathcal{X}_{0}$.
Proof. Let $P_{+}=T D_{+}$. Then by Proposition 5.25 the corresponding $P_{-}$can be obtained by $\left(T^{*}\right)^{-1} D_{-}$. The mapping

$$
\Xi:\left\{\begin{aligned}
\mathcal{X}_{+} \times \mathcal{X}_{-} \times \mathcal{X}_{0} & \rightarrow \mathcal{Y}_{+} \times \mathcal{Y}_{-} \times \mathcal{X}_{0}, \\
{\left[\begin{array}{l}
f \\
g \\
z
\end{array}\right] } & \mapsto\left[\begin{array}{ccc}
T_{+} & 0 & 0 \\
0 & \left(T^{*}\right)_{-}^{-1} & 0 \\
0 & 0 & S^{-1}
\end{array}\right]\left[\begin{array}{l}
f \\
g \\
z
\end{array}\right]
\end{aligned}\right.
$$

is linear bounded and boundedly invertible, where $\mathcal{Y}_{ \pm}$is the completion of $P_{ \pm}$as in Proposition 5.25. Since $\left(\mathcal{Y}_{+}, \mathcal{X}_{0}, \mathcal{Y}_{-}\right)$is a quasi Gelfand triple,

$$
\left[\begin{array}{ll}
\iota_{P_{+}} & \iota_{P_{-}}
\end{array}\right]=\left\{\left[\begin{array}{c}
T f \\
\left(T^{*}\right)^{-1} g \\
T f+\left(T^{*}\right)^{-1} g
\end{array}\right]: f \in D_{+}, g \in D_{-}\right\}
$$

is closed in $\mathcal{Y}_{+} \times \mathcal{Y}_{-} \times \mathcal{X}_{0}$ (Lemma 5.20) and therefore also its pre-image under $\Xi$

$$
\Xi^{-1}\left(\left[\begin{array}{ll}
\iota_{P_{+}} & \iota_{P_{-}}
\end{array}\right]\right)=\left[\begin{array}{ccc}
T^{-1} & 0 & 0 \\
0 & T^{*} & 0 \\
0 & 0 & S
\end{array}\right]\left[\begin{array}{ll}
\iota_{P_{+}} & \iota_{P_{-}}
\end{array}\right]=\left[\begin{array}{ll}
S T \iota_{+} & S\left(T^{*}\right)^{-1} \iota_{-}
\end{array}\right]
$$

is closed in $\mathcal{X}_{+} \times \mathcal{X}_{-} \times \mathcal{X}_{0}$. Furthermore, by Corollary 5.24

$$
\operatorname{ran}\left[\left.\left.S T\right|_{D_{+}} \quad S\left(T^{*}\right)^{-1}\right|_{D_{-}}\right]=S \operatorname{ran}\left[\begin{array}{ll}
\iota_{P_{+}} & \iota_{P-}
\end{array}\right]=S \mathcal{X}_{0}=\mathcal{X}_{0} .
$$

Lemma 5.28. Let $A_{0}$ be a densely defined, closed, skew-symmetric operator on $\mathcal{X}_{0}$, $\mathcal{Y}_{0}$ be a Hilbert space, and let $T: \mathcal{X}_{0} \rightarrow \mathcal{Y}_{0}$ be a bounded and boundedly invertible. Let ( $\mathcal{X}_{+}, \mathcal{X}_{0}, \mathcal{X}_{-}$) be a quasi Gelfand triple such that $\left(\mathcal{X}_{+}, B_{1}, \Psi B_{2}\right)$ is a boundary triple for $A_{0}^{*}$. Furthermore, let $\mathcal{Y}_{+}$and $\mathcal{Y}_{-}$be as defined in Proposition 5.25. Then $\left(\mathcal{Y}_{+}, \mathcal{Y}_{0}, \mathcal{Y}_{-}\right)$is also a quasi Gelfand triple such that $\left(\mathcal{Y}_{+}, T_{+} B_{1}, \Phi\left(T^{*}\right)_{-}^{-1} B_{2}\right)$ is a boundary triple for $A_{0}^{*}$, where $\Phi$ denotes the duality map of $\left(\mathcal{Y}_{+}, \mathcal{Y}_{0}, \mathcal{Y}_{-}\right)$.

Proof. By Proposition $5.25\left(\mathcal{Y}_{+}, \mathcal{Y}_{0}, \mathcal{Y}_{-}\right)$is a quasi Gelfand triple. For $x, y \in \operatorname{dom} A_{0}^{*}$ we have, by Corollary 5.26,

$$
\begin{aligned}
\left\langle B_{1} x, \Psi B_{2} y\right\rangle_{\mathcal{X}_{+}} & =\left\langle B_{1} x, B_{2} y\right\rangle_{\mathcal{X}_{+}, \mathcal{X}_{-}}=\left\langle T_{+} B_{1} x,\left(T^{*}\right)_{-}^{-1} B_{2} y\right\rangle_{\mathcal{Y}_{+}, \mathcal{Y}_{-}} \\
& =\left\langle T_{+} B_{1} x, \Phi\left(T^{*}\right)_{-}^{-1} B_{2} y\right\rangle_{\mathcal{Y}_{+}} .
\end{aligned}
$$

Since $T_{+}: \mathcal{X}_{+} \rightarrow \mathcal{Y}_{+}$and $\left(T^{*}\right)_{-}^{-1}: \mathcal{X}_{-} \rightarrow \mathcal{Y}_{-}$are surjective, the surjectivity of $\left[\begin{array}{c}T_{+} B_{1} \\ \Phi\left(T^{*}\right)_{-}^{1} B_{2}\end{array}\right]=\left[\begin{array}{cc}T_{+} & 0 \\ 0 & \Phi\left(T^{*}\right)_{-}^{-1} \Psi^{-1}\end{array}\right]\left[\begin{array}{c}B_{1} \\ \Psi B_{2}\end{array}\right]$ follows from the surjectivity of $\left[\begin{array}{c}B_{1} \\ \Psi B_{2}\end{array}\right]$.

Remark 5.29. In the setting of Lemma 5.28 the duality map $\Phi$ can be described by $\Phi=T_{+} \Psi\left(\left(T^{*}\right)_{-}^{-1}\right)^{-1}$. Note that $\left(\left(T^{*}\right)_{-}^{-1}\right)^{-1}$ can be described by the continuous extension of $\left.T^{*}\right|_{P_{-}}: P_{-} \subseteq \mathcal{Y}_{-} \rightarrow \mathcal{X}_{-}$. We denote this extension by $T_{-}^{*}$. Hence, $\Phi=T_{+} \Psi T_{-}^{*}$.
6. Boundary spaces. In this section we will construct a suitable boundary space $\mathcal{V}_{L}$ (Definition 6.5), such that we can extend the integration by parts formula (Lemma 3.8). We will formulate the boundary conditions in this space in section 7. This space will provide a quasi Gelfand triple with a subspace of $L^{2}(\partial \Omega)$ as pivot space. In order to impose different boundary conditions on different parts of the boundary we introduce boundary operators that only act on a part of the boundary and their boundary spaces $\mathcal{V}_{L, \Gamma_{1}}$.

Definition 6.1. We say $\left(\Gamma_{j}\right)_{j=1}^{k}$, where $\Gamma_{j} \subseteq \partial \Omega$, is a splitting with thin boundaries of $\partial \Omega$, if
(i) $\bigcup_{j=1}^{k} \overline{\Gamma_{j}}=\partial \Omega$,
(ii) the sets $\Gamma_{j}$ are pairwise disjoint,
(iii) the sets $\Gamma_{j}$ are relatively open in $\partial \Omega$,
(iv) the boundaries of $\Gamma_{j}$ have zero measure w.r.t. the surface measure of $\partial \Omega$.

For $\Gamma \subseteq \partial \Omega$ we will denote by $P_{\Gamma}$ the orthogonal projection from $L^{2}(\partial \Omega)^{m_{1}}$ on $L_{\pi}^{2}(\Gamma):=\overline{\operatorname{ran} \mathbb{1}_{\Gamma} L_{\nu}} \subseteq L^{2}(\Gamma)^{m_{1}}$, where $\mathbb{1}_{M}$ denotes the indicator function for a set $M$. We endow $L_{\pi}^{2}(\Gamma)$ with the inner product of $L^{2}(\partial \Omega)^{m_{1}}$. Therefore, we can adapt (3.1) to obtain

$$
\begin{equation*}
\left\langle L_{\partial} f, g\right\rangle_{L^{2}(\Omega)^{m_{1}}}+\left\langle f, L_{\partial}^{\mathrm{H}} g\right\rangle_{L^{2}(\Omega)^{m_{2}}}=\langle L_{\nu} \gamma_{0} f, \underbrace{P_{\partial \Omega} \gamma_{0} g}_{\pi_{L} g}\rangle_{L^{2}(\partial \Omega)^{m_{1}}} \tag{6.1}
\end{equation*}
$$

We define $\pi_{L}^{\Gamma}: H^{1}(\Omega)^{m_{1}} \rightarrow L_{\pi}^{2}(\Gamma)$ by $\pi_{L}^{\Gamma}:=P_{\Gamma} \gamma_{0}$ and $\pi_{L}:=\pi_{L}^{\partial \Omega}$. Since both $P_{\Gamma}$ and $\gamma_{0}$ are continuous, the mapping $\pi_{L}^{\Gamma}$ is also continuous. Therefore, $\operatorname{ker} \pi_{L}^{\Gamma}$ is closed. Note that $P_{\Gamma}=\mathbb{1}_{\Gamma} P_{\partial \Omega}$ and consequently $\pi_{L}^{\Gamma}=\mathbb{1}_{\Gamma} \pi_{L}$, and $\mathbb{1}_{\Gamma} L_{\nu}=L_{\nu} \mathbb{1}_{\Gamma}$.

Example 6.2. Let $L$ be as in Example 3.3. Then $L_{\nu} f=\nu \cdot f$ and $L_{\nu}$ is certainly surjective. Therefore, $L_{\pi}^{2}(\partial \Omega)=L^{2}(\partial \Omega), \pi_{L}=\gamma_{0}$ and $\pi_{L}^{\Gamma}=\mathbb{1}_{\Gamma} \gamma_{0}$. Since $L_{\partial}^{\mathrm{H}}=\operatorname{grad}$, we have $H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)=H^{1}(\Omega)$.
Lemma 6.3. Let $\Gamma \subseteq \partial \Omega$ be relatively open and let the boundary of $\Gamma$ have zero measure (w.r.t. the surface measure of $\partial \Omega$ ). Then $\operatorname{ker} \pi_{L}^{\Gamma}$ is closed as subspace of $H^{1}(\Omega)^{m_{1}}$ endowed with the trace topology of $\|\cdot\|_{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)}$, i.e.

$$
{\overline{\operatorname{ker} \pi_{L}^{\Gamma}}}^{\|\cdot\|_{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)} \cap H^{1}(\Omega)^{m_{1}}=\operatorname{ker} \pi_{L}^{\Gamma} . . . .}
$$

Proof. Clearly, $\overline{\operatorname{ker} \pi_{L}^{\Gamma}}{ }^{\|\cdot\|_{H\left(L_{\partial}^{H}, \Omega\right)} \cap H^{1}(\Omega)^{m_{1}} \supseteq \operatorname{ker} \pi_{L}^{\Gamma} \text {. So we will show the other }}$ inclusion. Note that for $\Upsilon \subseteq \partial \Omega$ we have

$$
H_{\Upsilon}^{1}(\Omega)^{m_{2}}:=\left\{f \in H^{1}(\Omega)^{m_{2}}: \mathbb{1}_{\Upsilon} \gamma_{0} f=0 \in L^{2}(\partial \Omega)^{m_{2}}\right\} .
$$

Hence, $H_{\partial \Omega \backslash \Gamma}^{1}(\Omega)^{m_{2}}=H_{\partial \Omega \backslash \bar{\Gamma}}^{1}(\Omega)^{m_{2}}$, since the boundary of $\Gamma$ has zero measure. Let $\left(g_{n}\right)_{n \in \mathbb{N}}$ be a sequence in ker $\pi_{L}^{\Gamma}$ which converges to $g \in H^{1}(\Omega)^{m_{1}}$ with respect to $\|\cdot\|_{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)}$. By Corollary 3.9 we have for an arbitrary $f \in H_{\partial \Omega \backslash \Gamma}^{1}(\Omega)^{m_{2}}$

$$
\left|\left\langle L_{\nu} \gamma_{0} f, \pi_{L}^{\Gamma}\left(g-g_{n}\right)\right\rangle_{L^{2}}\right|=\left|\left\langle L_{\nu} \gamma_{0} f, \pi_{L}\left(g-g_{n}\right)\right\rangle_{L^{2}}\right| \leq\|f\|_{H\left(L_{\partial}, \Omega\right)}\left\|g-g_{n}\right\|_{H\left(L_{\partial}^{H}, \Omega\right)}
$$

Since $\pi_{L}^{\Gamma}\left(g-g_{n}\right)=\pi_{L}^{\Gamma} g$ and the right-hand-side converges to 0 , we can see that $\pi_{L}^{\Gamma} g \perp L_{\nu} \gamma_{0} H_{\partial \Omega \backslash \Gamma}^{1}(\Omega)^{m_{2}}$. By [13, Th. 13.6.10, Re. 13.6.12] $\gamma_{0} H_{\partial \Omega \backslash \Gamma}^{1}(\Omega)^{m_{2}}$ is dense in $L^{2}(\Gamma)^{m_{2}}$, which implies $\pi_{L}^{\Gamma} g \perp \operatorname{ran} \mathbb{1}_{\Gamma} L_{\nu}$. By definition $\pi_{L}^{\Gamma} g$ is also in $\overline{\operatorname{ran} \mathbb{1}_{\Gamma} L_{\nu}}$, which leads to $\pi_{L}^{\Gamma} g=0$. Hence, ker $\pi_{L}^{\Gamma}$ is closed in $H^{1}(\Omega)^{m_{1}}$ with respect to $\|\cdot\|_{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)}$.

By the previous lemma

$$
\|\phi\|_{M_{\Gamma}}:=\inf \left\{\|g\|_{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)}: \pi_{L}^{\Gamma} g=\phi\right\}
$$

is a norm on $M_{\Gamma}:=\operatorname{ran} \pi_{L}^{\Gamma}$. The next lemma will show that this norm is induced by an inner product.
Lemma 6.4. Let $\Gamma \subseteq \partial \Omega$ be relatively open and let the boundary of $\Gamma$ have zero measure (w.r.t. the surface measure of $\partial \Omega$ ). Then the space $\left(M_{\Gamma},\|\cdot\|_{M_{\Gamma}}\right)$ is a preHilbert space. Furthermore, its completion denoted by $\left(\overline{M_{\Gamma}},\|\cdot\|_{\overline{M_{\Gamma}}}\right)$ is isomorphic to the Hilbert space $H\left(L_{\partial}^{\mathrm{H}}, \Omega\right) /{\overline{\operatorname{ker} \pi_{L}^{\Gamma}}}^{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)}$. The mapping $\pi_{L}^{\Gamma}: H^{1}(\Omega)^{m_{1}} \rightarrow M_{\Gamma}$ can be continuously extended to a surjective contraction $\bar{\pi}_{L}^{\Gamma}: H\left(L_{\partial}^{\mathrm{H}}, \Omega\right) \rightarrow \overline{M_{\Gamma}}$. The kernel of $\bar{\pi}_{L}^{\Gamma}$ satisfies $\operatorname{ker} \bar{\pi}_{L}^{\Gamma}={\overline{\operatorname{ker} \pi_{L}^{\Gamma}}}^{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)}$.

Instead of $\bar{\pi}_{L}^{\partial \Omega}$ we will just write $\bar{\pi}_{L}$.
Proof. By Lemma $6.3 \operatorname{ker} \pi_{L}^{\Gamma}$ is closed in $H^{1}(\Omega)^{m_{1}}$ with respect to trace topology of $\|\cdot\|_{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)}$, which implies that $\left(H^{1}(\Omega)^{m_{1}} / \operatorname{ker} \pi_{L}^{\Gamma},\|\cdot\|_{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)} / \operatorname{ker} \pi_{L}^{\Gamma}\right)$ is a normed space (normed space factorized by a closed subspace is again a normed space). Since

$$
\left\|[g]_{\sim}\right\|_{H\left(L_{\partial}^{H}, \Omega\right)}{\operatorname{ker} \pi_{L}^{\Gamma}}=\left\|\pi_{L}^{\Gamma} g\right\|_{M_{\Gamma}},
$$

it is straight forward that $[g]_{\sim} \mapsto \pi_{L}^{\Gamma} g$ is an isometry from $\left(H^{1}(\Omega)^{m_{1}} / \operatorname{ker} \pi_{L}^{\Gamma}\right.$, $\left.\|\cdot\|_{H\left(L_{\partial}^{H}, \Omega\right)}{\operatorname{ker} \pi_{L}^{\Gamma}}\right)$ onto $\left(M_{\Gamma},\|\cdot\|_{M_{\Gamma}}\right)$.

Clearly, $\left(M_{\Gamma},\|\cdot\|_{M_{\Gamma}}\right)$ has a completion $\left(\overline{M_{\Gamma}},\|\cdot\|_{\overline{M_{\Gamma}}}\right)$. By definition of the norm $\|\cdot\|_{M_{\Gamma}}$ we have for every $g \in H^{1}(\Omega)^{m_{1}}$

$$
\left\|\pi_{L}^{\Gamma} g\right\|_{\overline{M_{\Gamma}}}=\left\|\pi_{L}^{\Gamma} g\right\|_{M_{\Gamma}} \leq\|g\|_{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)} .
$$

Therefore, we can extend $\pi_{L}^{\Gamma}$ by continuity on $H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$. This extension is denoted by $\bar{\pi}_{L}^{\Gamma}$ and is a contraction by the previous equation.

Let $g \in H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$. Then by Theorem 3.18 there exists a sequence $\left(g_{n}\right)_{n \in \mathbb{N}}$ in $H^{1}(\Omega)^{m_{1}}$, which converges to $g$. Therefore, we have

$$
\left\|\bar{\pi}_{L}^{\Gamma} g\right\|_{\overline{M_{\Gamma}}}=\lim _{n \rightarrow \infty}\left\|\pi_{L}^{\Gamma} g_{n}\right\|_{M_{\Gamma}}=\lim _{n \rightarrow \infty} \inf _{k \in \operatorname{ker} \pi_{L}^{\Gamma}}\left\|g_{n}+k\right\|_{H\left(L_{\partial}^{H}, \Omega\right)}
$$

The triangular inequality yields

$$
\inf _{k \in \operatorname{ker} \pi_{L}^{\Gamma}}\|g+k\|-\left\|g_{n}-g\right\| \leq \inf _{k \in \operatorname{ker} \pi_{L}^{\Gamma}}\left\|g_{n}+k\right\| \leq \inf _{k \in \operatorname{ker} \pi_{L}^{\Gamma}}\|g+k\|+\left\|g_{n}-g\right\| .
$$

Hence, we have

$$
\begin{equation*}
\left\|\bar{\pi}_{L}^{\Gamma} g\right\|_{\overline{M_{\Gamma}}}=\inf _{k \in \operatorname{ker} \pi_{L}^{\Gamma}}\|g+k\|_{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)}=\inf _{k \in \overline{\operatorname{ker} \pi_{L}^{\Gamma}}}\|g+k\|_{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)} \tag{6.2}
\end{equation*}
$$

and consequently $H\left(L_{\partial}^{\mathrm{H}}, \Omega\right) / \overline{\operatorname{ker} \pi_{L}^{\Gamma}}$ is isomorphic to ran $\bar{\pi}_{L}^{\Gamma}$. Since $H\left(L_{\partial}^{\mathrm{H}}, \Omega\right) / \overline{\operatorname{ker} \pi_{L}^{\Gamma}}$ is a Hilbert space, in particular complete, and $M_{\Gamma} \subseteq \operatorname{ran} \bar{\pi}_{L}^{\Gamma} \subseteq \overline{M_{\Gamma}}$, we have $\overline{M_{\Gamma}}=$ $\operatorname{ran} \bar{\pi}_{L}^{\Gamma}$. This makes $\overline{M_{\Gamma}}$ also a Hilbert space and $M_{\Gamma}$ a pre-Hilbert space.

Finally, equation (6.2) implies $\operatorname{ker} \bar{\pi}_{L}^{\Gamma}=\overline{\operatorname{ker} \pi_{L}^{\Gamma}}$.

Now we are able to define a complete subspace of $H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$ that is in some sense 0 at one part of the boundary and the corresponding boundary space for the other part of the boundary.

Definition 6.5. Let $\Gamma_{0}, \Gamma_{1} \subseteq \partial \Omega$ be a splitting with thin boundaries and $\bar{\pi}_{L}$ the extension of $\pi_{L}$ introduced in Lemma 6.4. Then we define

$$
H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right):=\operatorname{ker} \bar{\pi}_{L}^{\Gamma_{0}} \quad \text { and } \quad \mathcal{V}_{L, \Gamma_{1}}:=\left.\operatorname{ran} \bar{\pi}_{L}\right|_{H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)},
$$

where we endow $H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$ with $\|\cdot\|_{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)}$ and $\mathcal{V}_{L, \Gamma_{1}}$ with $\|\cdot\|_{\mathcal{V}_{L, \Gamma_{1}}}:=\|\cdot\|_{\overline{M_{\partial \Omega}}}$. Instead of $\mathcal{V}_{L, \partial \Omega}=\operatorname{ran} \bar{\pi}_{L}=\overline{M_{\partial \Omega}}$ we just write $\mathcal{V}_{L}$.

From now on until the end of this section we will assume that $\Gamma_{0}, \Gamma_{1} \subseteq \partial \Omega$ is a splitting with thin boundaries. By Lemma $6.4 \mathcal{V}_{L}$ is a Hilbert space.

Note that $\mathcal{V}_{L, \Gamma_{1}}$ and $\overline{M_{\Gamma_{1}}}$ are not necessarily the same space. Although, we have $\bar{\pi}_{L}^{\Gamma_{1}} g=\bar{\pi}_{L} g$ (in $L^{2}(\partial \Omega)^{m_{1}}$ ) for $g \in H^{1}(\Omega)^{m_{1}} \cap H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$, but we can only say $\left\|\bar{\pi}_{L}^{\Gamma_{1}} g\right\|_{M_{\Gamma_{1}}} \leq\left\|\bar{\pi}_{L} g\right\|_{\mathcal{V}_{L, \Gamma_{1}}}$.

Example 6.6. Continuing Example 6.2 yields $H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)=H_{\Gamma_{0}}^{1}(\Omega)^{m_{1}}=\{f \in$ $\left.H^{1}(\Omega)^{m_{1}}: \mathbb{1}_{\Gamma_{1}} \gamma_{0} f=0\right\}$ which already appeared in the proof of Lemma 6.3. Moreover, we have $\bar{\pi}_{L}=\gamma_{0}, \bar{\pi}_{L}^{\Gamma_{1}}=\mathbb{1}_{\Gamma_{1}} \gamma_{0}, \mathcal{V}_{L}=H^{1 / 2}(\partial \Omega)$, and $\mathcal{V}_{L, \Gamma_{1}}=\left\{f \in H^{1 / 2}(\partial \Omega)\right.$ : $\left.\left.f\right|_{\Gamma_{0}}=0\right\}$.

Lemma 6.7. The space $H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$ equipped with $\langle., .\rangle_{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)}$ is a Hilbert space and $H^{1}(\Omega)^{m_{1}} \cap H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$ is dense in $H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$. Moreover, $\mathcal{V}_{L, \Gamma_{1}}$ is a closed subspace of $\mathcal{V}_{L}$ and therefore also a Hilbert space.

Proof. By definition of $H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$ and Lemma 6.4 we have

$$
H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)=\operatorname{ker} \bar{\pi}_{L}^{\Gamma_{0}}=\overline{\operatorname{ker} \pi_{L}^{\Gamma_{0}}}=\overline{H^{1}(\Omega)^{m_{1}} \cap H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)} .
$$

Note that ker $\pi_{L} \subseteq \operatorname{ker} \pi_{L}^{\Gamma_{0}}$, since $\pi_{L}^{\Gamma_{0}}=\mathbb{1}_{\Gamma_{0}} \pi_{L}$. Again by Lemma 6.4, we have

$$
\operatorname{ker} \bar{\pi}_{L}=\overline{\operatorname{ker} \pi_{L}} \subseteq \overline{\operatorname{ker} \pi_{L}^{\Gamma_{0}}}=\operatorname{ker} \bar{\pi}_{L}^{\Gamma_{0}}
$$

Therefore, $\bar{\pi}_{L}^{\Gamma_{0}} \circ \bar{\pi}_{L}^{-1}: \mathcal{V}_{L} \rightarrow \overline{M_{\Gamma_{0}}}$ is single-valued (well-defined). For arbitrary $\phi \in \mathcal{V}_{L}$ and $g \in \bar{\pi}_{L}^{-1} \phi$ we have

$$
\left\|\bar{\pi}_{L}^{\Gamma_{0}} \circ \bar{\pi}_{L}^{-1} \phi\right\|_{\bar{M}_{\Gamma_{0}}}=\inf _{k \in \operatorname{ker} \bar{\pi}_{L}^{\Gamma_{0}}}\|g+k\|_{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)} \leq \inf _{k \in \operatorname{ker} \bar{\pi}_{L}}\|g+k\|_{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)}=\|\phi\|_{\mathcal{V}_{L}} .
$$

Hence, $\bar{\pi}_{L}^{\Gamma_{0}} \circ \bar{\pi}_{L}^{-1}$ is continuous and $\operatorname{ker} \bar{\pi}_{L}^{\Gamma_{0}} \circ \bar{\pi}_{L}^{-1}$ is closed in $\mathcal{V}_{L}$ and therefore also a Hilbert space endowed with $\langle., .\rangle_{\mathcal{V}_{L}}$. The equivalences

$$
\phi \in \operatorname{ker} \bar{\pi}_{L}^{\Gamma_{0}} \circ \bar{\pi}_{L}^{-1} \quad \Leftrightarrow \quad \bar{\pi}_{L}^{-1} \phi \subseteq \operatorname{ker} \bar{\pi}_{L}^{\Gamma_{0}} \quad \Leftrightarrow \quad \phi \in \underbrace{\left.\operatorname{ran} \bar{\pi}_{L}\right|_{\mathrm{ker}} \bar{\pi}_{L}^{\Gamma_{0}}}_{=\mathcal{V}_{L, \Gamma_{1}}}
$$

imply that $\mathcal{V}_{L, \Gamma_{1}}$ is closed and therefore a Hilbert space.
Proposition 6.8. The mapping $\mathbb{1}_{\Gamma_{1}} L_{\nu} \gamma_{0}: H^{1}(\Omega)^{m_{2}} \rightarrow L_{\pi}^{2}\left(\Gamma_{1}\right)$ can be extended to a linear continuous mapping

$$
\bar{L}_{\nu}^{\Gamma_{1}}: H\left(L_{\partial}, \Omega\right) \rightarrow \mathcal{V}_{L, \Gamma_{1}}^{\prime}
$$

such that $\left\|\bar{L}_{\nu}^{\Gamma_{1}} f\right\|_{\mathcal{V}_{L, \Gamma_{1}}^{\prime}} \leq\|f\|_{H\left(L_{\partial}, \Omega\right)}$.

Proof. Let $f \in H^{1}(\Omega)^{m_{2}}$. For $g \in H^{1}(\Omega)^{m_{1}} \cap H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$ we have by Corollary 3.9

$$
\left|\left\langle\mathbb{1}_{\Gamma_{1}} L_{\nu} \gamma_{0} f, \bar{\pi}_{L} g\right\rangle_{L^{2}\left(\Gamma_{1}\right)^{m_{1}}}\right|=\left|\left\langle L_{\nu} \gamma_{0} f, \bar{\pi}_{L} g\right\rangle_{L^{2}(\partial \Omega)^{m_{1}}}\right| \leq\|f\|_{H\left(L_{\partial}, \Omega\right)}\|g\|_{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)}
$$

By Lemma 6.7 the subspace $M:=\left.\operatorname{ran} \bar{\pi}_{L}\right|_{H^{1}(\Omega)^{m_{1} \cap H_{\Gamma_{0}}\left(L_{\partial}^{H}, \Omega\right)}} \subseteq L_{\pi}^{2}\left(\Gamma_{1}\right)^{m_{1}}$ of $\mathcal{V}_{L, \Gamma_{1}}$ is dense in $\mathcal{V}_{L, \Gamma_{1}}$. For $\phi \in M$ there exists at least one $g \in H^{1}(\Omega)^{m_{1}} \cap H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$ such that $\pi_{L} g=\phi$. Hence, we can rewrite the inequality as

$$
\begin{aligned}
\left|\left\langle\mathbb{1}_{\Gamma_{1}} L_{\nu} \gamma_{0} f, \phi\right\rangle_{L^{2}\left(\Gamma_{1}\right)^{m_{1}}}\right| & \leq\|f\|_{H\left(L_{\partial}, \Omega\right)} \inf _{g \in H^{1}(\Omega)^{m_{1} \cap H_{\Gamma_{0}}\left(L_{\partial}^{H}, \Omega\right)}}^{\bar{\pi}_{L} g=\phi} \\
& \|g\|_{H\left(L_{\partial}^{H}, \Omega\right)} \\
& \|f\|_{H\left(L_{\partial}, \Omega\right)}\|\phi\|_{\mathcal{V}_{L, \Gamma_{1}}} .
\end{aligned}
$$

We will extend the mapping $\phi \mapsto\left\langle\mathbb{1}_{\Gamma_{1}} L_{\nu} \gamma_{0} f, \phi\right\rangle_{L^{2}\left(\Gamma_{1}\right)^{m_{1}}}$ by continuity on $\mathcal{V}_{L, \Gamma_{1}}$. We will denote this extension by $\Xi_{f}$. Therefore, we have

$$
\left|\Xi_{f}(\phi)\right| \leq\|f\|_{H\left(L_{\partial}, \Omega\right)}\|\phi\|_{\mathcal{V}_{L, \Gamma_{1}}}
$$

This means that the mapping $f \mapsto \Xi_{f}$ from $H^{1}(\Omega)^{m_{2}}$ to $\mathcal{V}_{L, \Gamma_{1}}^{\prime}$ is continuous, if we endow $H^{1}(\Omega)^{m_{2}}$ with $\|\cdot\|_{H\left(L_{z}, \Omega\right)}$. Once again, we will extend this mapping by continuity on $H\left(L_{\partial}, \Omega\right)$ and denote it by $\bar{L}_{\nu}^{\Gamma_{1}}$.

Instead of writing $\bar{L}_{\nu}^{\partial \Omega}$ we will just write $\bar{L}_{\nu}$.
Remark 6.9. Since $\mathcal{V}_{L, \Gamma_{1}}$ is a subspace of $\mathcal{V}_{L, \partial \Omega}=\mathcal{V}_{L}$ every element of $\mathcal{V}_{L}^{\prime}$ can also be treated as an element of $\mathcal{V}_{L, \Gamma_{1}}^{\prime}$. By definition of $\bar{L}_{\nu}^{\Gamma_{1}}$ and $\bar{L}_{\nu}$ it is easy to see that $\bar{L}_{\nu}^{\Gamma_{1}} f=\left.\bar{L}_{\nu} f\right|_{\mathcal{V}_{L, \Gamma_{1}}}$ or equivalently $\bar{L}_{\nu}^{\Gamma_{1}} f$ and $\bar{L}_{\nu} f$ coincide as elements of $\mathcal{V}_{L, \Gamma_{1}}^{\prime}$ for $f \in H\left(L_{\partial}, \Omega\right)$. Hence, we can say $\left.\mathcal{V}_{L}^{\prime}\right|_{\mathcal{V}_{L, \Gamma_{1}}} \subseteq \mathcal{V}_{L, \Gamma_{1}}^{\prime}$. Since Hahn-Banach gives the reverse inclusion we can even say $\left.\mathcal{V}_{L}^{\prime}\right|_{\mathcal{V}_{L, \Gamma_{1}}}=\mathcal{V}_{L, \Gamma_{1}}^{\prime}$.

The reason for even defining $\bar{L}_{\nu}^{\Gamma_{1}}$ instead of just using $\bar{L}_{\nu}$ is that the range of its restriction to $H^{1}(\Omega)^{m_{2}}$ is also contained in $L_{\pi}^{2}\left(\Gamma_{1}\right)$, which will be important for getting a quasi Gelfand triple.

Corollary 6.10. For $f \in H\left(L_{\partial}, \Omega\right)$ and $g \in H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$ we have

$$
\left\langle L_{\partial} f, g\right\rangle_{L^{2}(\Omega)^{m_{1}}}+\left\langle f, L_{\partial}^{\mathrm{H}} g\right\rangle_{L^{2}(\Omega)^{m_{2}}}=\left\langle\bar{L}_{\nu} f, \bar{\pi}_{L} g\right\rangle_{\mathcal{V}_{L, \Gamma_{1}}^{\prime}, \mathcal{V}_{L, \Gamma_{1}}} .
$$

For $f \in H\left(L_{\partial}, \Omega\right)$ and $g \in H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$ we have

$$
\begin{aligned}
\left\langle L_{\partial} f, g\right\rangle_{L^{2}(\Omega)^{m_{1}}}+\left\langle f, L_{\partial}^{\mathrm{H}} g\right\rangle_{L^{2}(\Omega)^{m_{2}}} & =\left\langle\bar{L}_{\nu} f, \bar{\pi}_{L} g\right\rangle_{\mathcal{V}_{L}^{\prime}, \mathcal{V}_{L}} \\
& =\left\langle\bar{\pi}_{L^{\mathrm{H}}} f, \bar{L}_{\nu}^{\mathrm{H}} g\right\rangle_{\mathcal{V}_{L^{\mathrm{H}}}, \mathcal{V}_{L^{\mathrm{H}}}^{\prime}} .
\end{aligned}
$$

Proof. Since $H^{1}(\Omega)^{m_{2}}$ is dense in $H\left(L_{\partial}, \Omega\right)$ and $H^{1}(\Omega)^{m_{1}} \cap H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$ is dense in $H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$, the first equation follows from (6.1) by continuity. The second equation is just the special case $\Gamma_{0}=\emptyset$ and switching the roles of $L_{\partial}$ and $L_{\partial}^{H}$ yields the last equation.
Theorem 6.11. The mapping $\bar{L}_{\nu}: H\left(L_{\partial}, \Omega\right) \rightarrow \mathcal{V}_{L}^{\prime}$ is linear, bounded and onto.
Proof. By Proposition 6.8 we already know that $\bar{L}_{\nu}$ is linear and bounded from $H\left(L_{\partial}, \Omega\right)$ to $\mathcal{V}_{L}^{\prime}$.

Let $\mu \in \mathcal{V}_{L}^{\prime}$ be arbitrary. Since $\bar{\pi}_{L}$ is continuous from $H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$ to $\mathcal{V}_{L}$, the mapping $g \mapsto\left\langle\mu, \bar{\pi}_{L} g\right\rangle_{\mathcal{V}_{L}^{\prime}, \mathcal{V}_{L}}$ is continuous from $H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$ to $\mathbb{C}$. Consequently, there exists an $h \in H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$ such that

$$
\langle h, g\rangle_{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)}=\left\langle\mu, \pi_{L} g\right\rangle_{\mathcal{V}_{L}^{\prime}, \mathcal{V}_{L}} \quad \text { for all } \quad g \in H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)
$$

For a test function $v \in \mathcal{D}(\Omega)^{m_{1}}$ we have

$$
\begin{aligned}
0 & =\left\langle\mu, \pi_{L} v\right\rangle_{\mathcal{V}_{L}^{\prime}}, \mathcal{V}_{L}=\langle h, v\rangle_{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)}=\langle h, v\rangle_{L^{2}(\Omega)^{m_{1}}}+\left\langle L_{\partial}^{\mathrm{H}} h, L_{\partial}^{\mathrm{H}} v\right\rangle_{L^{2}(\Omega)^{m_{2}}} \\
& =\langle h, v\rangle_{\mathcal{D}^{\prime}(\Omega)^{m_{1}}, \mathcal{D}(\Omega)^{m_{1}}}+\left\langle L_{\partial}^{\mathrm{H}} h, L_{\partial}^{\mathrm{H}} v\right\rangle_{\mathcal{D}^{\prime}(\Omega)^{m_{2}}, \mathcal{D}(\Omega)^{m_{2}}} \\
& =\left\langle\left(\mathrm{I}-L_{\partial} L_{\partial}^{\mathrm{H}}\right) h, v\right\rangle_{\mathcal{D}^{\prime}(\Omega)^{m_{1}}, \mathcal{D}(\Omega)^{m_{1}}}
\end{aligned}
$$

This means $L_{\partial} L_{\partial}^{\mathrm{H}} h=h$ in the sense of distributions. However, $h \in H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$ implies $h \in L^{2}(\Omega)$, which in turn gives $L_{\partial} L_{\partial}^{\mathrm{H}} h \in L^{2}(\Omega)^{m_{1}}$, and $L_{\partial}^{\mathrm{H}} h \in L^{2}(\Omega)^{m_{2}}$. Therefore, $f:=L_{\partial} h \in H\left(L_{\partial}, \Omega\right)$. By Corollary 6.10 for $f=L_{\partial}^{\mathrm{H}} h \in H\left(L_{\partial}, \Omega\right)$ and $g \in H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$ we have

$$
\begin{aligned}
\left\langle\mu, \pi_{L} g\right\rangle_{\mathcal{V}_{L}^{\prime}, \mathcal{V}_{L}} & =\langle h, g\rangle_{H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)}=\langle h, g\rangle_{L^{2}(\Omega)^{m_{1}}}+\left\langle L_{\partial}^{\mathrm{H}} h, L_{\partial}^{\mathrm{H}} g\right\rangle_{L^{2}(\Omega)^{m_{2}}} \\
& =\left\langle\left(\mathrm{I}-L_{\partial} L_{\partial}^{\mathrm{H}}\right) h, g\right\rangle_{L^{2}(\Omega)^{m_{1}}}+\left\langle\bar{L}_{\nu} L_{\partial}^{\mathrm{H}} h, \pi_{L} g\right\rangle_{\mathcal{V}_{L}^{\prime}, \mathcal{V}_{L}} \\
& =\langle\bar{L}_{\nu} \underbrace{\left(L_{\partial}^{\mathrm{H}} h\right)}_{=f}, \pi_{L} g\rangle_{\mathcal{V}_{L}^{\prime}, \mathcal{V}_{L}} .
\end{aligned}
$$

Hence, $\bar{L}_{\nu} f=\mu$ and $\bar{L}_{\nu}$ is onto.

Corollary 6.12. The mapping $\bar{L}_{\nu}^{\Gamma_{1}}: H\left(L_{\partial}, \Omega\right) \rightarrow \mathcal{V}_{L, \Gamma_{1}}^{\prime}$ is linear, bounded and onto.
Proof. By Proposition 6.8 we already know that $\bar{L}_{\nu}^{\Gamma_{1}}$ is linear and bounded form $H\left(L_{\partial}, \Omega\right)$ to $\mathcal{V}_{L}^{\prime}$. Remark 6.9 gives $\left.\bar{L}_{\nu} f\right|_{\mathcal{V}_{L, \Gamma_{1}}}=\bar{L}_{\nu}^{\Gamma_{1}} f$ for $f \in H\left(L_{\partial}, \Omega\right)$ and $\mathcal{V}_{L, \Gamma_{1}}^{\prime}=$ $\left.\mathcal{V}_{L}^{\prime}\right|_{\mathcal{V}_{L, \Gamma_{1}}}$, which completes the proof.

Theorem 6.13. $\left(\mathcal{V}_{L, \Gamma_{1}}, L_{\pi}^{2}\left(\Gamma_{1}\right), \mathcal{V}_{L, \Gamma_{1}}^{\prime}\right)$ is a quasi Gelfand triple.
Proof. Let $\tilde{D}_{+}:=\left.\operatorname{ran} \pi_{L}\right|_{H_{\Gamma_{0}}^{1}(\Omega)^{m_{1}}}$ equipped with $\|\cdot\|_{\mathcal{X}_{+}}=\|\cdot\| \mathcal{V}_{L, \Gamma_{1}}$ and let $D_{-}$ denote the corresponding set from Definition 5.2 with $\mathcal{X}_{0}=L_{\pi}^{2}\left(\Gamma_{1}\right)$. Then by Remark $5.3\|g\|_{\mathcal{X}_{-}}=\|g\|_{\mathcal{V}_{L, \Gamma_{1}}^{\prime}}$ for $g \in D_{-}$and $\operatorname{ran} \mathbb{1}_{\Gamma_{1}} L_{\nu} \gamma_{0} \subseteq D_{-}$(by Proposition 6.8). By definition $\operatorname{ran} \mathbb{1}_{\Gamma_{1}} L_{\nu} \gamma_{0}$ is dense in $L_{\pi}^{2}\left(\Gamma_{1}\right)$ and by Proposition 6.8 and Corollary 6.12 also dense in $\mathcal{V}_{L, \Gamma_{1}}^{\prime}$. Consequently, also $D_{-}$is dense in both $L_{\pi}^{2}\left(\Gamma_{1}\right)$ and $\mathcal{V}_{L, \Gamma_{1}}^{\prime}$. Hence, assertion (iv) of Proposition 5.8 is satisfied, and by Remark 5.18 the completions of $\tilde{D}_{+}$and $D_{-}$form a quasi Gelfand triple with pivot space $L_{\pi}^{2}\left(\Gamma_{1}\right)$. By construction the completion of $\tilde{D}_{+}$is $\mathcal{V}_{L, \Gamma_{1}}$. By the density of $D_{-}$in $\mathcal{V}_{L, \Gamma_{1}}^{\prime}$ and $\|g\|_{\mathcal{X}_{-}}=\|g\|_{\mathcal{V}_{L, \Gamma_{1}}^{\prime}}$ for $g \in D_{-}$the completion of $D_{-}$is $\mathcal{V}_{L, \Gamma_{1}}^{\prime}$.

Corollary 6.14. $H_{0}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)=H_{\partial \Omega}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)=\operatorname{ker} \bar{\pi}_{L}=\operatorname{ker} \bar{L}_{\nu}^{\mathrm{H}}$ and $H_{0}\left(L_{\partial}, \Omega\right)=$ $H_{\partial \Omega}\left(L_{\partial}, \Omega\right)=\operatorname{ker} \bar{\pi}_{L^{H}}=\operatorname{ker} \bar{L}_{\nu}$.

Proof. For $g \in H_{0}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$ there is a sequence $\left(g_{n}\right)_{n \in \mathbb{N}}$ in $\mathcal{D}(\Omega)$ converging to $g$, which implies $\bar{\pi}_{L} g=\lim _{n \rightarrow \infty} \bar{\pi}_{L} g_{n}=0$. Therefore, $H_{0}\left(L_{\partial}^{\mathrm{H}}, \Omega\right) \subseteq \operatorname{ker} \bar{\pi}_{L}=H_{\partial \Omega}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$. On the other hand, if $g \in H_{\partial \Omega}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$, then

$$
\left\langle L_{\partial} f, g\right\rangle_{L^{2}(\Omega)^{m_{1}}}+\left\langle f, L_{\partial}^{\mathrm{H}} g\right\rangle_{L^{2}(\Omega)^{m_{2}}}=\left\langle\bar{L}_{\nu} f, \bar{\pi}_{L} g\right\rangle_{\mathcal{V}_{L}^{\prime}, \mathcal{V}_{L}}=0
$$

for all $f \in H\left(L_{\partial}, \Omega\right)$. Hence, by Lemma $3.17 g \in H_{0}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$. Consequently, $H_{0}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)=H_{\partial \Omega}\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$. The second equality of the statement holds by definition and the third will be proven by the following equivalences

$$
\begin{aligned}
g \in \operatorname{ker} \pi_{L} & \Leftrightarrow\left\langle\bar{\pi}_{L} g, \psi\right\rangle_{\mathcal{V}_{L}^{\prime}, \mathcal{V}_{L}}=0 \quad \text { for all } \psi \in \mathcal{V}_{L}^{\prime} \\
& \Leftrightarrow\left\langle\bar{\pi}_{L} g, \bar{L}_{\nu} f\right\rangle_{\mathcal{V}_{L}^{\prime}, \mathcal{V}_{L}}=0 \quad \text { for all } f \in H\left(L_{\partial}, \Omega\right) \\
& \stackrel{C .6 .10}{\Leftrightarrow}\left\langle\bar{L}_{\nu}^{\mathrm{H}} g, \bar{\pi}_{L^{\mathrm{H}}} f\right\rangle_{\mathcal{V}_{L^{\mathrm{H}}}^{\prime}}, \mathcal{V}_{L^{H}}=0 \quad \text { for all } f \in H\left(L_{\partial}, \Omega\right) \\
& \Leftrightarrow\left\langle\bar{L}_{\nu}^{\mathrm{H}} g, \phi\right\rangle_{\mathcal{V}_{L^{H}}^{\prime}}, \mathcal{V}_{L^{\mathrm{H}}}=0 \quad \text { for all } \phi \in \mathcal{V}_{L^{\mathrm{H}}} \\
& \Leftrightarrow g \in \operatorname{ker} \bar{L}_{\nu}^{\mathrm{H}} .
\end{aligned}
$$

Switching $L$ with $L^{\mathrm{H}}$ yields $H_{0}\left(L_{\partial}, \Omega\right)=H_{\partial \Omega}\left(L_{\partial}, \Omega\right)=\operatorname{ker} \bar{\pi}_{L^{H}}=\operatorname{ker} \bar{L}_{\nu}$.
7. Existence and uniqueness via boundary triples. In this section we will show that there is a boundary triple associated to the port-Hamiltonian differential operator $\left(P_{\partial}+P_{0}\right) \mathcal{H}$, which enables us to formulate boundary conditions that admit existence and uniqueness of solutions. Moreover, we will parameterize all boundary conditions that provide unique solutions that are non-increasing in the Hamiltonian.

Recall the setting in section 4. Using $\pi_{P}=\left[\begin{array}{cc}\pi_{L} & 0 \\ 0 & \pi_{L^{\mathrm{H}}}\end{array}\right]$, Lemma 6.4 and Proposition 6.8, it is easy to see that $\mathcal{V}_{P}=\mathcal{V}_{L} \times \mathcal{V}_{L^{H}}$ and therefore $\mathcal{V}_{P}^{\prime}=\mathcal{V}_{L}^{\prime} \times \mathcal{V}_{L^{H}}^{\prime}$. Furthermore, for $\bar{P}_{\nu}: H\left(P_{\partial}, \Omega\right) \rightarrow \mathcal{V}_{P}^{\prime}$ and $\bar{\pi}_{P}: H\left(P_{\partial}, \Omega\right) \rightarrow \mathcal{V}_{P}$ we have

$$
\bar{P}_{\nu}=\left[\begin{array}{cc}
0 & \bar{L}_{\nu} \\
\bar{L}_{\nu}^{\mathrm{H}} & 0
\end{array}\right] \quad \text { and } \quad \bar{\pi}_{P}=\left[\begin{array}{cc}
\bar{\pi}_{L} & 0 \\
0 & \bar{\pi}_{L^{\mathrm{H}}}
\end{array}\right] .
$$

Recall the splitting $x=\left[\begin{array}{c}x_{L^{H}} \\ x_{L}\end{array}\right]$. Accordingly, we introduce $\mathcal{H} x=\left[\begin{array}{c}(\mathcal{H} x)_{L^{H}} \\ (\mathcal{H} x)_{L}\end{array}\right]$ for $x \in \mathcal{H}^{-1}\left(H\left(P_{\partial}, \Omega\right)\right)$, so that

$$
P_{\partial} \mathcal{H} x=\left[\begin{array}{c}
L_{\partial}(\mathcal{H} x)_{L} \\
L_{\partial}^{H}(\mathcal{H} x)_{L^{\mathrm{H}}}
\end{array}\right], \quad\left[\begin{array}{ll}
0 & \bar{L}_{\nu}
\end{array}\right] \mathcal{H} x=\bar{L}_{\nu}(\mathcal{H} x)_{L}, \quad\left[\begin{array}{ll}
\bar{\pi}_{L} & 0
\end{array}\right] \mathcal{H} x=\bar{\pi}_{L}(\mathcal{H} x)_{L^{H}} .
$$

Theorem 7.1. The operator

$$
A_{0}:=-\left(P_{\partial}+P_{0}\right) \mathcal{H}, \quad \operatorname{dom} A_{0}:=\mathcal{H}^{-1}\left(\operatorname{ker} \bar{P}_{\nu}\right)
$$

is closed, skew-symmetric, and densely defined on $\mathcal{X}_{\mathcal{H}}$. Its adjoint is

$$
A_{0}^{*}=\left(P_{\partial}+P_{0}\right) \mathcal{H}, \quad \operatorname{dom} A_{0}^{*}=\mathcal{H}^{-1}\left(H\left(P_{\partial}, \Omega\right)\right) .
$$

Let $B_{1}=\left[\begin{array}{ll}\bar{\pi}_{L} & 0\end{array}\right] \mathcal{H}, B_{2}=\left[\begin{array}{ll}0 & \bar{L}_{\nu}\end{array}\right] \mathcal{H}$ and $\Psi$ be the duality map of $\left(\mathcal{V}_{L}, L_{\pi}^{2}(\partial \Omega), \mathcal{V}_{L}^{\prime}\right)$. Then $\left(\mathcal{V}_{L}, B_{1}, \Psi B_{2}\right)$ is a boundary triple for $A_{0}^{*}$.
Proof. We define $\tilde{A}$ as $\left(P_{\partial}+P_{0}\right) \mathcal{H}$ with $\operatorname{dom} \tilde{A}=\mathcal{H}^{-1}\left(H\left(P_{\partial}, \Omega\right)\right)$ on $\mathcal{X}_{\mathcal{H}}$. By Lemma $3.5 P_{\partial}: H\left(P_{\partial}, \Omega\right) \subseteq L^{2}(\Omega)^{m} \rightarrow L^{2}(\Omega)^{m}$ is a closed operator. Since $\mathcal{H}$ is a bounded operator on $L^{2}(\Omega)^{m}$, and $\mathcal{X}_{\mathcal{H}}$ and $L^{2}(\Omega)^{m}$ have equivalent norms, it is easy to see that $\tilde{A}: \mathcal{H}^{-1}\left(H\left(P_{\partial}, \Omega\right)\right) \subseteq \mathcal{X}_{\mathcal{H}} \rightarrow \mathcal{X}_{\mathcal{H}}$ is closed. Let $B^{*_{H}}$ denote the adjoint of $B$ with respect to $\langle., .\rangle_{H}$ for any Hilbert space $H$. The adjoint of $\tilde{A}$ can be calculated by

$$
\tilde{A}^{*}=\left(\left(P_{\partial}+P_{0}\right) \mathcal{H}\right)^{* \chi_{\mathcal{H}}}=\mathcal{H}^{-1}\left(\left(P_{\partial}+P_{0}\right) \mathcal{H}\right)^{* L^{2}} \mathcal{H}=\left(P_{\partial}^{* L^{2}}+P_{0}^{*} L^{2}\right) \mathcal{H}
$$

and according to Remark 3.7 we have $P_{\partial}^{* L^{2}}=-\left.P_{\partial}\right|_{\operatorname{dom} P_{\partial}^{*} L^{2}}$ where $\operatorname{dom} P_{\partial}^{*_{L^{2}}} \subseteq$ $H\left(P_{\partial}, \Omega\right)$. Hence,

$$
\tilde{A}^{*}=-\left.\left(P_{\partial}+P_{0}\right) \mathcal{H}\right|_{\mathcal{H}^{-1}\left(\operatorname{dom} P_{\partial}^{*} L^{2}\right)}=-\left.\tilde{A}\right|_{\mathcal{H}^{-1}\left(\operatorname{dom} P_{\partial}^{*} L^{2}\right)} \subseteq-\tilde{A}
$$

Consequently, $\tilde{A}^{*}$ is skew-symmetric on $\mathcal{X}_{\mathcal{H}}$. Since $\tilde{A}$ is closed, we have $\tilde{A}^{* *}=\tilde{A}$.
Now we know that $\tilde{A}$ is the adjoint of a skew-symmetric operator. So we can talk about boundary triples for $\tilde{A}$. First we note that

$$
\operatorname{ran}\left[\begin{array}{c}
B_{1} \\
\Psi B_{2}
\end{array}\right]=\operatorname{ran} \bar{\pi}_{L} \times \operatorname{ran} \Psi \bar{L}_{\nu}=\mathcal{V}_{L} \times \mathcal{V}_{L}
$$

Since $\mathcal{H}$ is self-adjoint and $P_{0}$ is skew-adjoint, we have for $x, y \in \operatorname{dom} \tilde{A}$

$$
\begin{aligned}
\langle\tilde{A} x, y\rangle_{\mathcal{X}_{\mathcal{H}}} & +\langle x, \tilde{A} y\rangle_{\mathcal{X}_{\mathcal{H}}} \\
& =\left\langle P_{\partial} \mathcal{H} x, \mathcal{H} y\right\rangle_{L^{2}}+\left\langle\mathcal{H} x, P_{\partial} \mathcal{H} y\right\rangle_{L^{2}}
\end{aligned}
$$

by the the identity $P_{\partial}=\left[\begin{array}{cc}0 & L_{\partial} \\ L_{\partial}^{\mathrm{H}} & 0\end{array}\right]$ and Corollary 6.10 we further have

$$
\begin{aligned}
= & \left\langle\left[\begin{array}{c}
L_{\partial}(\mathcal{H} x)_{L} \\
L_{\partial}^{\mathrm{H}}(\mathcal{H} x)_{L^{\mathrm{H}}}
\end{array}\right],\left[\begin{array}{c}
(\mathcal{H} y)_{L^{\mathrm{H}}} \\
(\mathcal{H} y)_{L}
\end{array}\right]\right\rangle_{L^{2}}+\left\langle\left[\begin{array}{c}
(\mathcal{H} x)_{L^{\mathrm{H}}} \\
(\mathcal{H} x)_{L}
\end{array}\right],\left[\begin{array}{c}
L_{\partial}(\mathcal{H} y)_{L} \\
L_{\partial}^{\mathrm{H}}(\mathcal{H} y)_{L^{\mathrm{H}}}
\end{array}\right]\right\rangle_{L^{2}} \\
= & \left\langle L_{\partial}(\mathcal{H} x)_{L},(\mathcal{H} y)_{L^{\mathrm{H}}}\right\rangle_{L^{2}}+\left\langle(\mathcal{H} x)_{L}, L_{\partial}^{\mathrm{H}}(\mathcal{H} y)_{L^{\mathrm{H}}}\right\rangle_{L^{2}} \\
& \quad+\left\langle L_{\partial}^{\mathrm{H}}(\mathcal{H} x)_{L^{\mathrm{H}}},(\mathcal{H} y)_{L}\right\rangle_{L^{2}}+\left\langle(\mathcal{H} x)_{L^{\mathrm{H}}}, L_{\partial}^{\mathrm{H}}(\mathcal{H} y)_{L}\right\rangle_{L^{2}} \\
= & \left\langle\bar{L}_{\nu}(\mathcal{H} x)_{L}, \bar{\pi}_{L}(\mathcal{H} y)_{L^{\mathrm{H}}}\right\rangle_{\mathcal{V}_{L}^{\prime}, \mathcal{V}_{L}}+\left\langle\bar{\pi}_{L}(\mathcal{H} x)_{L^{\mathrm{H}}}, \bar{L}_{\nu}(\mathcal{H} y)_{L}\right\rangle_{\mathcal{V}_{L}, \mathcal{V}_{L}^{\prime}} \\
= & \left\langle\Psi B_{2} x, B_{1} y\right\rangle_{\mathcal{V}_{L}}+\left\langle B_{1} x, \Psi B_{2} y\right\rangle_{\mathcal{V}_{L}} .
\end{aligned}
$$

Therefore, $\left(\mathcal{V}_{L}, B_{1}, \Psi B_{2}\right)$ is a boundary triple for $\tilde{A}$.
By Lemma $2.2 \operatorname{dom} \tilde{A}^{*}=\operatorname{ker} B_{1} \cap \operatorname{ker} B_{2}$, which is equal to

$$
\operatorname{ker} B_{1} \cap \operatorname{ker} B_{2}=\mathcal{H}^{-1}\left(\operatorname{ker}\left[\begin{array}{ll}
\bar{\pi}_{L} & 0
\end{array}\right] \cap \operatorname{ker}\left[\begin{array}{ll}
0 & \bar{L}_{\nu}
\end{array}\right]\right)=\mathcal{H}^{-1}\left(\operatorname{ker} \bar{\pi}_{L} \times \operatorname{ker} \bar{L}_{\nu}\right)
$$

By Corollary 6.14 this is equal to $\mathcal{H}^{-1}\left(\operatorname{ker} \bar{L}_{\nu}^{\mathrm{H}} \times \operatorname{ker} \bar{L}_{\nu}\right)=\mathcal{H}^{-1}\left(\operatorname{ker} \bar{P}_{\nu}\right)$. Hence, $\tilde{A}^{*}=A_{0}$ and $A_{0}^{*}=\tilde{A}$.

Remark 7.2. We can replace $\left(\mathcal{V}_{L}, B_{1}, \Psi B_{2}\right)$ by $\left(\mathcal{V}_{L}^{\prime}, \Psi^{*} B_{1}, B_{2}\right)$ in the previous theorem.

Theorem 7.3. Let $A_{0}^{*}$ be the operator from the previous theorem and $\Psi_{\Gamma_{1}}$ the duality map associated to the quasi Gelfand triple $\left(\mathcal{V}_{L, \Gamma_{1}}, L_{\pi}^{2}\left(\Gamma_{1}\right), \mathcal{V}_{L, \Gamma_{1}}^{\prime}\right)$. Then we have $\left(\mathcal{V}_{L, \Gamma_{1}},\left[\begin{array}{ll}\bar{\pi}_{L} & 0\end{array}\right] \mathcal{H}, \Psi_{\Gamma_{1}}\left[\begin{array}{ll}0 & \bar{L}_{\nu}^{\Gamma_{1}}\end{array}\right] \mathcal{H}\right)$ as a boundary triple for

$$
A:=\left.A_{0}^{*}\right|_{\mathcal{H}^{-1}\left(H_{\Gamma_{0}}\left(L_{\partial}^{H}, \Omega\right) \times H\left(L_{\partial}, \Omega\right)\right)} .
$$

Proof. Since we already have a boundary triple for $A_{0}^{*}$, we can show that $A$ is the adjoint of a skew-symmetric operator by Proposition 2.3 (iii). Hence, we have to check, whether $\left[\begin{array}{cc}0 & \mathrm{I} \\ \mathrm{I} & 0\end{array}\right] \mathcal{C}^{\perp} \subseteq \mathcal{C}$ in $\mathcal{V}_{L} \times \mathcal{V}_{L}$, where $\mathcal{C}$ is the corresponding relation to the domain of $A$ according to Proposition 2.3. For $B_{1}, B_{2}$ being the mappings from the previous theorem we have (Note that $\mathcal{V}_{L, \Gamma_{1}}$ is a subspace of $\mathcal{V}_{L}$; Lemma 6.7)

$$
\begin{aligned}
\mathcal{C} & =\left[\begin{array}{c}
B_{1} \\
\Psi B_{2}
\end{array}\right] \operatorname{dom} A=\mathcal{V}_{L, \Gamma_{1}} \times \mathcal{V}_{L} \\
{\left[\begin{array}{ll}
0 & \mathrm{I} \\
\mathrm{I} & 0
\end{array}\right] \mathcal{C}^{\perp} } & =\{0\} \times \mathcal{V}_{L, \Gamma_{1}}^{\perp} \subseteq \mathcal{V}_{L, \Gamma_{1}} \times \mathcal{V}_{L}=\mathcal{C}
\end{aligned}
$$

For $x, y \in \operatorname{dom} A$ we have, using Remark 6.9,

$$
\begin{aligned}
\left\langle B_{1} x, \Psi B_{2} y\right\rangle_{\mathcal{V}_{L}} & =\left\langle\bar{\pi}_{L}(\mathcal{H} x)_{L^{\mathrm{H}}}, \bar{L}_{\nu}(\mathcal{H} y)_{L}\right\rangle_{\mathcal{V}_{L}, \mathcal{V}_{L}^{\prime}}=\left\langle\bar{\pi}_{L}(\mathcal{H} x)_{L^{\mathrm{H}}}, \bar{L}_{\nu}^{\Gamma_{1}}(\mathcal{H} y)_{L}\right\rangle_{\mathcal{V}_{L, \Gamma_{1}}, \mathcal{V}_{L, \Gamma_{1}}^{\prime}} \\
& =\left\langle\left[\begin{array}{ll}
\bar{\pi}_{L} & 0
\end{array}\right] \mathcal{H} x, \Psi_{\Gamma_{1}}\left[\begin{array}{ll}
0 & \bar{L}_{\nu}^{\Gamma_{1}}
\end{array}\right] \mathcal{H} y\right\rangle_{\mathcal{V}_{L, \Gamma_{1}}}
\end{aligned}
$$

which yields item (ii) in Definition 2.1. By $\left.\operatorname{ran}\left[\begin{array}{cc}\bar{\pi}_{L} & 0 \\ 0 & \Psi_{\Gamma_{1}} \bar{L}_{\nu}^{\Gamma_{1}}\end{array}\right]\right|_{H_{\Gamma_{0}}\left(L_{\partial}^{H}, \Omega\right) \times H\left(L_{\partial}, \Omega\right)}=$ $\mathcal{V}_{L, \Gamma_{1}} \times \mathcal{V}_{L, \Gamma_{1}}$, the remaining item (i) is fulfilled.

The next theorem is [9, Theorem 2.5].
Theorem 7.4. Let $A_{0}$ be a skew-symmetric operator on a Hilbert space $X$ and $\left(\mathcal{B}, B_{1}, B_{2}\right)$ be a boundary triple for $A_{0}^{*}$. Furthermore let $\mathcal{K}$ be a Hilbert space, $W_{B}=\left[\begin{array}{ll}W_{1} & W_{2}\end{array}\right]$, where $W_{1}, W_{2} \in \mathcal{L}(\mathcal{B}, \mathcal{K})$, and $A:=\left.A_{0}^{*}\right|_{\operatorname{dom} A}$, where $\operatorname{dom} A=$ $\operatorname{ker} W_{B}\left[\begin{array}{l}B_{1} \\ B_{2}\end{array}\right]$. If ran $W_{1}-W_{2} \subseteq \operatorname{ran} W_{1}+W_{2}$ then the following assertions are equivalent.
(i) The operator $A$ generates a contraction semigroup on $X$.
(ii) The operator $A$ is dissipative.
(iii) The operator $W_{1}+W_{2}$ is injective and the following operator inequality holds

$$
W_{1} W_{2}^{*}+W_{2} W_{1}^{*} \geq 0
$$

We will reformulate this theorem to fit our situation.
Corollary 7.5. Let $\mathcal{K}$ be some Hilbert space and $W=\left[\begin{array}{ll}W_{1} & W_{2}\end{array}\right]: \mathcal{V}_{L, \Gamma_{1}} \times \mathcal{V}_{L, \Gamma_{1}} \rightarrow$ $\mathcal{K}$ a bounded linear mapping such that $\operatorname{ran} W_{1}-W_{2} \subseteq \operatorname{ran} W_{1}+W_{2}$. Let

$$
\begin{aligned}
& D:=\left\{x \in \mathcal{H}^{-1}\left(H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right) \times H\left(L_{\partial}, \Omega\right)\right)\right. \\
& \\
& \left.: W_{1}\left[\begin{array}{ll}
\bar{\pi}_{L} & 0
\end{array}\right] \mathcal{H} x+W_{2} \Psi\left[\begin{array}{ll}
0 & \bar{L}_{\nu}^{\Gamma_{1}}
\end{array}\right] \mathcal{H} x=0\right\}
\end{aligned}
$$

where $\Psi: \mathcal{V}_{L, \Gamma_{1}}^{\prime} \rightarrow \mathcal{V}_{L, \Gamma_{1}}$ is the duality mapping corresponding to the quasi Gelfand triple. Then the following assertions are equivalent.
(i) $\left.\left(P_{\partial}+P_{0}\right) \mathcal{H}\right|_{D}$ generates a contraction semigroup.
(ii) $\left.\left(P_{\partial}+P_{0}\right) \mathcal{H}\right|_{D}$ is dissipative.
(iii) The operator $W_{1}+W_{2}$ is injective and the following operator inequality holds

$$
W_{1} W_{2}^{*}+W_{2} W_{1}^{*} \geq 0
$$

Corollary 7.5 already gives a parameterization via $W$ for all boundary conditions that make $\left(P_{\partial}+P_{0}\right) \mathcal{H}$ a generator of a contraction semigroup. In particular the corresponding PDEs have unique solutions that continuously depend on the initial state and don't grow in the Hamiltonian. However, checking continuity for boundary operators which map into $\mathcal{V}_{L}$ can be difficult. Hence, it would be appreciated to reduce the conditions on the boundary operators to conditions on better known spaces like the pivot space $L^{2}(\partial \Omega)$. The next theorem will provide this.

The following result is a generalization of [9, Theorem 2.6] for quasi Gelfand triple and also fixes some minor issues, like the specific choice of $\Psi$ and the closedness of $\left[\left.\left.V_{1}\right|_{\mathcal{B}_{+} \cap \mathcal{B}_{0}} \quad V_{2}\right|_{\mathcal{B}_{-} \cap \mathcal{B}_{0}}\right]$ as an operator from $\mathcal{B}_{+} \times \mathcal{B}_{-}$to $\mathcal{K}$.
Theorem 7.6. Let $\left(\mathcal{B}_{+}, \mathcal{B}_{0}, \mathcal{B}_{-}\right)$be a quasi Gelfand triple, $A_{0}$ be a closed skewsymmetric operator and $\left(\mathcal{B}_{+}, B_{1}, \Psi B_{2}\right)$ be a boundary triple for $A_{0}^{*}$, where $\Psi$ is the duality map of the quasi Gelfand triple. For $V_{1}, V_{2} \in \mathcal{L}\left(\mathcal{B}_{0}, \mathcal{K}\right)$ we define

$$
D:=\left\{a \in \operatorname{dom} A_{0}^{*}: B_{1} a, B_{2} a \in \mathcal{B}_{0} \quad \text { and } \quad\left[\begin{array}{ll}
V_{1} & V_{2}
\end{array}\right]\left[\begin{array}{l}
B_{1} \\
B_{2}
\end{array}\right] a=0\right\}
$$

and the operator $A:=\left.A_{0}^{*}\right|_{D}$. If
(i) $\left[\left.\left.V_{1}\right|_{\mathcal{B}_{0} \cap \mathcal{B}_{+}} \quad V_{2}\right|_{\mathcal{B}_{0} \cap \mathcal{B}_{-}}\right]$is closed as an operator from $\mathcal{B}_{+} \times \mathcal{B}_{-}$to $\mathcal{K}$,
(ii) $\operatorname{ker}\left[\begin{array}{ll}V_{1} & V_{2}\end{array}\right]$ is dissipative as linear relation on $\mathcal{B}_{0}$,
(iii) $V_{1} V_{2}^{*}+V_{2} V_{1}^{*} \geq 0$ as operator on $\mathcal{K}$,
then $A$ is a generator of a contraction semigroup.
Proof. It is sufficient to show that $A$ is closed, and $A$ and $A^{*}$ are dissipative.
Step 1. Showing that $A$ is closed and dissipative. We have

$$
\begin{aligned}
a \in D & \Leftrightarrow\left[\begin{array}{l}
B_{1} \\
B_{2}
\end{array}\right] a \in\left(\mathcal{B}_{0} \times \mathcal{B}_{0}\right) \cap \operatorname{ker}\left[\begin{array}{ll}
V_{1} & V_{2}
\end{array}\right] \\
& \Leftrightarrow\left[\begin{array}{c}
B_{1} \\
\Psi B_{2}
\end{array}\right] a \in \underbrace{\operatorname{ker}\left[\left.\left.V_{1}\right|_{\mathcal{B}_{0} \cap \mathcal{B}_{+}} \quad V_{2} \Psi^{*}\right|_{\Psi\left(\mathcal{B}_{0} \cap \mathcal{B}_{-}\right)}\right]}_{=: \mathcal{C}} .
\end{aligned}
$$

We can write

$$
\mathcal{C}=\left\{\left[\begin{array}{l}
q \\
p
\end{array}\right] \in \mathcal{B}_{+} \times \mathcal{B}_{+}: q \in \mathcal{B}_{0}, \exists \tilde{p} \in \mathcal{B}_{0}: p=\Psi \tilde{p}, V_{1} q+V_{2} \Psi^{*} p=0\right\}
$$

For $\left[\begin{array}{l}q \\ p\end{array}\right] \in \mathcal{C}$ we have

$$
\operatorname{Re}\langle q, p\rangle_{\mathcal{B}_{+}}=\operatorname{Re}\langle q, \Psi \tilde{p}\rangle_{\mathcal{B}_{+}}=\operatorname{Re}\langle q, \tilde{p}\rangle_{\mathcal{B}_{+}, \mathcal{B}_{-}}=\operatorname{Re}\langle q, \tilde{p}\rangle_{\mathcal{B}_{0}} \leq 0
$$

which implies the dissipativity of $A$ by Proposition 2.3. Assumption (i) implies that $\mathcal{C}$ is closed in $\mathcal{B}_{+}^{2}$, which implies the closedness of $A$ by Proposition 2.3.
Step 2. Showing that $A^{*}$ is dissipative. By Proposition 2.3 we can characterize the domain of $A^{*}$ by

$$
\begin{aligned}
d \in \operatorname{dom} A^{*} & \Leftrightarrow\left[\begin{array}{c}
B_{1} \\
\Psi B_{2}
\end{array}\right] d \in\left[\begin{array}{cc}
0 & \mathrm{I} \\
\mathrm{I} & 0
\end{array}\right] \mathcal{C}^{\perp_{\mathcal{B}_{+}^{2}}} \\
& \Leftrightarrow\left[\begin{array}{c}
\Psi B_{2} \\
B_{1}
\end{array}\right] d \in \operatorname{ran}\left[\begin{array}{c}
\left(\left.V_{1}\right|_{\mathcal{B}_{0} \cap \mathcal{B}_{+}}\right)^{*_{\mathcal{B}_{+}}} \\
\left(\left.V_{2} \Psi^{*}\right|_{\Psi\left(\mathcal{B}_{0} \cap \mathcal{B}_{-}\right)}\right)^{* \mathcal{B}_{+}}
\end{array}\right]^{\mathcal{B}_{+}^{2}} .
\end{aligned}
$$

The second equivalence needed the closedness in assumption (i), since $(\operatorname{ker} T)^{\perp}=$ $\overline{\operatorname{ran} T^{*}}$ for a linear relation (or even unbounded operator) $T$ is not true in general. Note that if $P$ is a bounded and everywhere defined operator, and $Q$ is a linear relation, then $(P Q)^{*}=Q^{*} P^{*}$. Hence, by Proposition 5.16

$$
\left(\left.V_{1}\right|_{\mathcal{B}_{0} \cap \mathcal{B}_{+}}\right)^{* \mathcal{B}_{+}}=\left(V_{1} \iota_{+}\right)^{*}=\iota_{+}^{*} V_{1}^{*}=\left.\Psi V_{1}^{*}\right|_{V_{1}^{*-1}\left(\mathcal{B}_{0} \cap \mathcal{B}_{-}\right)}
$$

where $\iota_{+}: \mathcal{B}_{+} \cap \mathcal{B}_{0} \subseteq \mathcal{B}_{+} \rightarrow \mathcal{B}_{0}$ is one embedding of the quasi Gelfand triple and

$$
\left(\left.V_{2} \Psi^{*}\right|_{\Psi\left(\mathcal{B}_{0} \cap \mathcal{B}_{-}\right)}\right)^{* \mathcal{B}_{+}}=\left(V_{2} \iota_{-} \Psi^{*}\right)^{*}=\left(\iota_{-} \Psi^{*}\right)^{*} V_{2}^{*},
$$

where $\iota_{-}: \mathcal{B}_{-} \cap \mathcal{B}_{0} \subseteq \mathcal{B}_{-} \rightarrow \mathcal{B}_{0}$ is the other embedding of the quasi Gelfand triple. From $\left(\Psi \iota_{-}^{*}\right)^{*}=\iota_{-} \Psi^{*}$ and $\iota_{-}^{*}=\Psi^{*} \iota_{+}^{-1}$ (Proposition 5.16) follows $\left(\iota_{-} \Psi^{*}\right)^{*}=\overline{\Psi \iota_{-}^{*}}=$ $\iota_{+}^{-1}$. Consequently,

$$
\left(\left.V_{2} \Psi^{*}\right|_{\Psi\left(\mathcal{B}_{0} \cap \mathcal{B}_{-}\right)}\right)^{* \mathcal{B}_{+}}=\iota_{+}^{-1} V_{2}^{*}=\left.V_{2}^{*}\right|_{V_{2}^{*-1}\left(\mathcal{B}_{0} \cap \mathcal{B}_{+}\right)} .
$$

Hence, for

$$
\begin{array}{r}
{\left[\begin{array}{l}
x \\
y
\end{array}\right] \in \operatorname{ran}\left[\begin{array}{c}
\left(\left.V_{1}\right|_{\mathcal{B}_{0} \cap \mathcal{B}_{+}}\right)^{* \mathcal{B}_{+}} \\
\left(\left.V_{2} \Psi^{*}\right|_{\Psi\left(\mathcal{B}_{0} \cap \mathcal{B}_{-}\right)}\right)^{* \mathcal{B}_{+}}
\end{array}\right]} \\
\left.=\left\{\begin{array}{c}
\Psi V_{1}^{*} \\
V_{2}^{*}
\end{array}\right] k: k \in V_{1}^{*-1}\left(\mathcal{B}_{0} \cap \mathcal{B}_{-}\right) \cap V_{2}^{*-1}\left(\mathcal{B}_{0} \cap \mathcal{B}_{+}\right)\right\}
\end{array}
$$

we have

$$
\begin{aligned}
\operatorname{Re}\langle x, y\rangle_{\mathcal{B}_{+}} & =\operatorname{Re}\left\langle\Psi V_{1}^{*} k, V_{2}^{*} k\right\rangle_{\mathcal{B}_{+}}=\operatorname{Re}\left\langle V_{1}^{*} k, V_{2}^{*} k\right\rangle_{\mathcal{B}_{-}, \mathcal{B}_{+}}=\operatorname{Re}\left\langle V_{1}^{*} k, V_{2}^{*} k\right\rangle_{\mathcal{B}_{0}} \\
& =\operatorname{Re}\left\langle V_{2} V_{1}^{*} k, k\right\rangle_{\mathcal{K}} \geq 0 .
\end{aligned}
$$

Therefore, $\mathcal{C}^{\perp}$ is accretive and by Proposition 2.3 also $\left.A_{0}\right|_{\text {dom } A^{*}}$ is accretive, which yields $A^{*}=-\left.A_{0}\right|_{\text {dom } A^{*}}$ is dissipative.

Remark 7.7. If we are already satisfied with the operator closure $\bar{A}$ is a generator (instead of $A$ ) in the previous theorem, then we can replace condition (i) by

$$
\operatorname{ker} \overline{\left[\begin{array}{ll}
\left.V_{1}\right|_{\mathcal{B}_{0} \cap \mathcal{B}_{+}} & \left.\left.V_{2}\right|_{\mathcal{B}_{0} \cap \mathcal{B}_{-}}\right]  \tag{7.1}\\
& \overline{\operatorname{ker}\left[\begin{array}{ll}
\left.V_{1}\right|_{\mathcal{B}_{0} \cap \mathcal{B}_{+}} & \left.V_{2}\right|_{\mathcal{B}_{0} \cap \mathcal{B}_{-}}
\end{array}\right]}{ }^{\mathcal{B}_{+} \times \mathcal{B}_{-}}
\end{array}, ., ~\right.}
$$

where $\overline{\left[\left.\left.V_{1}\right|_{\mathcal{B}_{0} \cap \mathcal{B}_{+}} V_{2}\right|_{\mathcal{B}_{0} \cap \mathcal{B}_{-}}\right]}$is the closure as linear relation (possibly multi-valued). Clearly, if (7.1) holds, then there is already equality.

Example 7.8. Let $\left(\mathcal{B}_{+}, \mathcal{B}_{0}, \mathcal{B}_{-}\right)$be a quasi Gelfand triple that satisfies all conditions of Theorem 7.6 and let $M \in \mathcal{L}\left(\mathcal{B}_{0}\right)$ be coercive (i.e. $M \geq c \mathrm{I}, c>0$ ). Then $V_{1}:=\mathrm{I}, V_{2}:=M$ fulfill all conditions of Theorem 7.6:
(i) Setting $S=M^{\frac{1}{2}}$ and $T=M^{-\frac{1}{2}}$ in Corollary 5.27 implies the closedness of $\left[\left.\left.\mathrm{I}\right|_{\mathcal{B}_{0} \cap \mathcal{B}_{+}} M\right|_{\mathcal{B}_{0} \cap \mathcal{B}_{-}}\right]$.
(ii) For $(x, y) \in \operatorname{ker}\left[\begin{array}{ll}V_{1} & V_{2}\end{array}\right]$ we have $x=-M y$. Since $M$ is positive this yields

$$
\operatorname{Re}\langle x, y\rangle_{\mathcal{B}_{0}}=\operatorname{Re}\langle-M y, y\rangle=-\langle M y, y\rangle \leq 0
$$

(iii) $V_{1} V_{2}^{*}+V_{2} V_{1}^{*}=M^{*}+M=2 \operatorname{Re} M \geq 0$.

Moreover, Corollary 5.27 also implies the surjectivity of $\left[\left.\left.\mathrm{I}\right|_{\mathcal{B}_{0} \cap \mathcal{B}_{+}} M\right|_{\mathcal{B}_{0} \cap \mathcal{B}_{-}}\right]$.
Actually, it would have been enough, if $M \in \mathcal{L}\left(\mathcal{B}_{0}\right)$ was boundedly invertible and accretive. Clearly, also $V_{1}:=M, V_{2}:=\mathrm{I}$ fulfill all conditions.
8. Port-hamiltonian systems as boundary control systems. We will recall the notion of boundary control systems, scattering passive and impedance passive in the manner of [12]. We will show that a port-Hamiltonian system can be described as such a system. This concept already provides solution theory (see i.e. [11, Lemma 2.6]). It is well known that every scattering passive boundary control system induces a scattering passive well-posed linear system.
Definition 8.1. A colligation $\Xi:=\left(\left[\begin{array}{c}G \\ L \\ K\end{array}\right] ;\left[\begin{array}{c}\mathcal{U} \\ \mathcal{X} \\ \mathcal{Y}\end{array}\right]\right)$ consists of the three Hilbert spaces $\mathcal{U}, \mathcal{X}$, and $\mathcal{Y}$, and the three linear maps $G, L$, and $K$, with the same domain $\mathcal{Z} \subseteq \mathcal{X}$ and with values in $\mathcal{U}, \mathcal{X}$, and $\mathcal{Y}$, respectively.
Definition 8.2. A colligation $\Xi:=\left(\left[\begin{array}{c}G \\ L \\ K\end{array}\right] ;\left[\begin{array}{c}\mathcal{U} \\ \mathcal{X} \\ \mathcal{Y}\end{array}\right]\right)$ is an (internally well-posed) boundary control system, if
(i) the operator $\left[\begin{array}{c}G \\ L \\ K\end{array}\right]$ is closed from $\mathcal{X}$ to $\left[\begin{array}{l}\mathcal{U} \\ \mathcal{X} \\ \mathcal{Y}\end{array}\right]$,
(ii) the operator $G$ is surjective, and
(iii) the operator $A:=\left.L\right|_{\operatorname{ker} G}$ generates a contraction semigroup on $\mathcal{X}$.

We think of the operators in this definition as determining a system via

$$
\begin{align*}
u(t) & =G x(t) \\
\dot{x}(t) & =L x(t), \quad x(0)=x_{0}  \tag{8.1}\\
y(t) & =K x(t)
\end{align*}
$$

We call $\mathcal{U}$ the input space, $\mathcal{X}$ the state space, $\mathcal{Y}$ the output space and $\mathcal{Z}$ the solution space.
Definition 8.3. Let $\Xi=\left(\left[\begin{array}{l}G \\ L \\ K\end{array}\right] ;\left[\begin{array}{l}\mathcal{U} \\ \mathcal{X} \\ \mathcal{Y}\end{array}\right]\right)$ be a colligation. If $\Xi$ is a boundary control system such that

$$
\begin{equation*}
2 \operatorname{Re}\langle L x, x\rangle_{\mathcal{X}}+\|K x\|_{\mathcal{Y}}^{2} \leq\|G x\|_{\mathcal{U}}^{2} \quad \text { for } \quad x \in \mathcal{Z} \tag{8.2}
\end{equation*}
$$

then it is scattering passive and it is scattering energy preserving if we have equality in (8.2).

We say $\Xi$ is impedance passive (energy preserving), if $\mathcal{Y}=\mathcal{U}^{\prime}, \Psi: \mathcal{U}^{\prime} \rightarrow \mathcal{U}$ is the unitary identification mapping and $\tilde{\Xi}:=\left(\left[\begin{array}{c}\frac{1}{\sqrt{2}}(G+\Psi K) \\ \frac{1}{\sqrt{2}}(G-\Psi K)\end{array}\right] ;\left[\begin{array}{c}\mathcal{U} \\ \mathcal{U}\end{array}\right]\right)$ is scattering passive (energy preserving).

Note that an impedance passive (energy preserving) colligation $\Xi$ does not need to be a boundary control system. If $\mathcal{U}=\mathcal{Y}$, then $\Psi$ is the identity mapping.

Corresponding to a port-Hamiltonian system we want to introduce the following operators

$$
\begin{array}{rll}
G_{\mathrm{p}}:=S_{+}\left[\begin{array}{ll}
\bar{\pi}_{L} & 0
\end{array}\right] \mathcal{H}: & \mathcal{H}^{-1}\left(H\left(P_{\partial}, \Omega\right)\right) \subseteq \mathcal{X}_{\mathcal{H}} \rightarrow S \mathcal{V}_{L} \\
L_{\mathrm{p}} & :=\left(P_{\partial}+P_{0}\right) \mathcal{H}: & \mathcal{H}^{-1}\left(H\left(P_{\partial}, \Omega\right)\right) \subseteq \mathcal{X}_{\mathcal{H}} \rightarrow \mathcal{X}_{\mathcal{H}} \\
K_{\mathrm{p}} & :=\left(S^{*}\right)_{-}^{-1}\left[\begin{array}{ll}
0 & \bar{L}_{\nu}
\end{array}\right] \mathcal{H}: & \mathcal{H}^{-1}\left(H\left(P_{\partial}, \Omega\right)\right) \subseteq \mathcal{X}_{\mathcal{H}} \rightarrow\left(S \mathcal{V}_{L}\right)^{\prime}
\end{array}
$$

where $S \in \mathcal{L}\left(L^{2}(\partial \Omega)^{m_{1}}\right)$ is boundedly invertible, and $S_{+}$and $\left(S^{*}\right)_{-}^{-1}$ denote their extension on $\mathcal{V}_{L}$ and $\mathcal{V}_{L}^{\prime}$ respectively (see Corollary 5.26). By Lemma 5.28 also $G_{\mathrm{p}}$ and $K_{\mathrm{p}}$ establish a boundary triple for $L_{\mathrm{p}}$ restricted to $H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right) \times H\left(L_{\partial}, \Omega\right)$ and $\left(S_{+} \mathcal{V}_{L, \Gamma_{1}}, S L_{\pi}^{2}\left(\Gamma_{1}\right),\left(S_{+} \mathcal{V}_{L, \Gamma_{1}}\right)^{\prime}\right)$ is a quasi Gelfand triple For simplification $S$ can be imagined to be the identity mapping. We still have $\Gamma_{0}, \Gamma_{1}$ as a splitting with thin boundaries of $\partial \Omega$.
Corollary 8.4. The colligation $\left(\left[\begin{array}{c}G_{\mathrm{p}} \\ L_{\mathrm{p}} \\ K_{\mathrm{p}}\end{array}\right] ;\left[\begin{array}{c}S_{+} \mathcal{V}_{L, \Gamma_{1}} \\ \mathcal{X}_{\mathcal{H}} \\ \left(S_{+} \mathcal{V}_{L, \Gamma_{1}}\right)^{\prime}\end{array}\right]\right)$ with solution space

$$
\mathcal{Z}=\mathcal{H}^{-1}\left(H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right) \times H\left(L_{\partial}, \Omega\right)\right)
$$

is a boundary control system.
Proof. Since $L_{\mathrm{p}}$ is closed on $\mathcal{X}_{\mathcal{H}}$ with domain $\mathcal{Z}$, and $G_{\mathrm{p}}$ and $K_{\mathrm{p}}$ are continuous with the graph norm of $L_{\mathrm{p}}$, we have $\left[G_{\mathrm{p}} L_{\mathrm{p}} K_{\mathrm{p}}\right]^{\top}$ is closed. By construction $G_{\mathrm{p}}$ with domain $\mathcal{Z}$ maps onto $S_{+} \mathcal{V}_{L, \Gamma_{0}}$. Since $G_{\mathrm{p}}$ is one operator of a boundary triple for $L_{\mathrm{p}}$, the restriction $\left.L_{\mathrm{p}}\right|_{\operatorname{ker} G_{\mathrm{p}}}$ is skew-adjoint and therefore a generator of a contraction semigroup.
Proposition 8.5. Let $R \in \mathcal{L}\left(S L_{\pi}^{2}\left(\Gamma_{1}\right)\right)$ be coercive. Then the colligation $\Xi=$ $\left(\left[\begin{array}{c}\frac{1}{\sqrt{2}}\left(G_{\mathrm{p}}+R K_{\mathrm{p}}\right) \\ L_{\mathrm{p}} \\ \frac{1}{\sqrt{2}}\left(G_{\mathrm{p}}-R K_{\mathrm{p}}\right)\end{array}\right] ;\left[\begin{array}{c}\mathcal{U} \\ \mathcal{X}_{\mathcal{H}} \\ \mathcal{Y}\end{array}\right]\right)$ with $\mathcal{U}=\mathcal{Y}=S L_{\pi}^{2}\left(\Gamma_{1}\right)$ endowed with $\|f\|_{\mathcal{U}}=\|f\|_{\mathcal{Y}}=$
$\left\|R^{-1 / 2} f\right\|_{L^{2}}$ and solution space

$$
\mathcal{Z}=\left\{x \in \mathcal{H}^{-1}\left(H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right) \times H\left(L_{\partial}, \Omega\right)\right): G_{\mathrm{p}} x, K_{\mathrm{p}} x \in S L_{\pi}^{2}\left(\Gamma_{1}\right)\right\} .
$$

is a scattering energy preserving boundary control system
Proof. Let $\left(x_{n},\left[\begin{array}{lll}G_{\mathrm{p}} x_{n} & L_{\mathrm{p}} x_{n} & K_{\mathrm{p}} x_{n}\end{array}\right]^{\top}\right)_{n \in \mathbb{N}}$ be a sequence in $\left[\begin{array}{lll}G_{\mathrm{p}} & L_{\mathrm{p}} & K_{\mathrm{p}}\end{array}\right]^{\top}$ (restricted to $\mathcal{Z})$ that converges to $\left(x,\left[\begin{array}{lll}f & y & g\end{array}\right]^{\top}\right) \in \mathcal{X}_{\mathcal{H}} \times \mathcal{U} \times \mathcal{X}_{\mathcal{H}} \times \mathcal{U}$. Since $L_{\mathrm{p}}$ with domain $H\left(P_{\partial}, \Omega\right)$ is a closed operator and $H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right) \times H\left(L_{\partial}, \Omega\right)$ is closed in $H\left(P_{\partial}, \Omega\right)$, we conclude that $x \in \mathcal{H}^{-1}\left(H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right) \times H\left(L_{\partial}, \Omega\right)\right)$ and $y=L_{\mathrm{p}} x$. Hence, $G_{\mathrm{p}} x_{n}$ converges in $S_{+} \mathcal{V}_{L, \Gamma_{1}}$ to $G_{\mathrm{p}} x$ and in $S L_{\pi}^{2}\left(\Gamma_{1}\right)$ to $f$. Since $\left(S_{+} \mathcal{V}_{L, \Gamma_{1}}, S L_{\pi}^{2}\left(\Gamma_{1}\right)\right.$, $\left.\left(S_{+} \mathcal{V}_{L, \Gamma_{1}}\right)^{\prime}\right)$ is a quasi Gelfand triple, we have $G_{\mathrm{p}} x=f$. Analogously, we conclude $K_{\mathrm{p}} x=g$. Therefore, $x \in \mathcal{Z}$ and $\left[\begin{array}{lll}G_{\mathrm{p}} & L_{\mathrm{p}} & K_{\mathrm{p}}\end{array}\right]^{\top}$ is closed, which implies that also $\left[\frac{1}{\sqrt{2}}\left(G_{\mathrm{p}}+R K_{\mathrm{p}}\right) \quad L_{\mathrm{p}} \quad \frac{1}{\sqrt{2}}\left(G_{\mathrm{p}}-R K_{\mathrm{p}}\right)\right]^{\top}$ is closed.

By Example 7.8 and Theorem $\left.7.6 L_{\mathrm{p}}\right|_{\operatorname{ker} \frac{1}{\sqrt{2}}\left(G_{\mathrm{p}}+R K_{\mathrm{p}}\right)}$ generates a contraction semigroup.

The surjectivity of $\left[\begin{array}{c}G_{\mathrm{p}} \\ K_{\mathrm{p}}\end{array}\right]$ and Example 7.8 gives the surjectivity of $\frac{1}{\sqrt{2}}\left(G_{\mathrm{p}}+R K_{\mathrm{p}}\right)$. Since $\left(\mathcal{V}_{L}, G_{\mathrm{p}}, \Psi K_{\mathrm{p}}\right)$ is a boundary triple for $L_{\mathrm{p}}$, we have

$$
\begin{aligned}
2 \operatorname{Re}\left\langle L_{\mathrm{p}} x, x\right\rangle_{\mathcal{X}_{\mathcal{H}}}= & 2 \operatorname{Re}\left\langle G_{\mathrm{p}} x, K_{\mathrm{p}} x\right\rangle_{\mathcal{V}_{L}, \mathcal{V}_{L}^{\prime}}=2 \operatorname{Re}\left\langle G_{\mathrm{p}} x, K_{\mathrm{p}} x\right\rangle_{L_{\pi}^{2}\left(\Gamma_{1}\right)} \\
= & \frac{1}{2}\left(\left\langle R^{-1} G_{\mathrm{p}} x, G_{\mathrm{p}} x\right\rangle_{L^{2}}+2 \operatorname{Re}\left\langle G_{\mathrm{p}} x, K_{\mathrm{p}} x\right\rangle_{L^{2}}+\left\langle R K_{\mathrm{p}} x, K_{\mathrm{p}} x\right\rangle_{L^{2}}\right) \\
& -\frac{1}{2}\left(\left\langle R^{-1} G_{\mathrm{p}} x, G_{\mathrm{p}} x\right\rangle_{L^{2}}-2 \operatorname{Re}\left\langle G_{\mathrm{p}} x, K_{\mathrm{p}} x\right\rangle_{L^{2}}+\left\langle R K_{\mathrm{p}} x, K_{\mathrm{p}} x\right\rangle_{L^{2}}\right) \\
= & \left\|\frac{1}{\sqrt{2}}\left(G_{\mathrm{p}}+R K_{\mathrm{p}}\right) x\right\|_{\mathcal{U}}^{2}-\left\|\frac{1}{\sqrt{2}}\left(G_{\mathrm{p}}-R K_{\mathrm{p}}\right) x\right\|_{\mathcal{Y}}^{2}
\end{aligned}
$$

which makes $\Xi$ scattering energy preserving.
Remark 8.6. Clearly, the previous proposition holds also true for the operator triple $\left[\frac{1}{\sqrt{2}}\left(R K_{\mathrm{p}}+G_{\mathrm{p}}\right) \quad L_{\mathrm{p}} \quad \frac{1}{\sqrt{2}}\left(R K_{\mathrm{p}}-G_{\mathrm{p}}\right)\right]^{\top}$ and for $G_{\mathrm{p}}$ and $K_{\mathrm{p}}$ being swapped. Moreover, replacing $L_{\mathrm{p}}$ by $L_{\mathrm{p}}+J$, where $J \in \mathcal{L}\left(\mathcal{X}_{\mathcal{H}}\right)$ is dissipative, yields a scattering passive system.

Hence, the port-Hamiltonian system with input $u$ and output $y$ described by the equations

$$
\begin{align*}
\sqrt{2} u(t, \zeta) & =\pi_{L}(\mathcal{H}(\zeta) x(t, \zeta))_{L^{H}}+R L_{\nu}(\mathcal{H}(\zeta) x(t, \zeta))_{L}, & & t \in \mathbb{R}_{+}, \zeta \in \Gamma_{1} \\
\frac{\partial}{\partial t} x(t, \zeta) & =\sum_{i=1}^{n} \frac{\partial}{\partial \zeta_{i}} P_{i}(\mathcal{H}(\zeta) x(t, \zeta))+P_{0}(\mathcal{H}(\zeta) x(t, \zeta)), & & t \in \mathbb{R}_{+}, \zeta \in \Omega \\
\sqrt{2} y(t, \zeta) & =\pi_{L}(\mathcal{H}(\zeta) x(t, \zeta))_{L^{H}}-R L_{\nu}(\mathcal{H}(\zeta) x(t, \zeta))_{L}, & & t \in \mathbb{R}_{+}, \zeta \in \Gamma_{1}  \tag{8.3}\\
0 & =\pi_{L}(\mathcal{H}(\zeta) x(t, \zeta))_{L^{H}}, & & t \in \mathbb{R}_{+}, \zeta \in \Gamma_{0} \\
x(0, \zeta) & =x_{0}(\zeta), & & \zeta \in \Omega
\end{align*}
$$

is scattering passive and in particular well-posed, as the following corollary will clarify. The mappings $\pi_{L}$ and $L_{\nu}$ are used a little bit sloppy. There is always a pointwise a.e. description for these mappings, but due to compact notation we use $\pi_{L}$ and $L_{\nu}$.

Corollary 8.7. The system (8.3) can be interpreted as the scattering energy preserving boundary control system

$$
\left(\left[\begin{array}{c}
\frac{1}{\sqrt{2}}\left(G_{\mathrm{p}}+R K_{\mathrm{p}}\right) \\
L_{\mathrm{p}} \\
\frac{1}{\sqrt{2}}\left(G_{\mathrm{p}}-R K_{\mathrm{p}}\right)
\end{array}\right] ;\left[\begin{array}{c}
\mathcal{U} \\
\mathcal{X}_{\mathcal{H}} \\
\mathcal{Y}
\end{array}\right]\right)
$$

with the assumptions of Proposition 8.5 and $S=\mathrm{I}$. Replacing $L_{\mathrm{p}}$ with $L_{\mathrm{p}}+J$ for a dissipative $J \in \mathcal{L}\left(\mathcal{X}_{\mathcal{H}}\right)$ yields a scattering passive boundary control system.

Corollary 8.8. With the setting of Proposition 8.5 the colligation

$$
\left(\left[\begin{array}{c}
G_{\mathrm{p}} \\
L_{\mathrm{p}} \\
K_{\mathrm{p}}
\end{array}\right] ;\left[\begin{array}{c}
S L_{\pi}^{2}\left(\Gamma_{1}\right) \\
\mathcal{X} \mathcal{H} \\
S L_{\pi}^{2}\left(\Gamma_{1}\right)
\end{array}\right]\right)
$$

with solution space

$$
\mathcal{Z}=\left\{x \in \mathcal{H}^{-1}\left(H_{\Gamma_{0}}\left(L_{\partial}^{\mathrm{H}}, \Omega\right) \times H\left(L_{\partial}, \Omega\right)\right): G_{\mathrm{p}} x, K_{\mathrm{p}} x \in S L_{\pi}^{2}\left(\Gamma_{1}\right)\right\}
$$

is impedance energy preserving.
Proof. This is a direct consequence of Proposition 8.5 for $R=\mathrm{I}$.
Note that the colligations in Corollary 8.4 and Corollary 8.8 are the same but the solution spaces are slightly different. The colligation in Corollary 8.8 is in general not necessarily a boundary control system.
Example 8.9 (Wave equation). Let $\rho \in L^{\infty}(\Omega)$ be the mass density and $T \in$ $L^{\infty}(\Omega)^{n \times n}$ be the Young modulus, such that $\frac{1}{\rho} \in L^{\infty}(\Omega), T(\zeta)^{\mathrm{H}}=T(\zeta)$ and $T(\zeta) \geq$ $\delta$ I for a $\delta>0$ and almost every $\zeta \in \Omega$. Then the wave equation

$$
\frac{\partial^{2}}{\partial t^{2}} w(t, \xi)=\frac{1}{\rho(\xi)} \operatorname{div}(T(\xi) \operatorname{grad} w(t, \xi))
$$

can be formulated as a port-Hamiltonian system by choosing the state variable $x(t, \zeta)=\left[\begin{array}{c}\rho(\xi) \frac{\partial}{\partial t} w(t, \zeta) \\ \operatorname{grad} w(t, \zeta)\end{array}\right]$. Then the PDE looks like

$$
\dot{x}=\underbrace{\left[\begin{array}{cc}
0 & \operatorname{div} \\
\operatorname{grad} & 0
\end{array}\right]}_{=P_{\partial}} \underbrace{\left[\begin{array}{cc}
\frac{1}{\rho} & 0 \\
0 & T
\end{array}\right]}_{=\mathcal{H}} x .
$$

This is shown in section 3 of [9]. This is exactly the port-Hamiltonian system we get from choosing $L$ as in Example 3.3. From Example 6.2 and Example 6.6 we know that the boundary operators are $\gamma_{0}$ and the extension of $\nu \cdot \gamma_{0}$. Therefore,

$$
\begin{aligned}
\sqrt{2} u(t, \zeta) & =\nu \cdot(T(\zeta) \operatorname{grad} w(t, \zeta))+\frac{\partial}{\partial t} w(t, \zeta), & & t \in \mathbb{R}_{+}, \zeta \in \Gamma_{1} \\
\frac{\partial^{2}}{\partial t^{2}} w(t, \xi) & =\frac{1}{\rho(\xi)} \operatorname{div}(T(\xi) \operatorname{grad} w(t, \xi)), & & t \in \mathbb{R}_{+}, \zeta \in \Omega \\
\sqrt{2} y(t, \zeta) & =\nu \cdot(T(\zeta) \operatorname{grad} w(t, \zeta))-\frac{\partial}{\partial t} w(t, \zeta), & & t \in \mathbb{R}_{+}, \zeta \in \Gamma_{1} \\
0 & =\frac{\partial}{\partial t} w(t, \zeta), & & t \in \mathbb{R}_{+}, \zeta \in \Gamma_{0}
\end{aligned}
$$

can be modeled by a scattering passive and well-posed boundary control system, by Corollary 8.7.

Example 8.10 (Maxwell's equations). Let $\Omega \subseteq \mathbb{R}^{3}$ be as in Assumption 3.1 and $L=\left(L_{i}\right)_{i=1}^{3}$ be as in Example 3.4. In this example we have already showed $L_{\partial}=\operatorname{rot}$ and $L_{\nu} f=\nu \times f$. The corresponding differential operator for the port-Hamiltonian PDE is

$$
P_{\partial}=\left[\begin{array}{cc}
0 & L_{\partial} \\
L_{\partial}^{\mathrm{H}} & 0
\end{array}\right]=\left[\begin{array}{cc}
0 & \operatorname{rot} \\
-\operatorname{rot} & 0
\end{array}\right] .
$$

We write the state as $x=\left[\begin{array}{l}\mathbf{D} \\ \mathbf{B}\end{array}\right]$, where $\mathbf{D}, \mathbf{B} \in \mathbb{K}^{3}$. We also want to introduce the positive scalar functions $\epsilon, \mu, g$ and $r$ such that

$$
\epsilon, \frac{1}{\epsilon}, \mu, \frac{1}{\mu}, g \in L^{\infty}(\Omega) \quad \text { and } \quad r, \frac{1}{r} \in L^{\infty}\left(\Gamma_{1}\right)
$$

Furthermore, we define the Hamiltonian density by $\mathcal{H}(\zeta):=\left[\begin{array}{cc}\frac{1}{\epsilon(\zeta)} & 0 \\ 0 & \frac{1}{\mu(\zeta)}\end{array}\right]$, where each block is a $3 \times 3$ matrix. At last we define $\left[\begin{array}{l}\mathbf{E} \\ \mathbf{H}\end{array}\right]:=\mathcal{H}\left[\begin{array}{l}\mathbf{D} \\ \mathbf{B}\end{array}\right]$, so that we have the same notation as in [16].

The projection on $\overline{\text { ran } L_{\nu}}$ is given by $g \mapsto(\nu \times g) \times \nu$, therefore $\bar{\pi}_{L}$ is the extension of $g \mapsto\left(\nu \times \gamma_{0} g\right) \times \nu$ to $H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$. The mapping $\pi_{\tau}$ from [16] can be compared with $\bar{\pi}_{L}$ but is not exactly the same, since they have different domains and codomains. We have $\pi_{\tau}: H^{1}(\Omega)^{3} \rightarrow V_{\tau} \subseteq L^{2}(\partial \Omega)^{3}$ and $\bar{\pi}_{L}: H(\operatorname{rot}, \Omega) \rightarrow \mathcal{V}_{L}$ is its extension, if we change the norms in the domain and codomain of $\pi_{\tau}$. However, $\mathcal{V}_{L}$ cannot be embedded into $L^{2}(\partial \Omega)^{3}$.

Note that by Example A. 4 neither $\bar{\pi}_{L}$ nor $\bar{L}_{\nu}^{\Gamma_{1}}$ map even into $L_{\pi}^{2}\left(\Gamma_{1}\right)$, therefore it is really necessary to use a quasi Gelfand triple instead of an "ordinary" Gelfand triple.

The corresponding boundary control system is a model for Maxwell's equations in the following form

$$
\begin{aligned}
\sqrt{2} u(t, \zeta) & =r(\zeta) \nu(\zeta) \times \mathbf{H}(t, \zeta)+(\nu(\zeta) \times \mathbf{E}(t, \zeta)) \times \nu(\zeta), & & t \in \mathbb{R}_{+}, \zeta \in \Gamma_{1}, \\
\frac{\partial}{\partial t} \mathbf{D}(t, \zeta) & =\operatorname{rot} \mathbf{H}(t, \zeta)-g(\zeta) \mathbf{E}(t, \zeta), & & t \in \mathbb{R}_{+}, \zeta \in \Omega \\
\frac{\partial}{\partial t} \mathbf{B}(t, \zeta) & =-\operatorname{rot} \mathbf{E}(t, \zeta), & & t \in \mathbb{R}_{+}, \zeta \in \Omega, \\
\sqrt{2} y(t, \zeta) & =r(\zeta) \nu(\zeta) \times \mathbf{H}(t, \zeta)-(\nu(\zeta) \times \mathbf{E}(t, \zeta)) \times \nu(\zeta), & & t \in \mathbb{R}_{+}, \zeta \in \Gamma_{1}, \\
0 & =(\nu(\zeta) \times \mathbf{E}(t, \zeta)) \times \nu(\zeta), & & t \in \mathbb{R}_{+}, \zeta \in \Gamma_{0},
\end{aligned}
$$

and is scattering passive by Corollary 8.7 , where we set $J=\left[\begin{array}{cc}-g & 0 \\ 0 & 0\end{array}\right] \mathcal{H}$.
Note that, following the trick in [16, Proposition 6.1], Gauß's law $\operatorname{div} \mathbf{D}=\rho$ is satisfied by simply defining $\rho$ by this formula and Gauß's law for magnetism $\operatorname{div} \mathbf{B}=0$ is automatically satisfied, if the initial condition satisfies it. This can be seen, if we apply div on both sides of $\frac{\partial}{\partial t} \mu \mathbf{H}=-\operatorname{rot} \mathbf{E}$ and noting that $\operatorname{div} \mu \mathbf{H}=$ $\operatorname{div} \mathbf{B}$ is constant in time (div rot $=0$ ). This has to be understood in the sense of distributions. However, for classical solutions this can also be understood in the classical sense.

Example 8.11 (Mindlin plate). Let $\Omega \subseteq \mathbb{R}^{2}$ be as in Assumption 3.1. Let us consider the differential operator $P_{\partial}$ and the skew-symmetric matrix $P_{0}$ given by

$$
P_{\partial}:=\left[\begin{array}{ccc|ccccc}
0 & 0 & 0 & 0 & 0 & 0 & \partial_{1} & \partial_{2} \\
0 & 0 & 0 & \partial_{1} & 0 & \partial_{2} & 0 & 0 \\
0 & 0 & 0 & 0 & \partial_{2} & \partial_{1} & 0 & 0 \\
\hline 0 & \partial_{1} & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & \partial_{2} & 0 & 0 & 0 & 0 & 0 \\
0 & \partial_{2} & \partial_{1} & 0 & 0 & 0 & 0 & 0 \\
\partial_{1} & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
\partial_{2} & 0 & 0 & 0 & 0 & 0 & 0 & 0
\end{array}\right], P_{0}:=\left[\begin{array}{cccccccc}
0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 1 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & 1 \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & -1 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & -1 & 0 & 0 & 0 & 0 & 0
\end{array}\right] .
$$

It is easy to derive the corresponding $P=\left(P_{i}\right)_{i=1}^{2}$ and $L=\left(L_{i}\right)_{i=1}^{2}$. We define a Hamiltonian density by

$$
\mathcal{H}=\left[\begin{array}{cccccccc}
\frac{1}{\rho h} & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & \frac{12}{\rho h^{3}} & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & \frac{12}{\rho h^{3}} & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & & & 0 & 0 \\
0 & 0 & 0 & \boldsymbol{D}_{b} & 0 & 0 \\
0 & 0 & 0 & & & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & \boldsymbol{D}_{s} \\
0 & 0 & 0 & 0 & 0 & 0 &
\end{array}\right]
$$

where $\rho, h$ are strictly positive functions, $\boldsymbol{D}_{b}(\zeta)$ is a coercive $3 \times 3$ matrix and $\boldsymbol{D}_{s}(\zeta)$ is a coercive $2 \times 2$ matrix, such that all conditions on $\mathcal{H}$ in Definition 4.1 are satisfied. We write the state variable $x$ as

$$
\boldsymbol{\alpha}:=\left[\begin{array}{llllllll}
\rho h v & \rho \frac{h^{3}}{12} w_{1} & \rho \frac{h^{3}}{12} w_{2} & \kappa_{1,1} & \kappa_{2,2} & \kappa_{1,2} & \gamma_{1,3} & \gamma_{2,3}
\end{array}\right]^{\top}
$$

where we stick to the notation in [2] except that we renamed the coordinates $x, y$ and $z$ as 1,2 and 3 . Furthermore, we have

$$
\mathbf{e}:=\mathcal{H} \boldsymbol{\alpha}=\left[\begin{array}{llllllll}
v & w_{1} & w_{2} & M_{1,1} & M_{2,2} & M_{1,2} & Q_{1} & Q_{2}
\end{array}\right]^{\top}
$$

We don't want to go into details about the physical meaning of these state variables. We just want to make it easier to translate the results into the notation of [2]. So the port-Hamiltonian PDE

$$
\frac{\partial}{\partial t} x=\left(P_{\partial}+P_{0}\right) \mathcal{H} x \quad \text { looks like } \quad \frac{\partial}{\partial t} \boldsymbol{\alpha}=\left(P_{\partial}+P_{0}\right) \mathbf{e}
$$

The corresponding boundary operator is

$$
L_{\nu} f=\left[\begin{array}{ccccc}
0 & 0 & 0 & \nu_{1} & \nu_{2} \\
\nu_{1} & 0 & \nu_{2} & 0 & 0 \\
0 & \nu_{2} & \nu_{1} & 0 & 0
\end{array}\right]\left[\begin{array}{l}
f_{1} \\
f_{2} \\
f_{3} \\
f_{4} \\
f_{5}
\end{array}\right]=\left[\begin{array}{l}
\nu \cdot\left[\begin{array}{l}
f_{4} \\
f_{5}
\end{array}\right] \\
\nu \cdot\left[\begin{array}{l}
f_{1} \\
f_{3}
\end{array}\right] \\
\nu \cdot\left[\begin{array}{l}
f_{3} \\
f_{2}
\end{array}\right]
\end{array}\right]
$$

Since $\|\nu(\zeta)\|=1$, at least $\nu_{1}(\zeta) \neq 0$ or $\nu_{2}(\zeta) \neq 0$. This can be used to show that $\operatorname{ran} L_{\nu}=L^{2}(\partial \Omega)^{3}$. Therefore, $\bar{\pi}_{L}$ is the extension of the boundary trace operator $\gamma_{0}$ to $H\left(L_{\partial}^{\mathrm{H}}, \Omega\right)$.

Since there is no direct physical meaning to the boundary variables

$$
\left[\begin{array}{ll}
0 & L_{\nu}
\end{array}\right] \mathbf{e}=\left[\begin{array}{c}
\nu \cdot\left[\begin{array}{l}
Q_{1} \\
Q_{2}
\end{array}\right] \\
\nu \cdot\left[\begin{array}{l}
M_{1,1} \\
M_{1,2}
\end{array}\right] \\
\nu \cdot\left[\begin{array}{l}
M_{1,2} \\
M_{2,2}
\end{array}\right]
\end{array}\right] \quad \text { and } \quad\left[\begin{array}{ll}
\pi_{L} & 0
\end{array}\right] \mathbf{e}=\left[\begin{array}{c}
v \\
w_{1} \\
w_{2}
\end{array}\right]
$$

we define $\eta:=\left[\begin{array}{c}-\nu_{2} \\ \nu_{1}\end{array}\right]$ and apply the unitary transformation $S=\left[\begin{array}{ccc}1 & 0 & 0 \\ 0 & \nu_{1} & \nu_{2} \\ 0 & -\nu_{2} & \nu_{1}\end{array}\right]$ to obtain

$$
\left[\begin{array}{c}
Q_{\nu} \\
M_{\nu, \nu} \\
M_{\nu, \eta}
\end{array}\right]:=S\left[\begin{array}{l}
\nu \cdot\left[\begin{array}{l}
Q_{1} \\
Q_{2}
\end{array}\right] \\
\nu \cdot\left[\begin{array}{l}
M_{1,1} \\
M_{1,2}
\end{array}\right] \\
\nu \cdot\left[\begin{array}{l}
M_{1,2} \\
M_{2,2}
\end{array}\right]
\end{array}\right] \quad \text { and } \quad\left[\begin{array}{c}
v \\
w_{\nu} \\
w_{\eta}
\end{array}\right]:=\underbrace{\left(S^{*}\right)^{-1}}_{=S}\left[\begin{array}{c}
v \\
w_{1} \\
w_{2}
\end{array}\right],
$$

which have a physical interpretation; see [2]. Hence, by Corollary 8.8 the system

$$
\begin{aligned}
u & =\left[\begin{array}{lll}
Q_{\nu} & M_{\nu, \nu} & M_{\nu, \eta}
\end{array}\right]^{\top}, & & \text { on } \mathbb{R}_{+} \times \Gamma_{1}, \\
\frac{\partial}{\partial t} \boldsymbol{\alpha} & =\left(P_{\partial}+P_{0}\right) \mathbf{e}, & & \text { on } \mathbb{R}_{+} \times \Omega \\
y & =\left[\begin{array}{lll}
v & w_{\nu} & w_{\eta}
\end{array}\right]^{\top}, & & \text { on } \mathbb{R}_{+} \times \Gamma_{1}, \\
0 & =\left[\begin{array}{lll}
v & w_{\nu} & w_{\eta}
\end{array}\right]^{\top}, & & \text { on } \mathbb{R}_{+} \times \Gamma_{0},
\end{aligned}
$$

for the Mindlin plate is impedance energy preserving, which is exactly the system in [2].

Appendix A. Counter examples and technical lemmas. The next example shows that it is possible to have item (i) and item (ii) of a "boundary triple" for an operator $A$ (Definition 2.1) without $A$ being the adjoint of a skew-symmetric operator. Moreover, it shows that in this situation Lemma 2.2 does not hold. This demonstrates the importance of $A$ being the adjoint of a skew-symmetric operator in the definition.
Example A.1. Let $A=\left[\begin{array}{cc}0 & \frac{d}{d \xi} \\ \frac{d}{d \xi} & 0\end{array}\right]$ be an operator on $L^{2}(0,1)^{2}$ with $\operatorname{dom} A=$ $H^{1}(0,1)^{2}$. By Remark 3.7 the operator $A$ is the adjoint of a skew-symmetric operator. Integration by parts yields

$$
\begin{aligned}
\langle A f, g\rangle+\langle f, A g\rangle & =\int_{0}^{1}\left\langle\left[\begin{array}{l}
f_{2}^{\prime} \\
f_{1}^{\prime}
\end{array}\right],\left[\begin{array}{l}
g_{1} \\
g_{2}
\end{array}\right]\right\rangle \mathrm{d} \xi+\int_{0}^{1}\left\langle\left[\begin{array}{l}
f_{1} \\
f_{2}
\end{array}\right],\left[\begin{array}{l}
g_{2}^{\prime} \\
g_{1}^{\prime}
\end{array}\right]\right\rangle \mathrm{d} \xi \\
& =\int_{0}^{1}\left(f_{2}^{\prime} g_{1}+f_{1}^{\prime} g_{2}+f_{1} g_{2}^{\prime}+f_{2} g_{1}^{\prime}\right) \mathrm{d} \xi=\left.f_{2} g_{1}\right|_{0} ^{1}+\left.f_{1} g_{2}\right|_{0} ^{1} \\
& =f_{2}(1) g_{1}(1)-f_{2}(0) g_{1}(0)+f_{1}(1) g_{2}(1)-f_{1}(0) g_{2}(0) \\
& =\langle\underbrace{\left[\begin{array}{c}
f_{2}(1) \\
-f_{2}(0)
\end{array}\right]}_{B_{2} f}, \underbrace{\left[\begin{array}{c}
g_{1}(1) \\
g_{1}(0)
\end{array}\right]}_{B_{1} g}\rangle+\langle\underbrace{\left[\begin{array}{c}
f_{1}(1) \\
f_{1}(0)
\end{array}\right]}_{B_{1} f}, \underbrace{\left[\begin{array}{c}
g_{2}(1) \\
-g_{2}(0)
\end{array}\right]}_{B_{2} g} / .
\end{aligned}
$$

Defining $B_{1} f:=\left[\begin{array}{l}f_{1}(1) \\ f_{1}(0)\end{array}\right]$ and $B_{2} f:=\left[\begin{array}{c}f_{2}(1) \\ -f_{2}(0)\end{array}\right]$ yields

$$
\begin{equation*}
\langle A f, g\rangle+\langle f, A g\rangle=\left\langle B_{1} f, B_{2} g\right\rangle+\left\langle B_{2} f, B_{1} g\right\rangle \tag{A.1}
\end{equation*}
$$

The mapping $\left[\begin{array}{c}B_{1} \\ B_{2}\end{array}\right]: \operatorname{dom} A \rightarrow \mathbb{R}^{4}$ is surjective (this can be seen by choosing $f_{1}$ and $f_{2}$ to be linear interpolations). So $\left(\mathbb{R}^{2}, B_{1}, B_{2}\right)$ is a boundary triple for $A$.

We define $\hat{A}$ as the restriction of $A$ on $H_{\{1\}=0}^{1}(0,1) \times H_{\{0\}=\{1\}}^{1}(0,1)$, where

$$
\begin{aligned}
H_{\{1\}=0}^{1}(0,1) & :=\left\{f \in H^{1}(0,1): f(1)=0\right\}, \quad \text { and } \\
H_{\{0\}=\{1\}}^{1}(0,1) & :=\left\{f \in H^{1}(0,1): f(0)=f(1)\right\} .
\end{aligned}
$$

Therefore, we can reformulate (A.1) for $f, g \in \operatorname{dom} \hat{A}$

$$
\langle\hat{A} f, g\rangle+\langle f, \hat{A} g\rangle=-f_{1}(0) g_{2}(0)+f_{2}(0)\left(-g_{1}(0)\right)
$$

By defining $F_{1} f:=-f_{1}(0)$ and $F_{2} f:=f_{2}(0)$ we again have that $\left[\begin{array}{c}F_{1} \\ F_{2}\end{array}\right]: \operatorname{dom} \hat{A} \rightarrow \mathbb{R}^{2}$ is surjective. However $\hat{A}$ is not the adjoint of a skew-symmetric operator. If it were, then $\left(\mathbb{R}^{2}, F_{1}, F_{2}\right)$ would be a boundary triple for $\hat{A}$ and

$$
\hat{A}^{*}=-\left.\hat{A}\right|_{\operatorname{ker} F_{1} \cap \operatorname{ker} F_{2}}=-\left.A\right|_{H_{0}^{1}(0,1)^{2}}=A^{*}
$$

which is not true since $\hat{A}$ is certainly not dense in $A$. In fact, with the boundary triple for $A$ we get that the adjoint of $\hat{A}$ is $-\left.A\right|_{H_{\{0\}=\{1\}}^{1}(0,1) \times H_{\{0\}=0}^{1}(0,1)}$.

Lemma A.2. Let $\left(x_{n}\right)_{n \in \mathbb{N}}$ be a sequence in a normed vector space $X$ that converges w.r.t. the weak-* topology to an $x_{0} \in X$. Then $\left(x_{n}\right)_{n \in \mathbb{N}}$ is bounded i.e. $\sup _{n \in \mathbb{N}}\left\|x_{n}\right\|_{X}<+\infty$.

Proof. Let $\iota$ denote the canonical embedding from $X$ into $X^{\prime \prime}$ that maps $x$ to $\langle x, .\rangle_{X, X^{\prime}}$. Then, by assumption, for every fixed $\phi \in X^{\prime}\left(\iota x_{n}\right)(\phi) \rightarrow\left(\iota x_{0}\right)(\phi)$, in particular $\sup _{n \in \mathbb{N}}\left|\left(\iota x_{n}\right)(\phi)\right|<\infty$. The principle of uniform boundedness yields $\sup _{n \in \mathbb{N}}\left\|\iota x_{n}\right\|_{X^{\prime \prime}}<+\infty$. Since $\|\iota x\|_{X^{\prime \prime}}=\|x\|_{X}$ for every $x \in X$, this proves the assertion.

Lemma A.3. Let $\left(x_{n}\right)_{n \in \mathbb{N}}$ be a weak convergent sequence in a Hilbert space $H$ with limit $x$. Then there exists a subsequence $\left(x_{n(k)}\right)_{k \in \mathbb{N}}$ such that

$$
\left\|\frac{1}{N} \sum_{k=1}^{N} x_{n(k)}-x\right\| \rightarrow 0
$$

Proof. We assume that $x=0$. For the general result we just need to replace $x_{n}$ by $x_{n}-x$.

We define the subsequence inductively: $n(1)=1$ and for $k>1$ we choose $n(k)$ such that

$$
\left|\left\langle x_{n(k)}, x_{n(j)}\right\rangle\right| \leq \frac{1}{k} \quad \text { for all } \quad j<k
$$

This is possible, because $\left(x_{n}\right)_{n \in \mathbb{N}}$ converges weakly to 0 . Note that in Hilbert spaces the weak topology and the weak-* topology are the same. Hence, by Lemma A. 2
$\sup _{n \in \mathbb{N}}\left\|x_{n}\right\| \leq C$. This yields

$$
\begin{aligned}
\left\|\frac{1}{N} \sum_{k=1}^{N} x_{n(k)}\right\|^{2} & =\frac{1}{N^{2}} \sum_{k=1}^{N} \sum_{j=1}^{N}\left\langle x_{n(k)}, x_{n(j)}\right\rangle \\
& =\frac{1}{N^{2}} \sum_{k=1}^{N}\left\|x_{n(k)}\right\|^{2}+\frac{1}{N^{2}} \sum_{j=1}^{N} \sum_{k=j+1}^{N} 2 \operatorname{Re}\left\langle x_{n(k)}, x_{n(j)}\right\rangle \\
& \leq \frac{1}{N} C^{2}+\frac{2}{N^{2}} \sum_{j=1}^{N} \sum_{k=j+1}^{N} \frac{1}{k} \leq \frac{C^{2}}{N}+\frac{1}{N} \ln (N) \rightarrow 0
\end{aligned}
$$

Example A.4. Let $\Omega=(0,1)^{3}$ and $F: \Omega \rightarrow \mathbb{R}$ be defined by

$$
F(x)=\frac{1}{\|x\|_{2}^{2 / 5}}=\left(x_{1}^{2}+x_{2}^{2}+x_{3}^{2}\right)^{-2 / 10}
$$

Then we define $f=\operatorname{grad} F$, which is

$$
f(x)=\left[\begin{array}{c}
-\frac{4}{10} x_{1}\left(x_{1}^{2}+x_{2}^{2}+x_{3}^{2}\right)^{-6 / 5} \\
-\frac{4}{10} x_{2}\left(x_{1}^{2}+x_{2}^{2}+x_{3}^{2}\right)^{-6 / 5} \\
-\frac{4}{10} x_{3}\left(x_{1}^{2}+x_{2}^{2}+x_{3}^{2}\right)^{-6 / 5}
\end{array}\right]
$$

Hence, $\operatorname{rot} f=\operatorname{rot} \operatorname{grad} F=0$. We will show that $f$ is in $L^{2}(\Omega)^{3}$ :

$$
\begin{aligned}
\int_{\Omega}\|f(x)\|_{2}^{2} \mathrm{~d} x & =\int_{\Omega} \sum_{i=1}^{3} \frac{16}{100} x_{i}^{2}\left(x_{1}^{2}+x_{2}^{2}+x_{3}^{2}\right)^{-12 / 5} \mathrm{~d} x=\frac{16}{100} \int_{\Omega}\left(x_{1}^{2}+x_{2}^{2}+x_{3}^{2}\right)^{-7 / 5} \mathrm{~d} x \\
& \leq \int_{B_{\sqrt{3}}(0)}\left(x_{1}^{2}+x_{2}^{2}+x_{3}^{2}\right)^{-7 / 5} \mathrm{~d} x=2 \pi \int_{-\frac{\pi}{2}}^{\frac{\pi}{2}} \int_{0}^{\sqrt{3}} r^{-14 / 5} r^{2} \cos \theta \mathrm{~d} r \mathrm{~d} \theta \\
& =4 \pi \int_{0}^{\sqrt{3}} r^{-4 / 5} \mathrm{~d} r=\left.4 \pi 5 r^{1 / 5}\right|_{0} ^{\sqrt{3}}<+\infty .
\end{aligned}
$$

Therefore, $f$ is even in $H(\operatorname{rot}, \Omega)$. Let $\nu$ denote the normal vector on $\partial \Omega$. Then we show that $\nu \times\left. f\right|_{\partial \Omega}$ is not in $L^{2}(\partial \Omega)^{3}$ : Note that $\nu(\zeta)=\left[\begin{array}{c}0 \\ 0 \\ -1\end{array}\right]$ on $[0,1] \times[0,1] \times\{0\}$. Therefore,

$$
\nu(\zeta) \times f(\zeta)=\left[\begin{array}{c}
-\frac{4}{10} \zeta_{2}\left(\zeta_{1}^{2}+\zeta_{2}^{2}\right)^{-6 / 5} \\
\frac{4}{10} \zeta_{1}\left(\zeta_{1}^{2}+\zeta_{2}^{2}\right)^{-6 / 5} \\
0
\end{array}\right] \quad \text { for } \quad \zeta \in[0,1] \times[0,1] \times\{0\}
$$

and consequently

$$
\begin{aligned}
\int_{\partial \Omega}\|\nu(\zeta) \times f(\zeta)\|_{2}^{2} \mathrm{~d} \zeta & \geq \int_{[0,1] \times[0,1] \times\{0\}}\|\nu(\zeta) \times f(\zeta)\|_{2}^{2} \mathrm{~d} \zeta \\
& =\frac{16}{100} \int_{[0,1] \times[0,1]}\left(\xi_{1}^{2}+\xi_{2}^{2}\right)^{-7 / 5} \mathrm{~d} \xi
\end{aligned}
$$

Since $[0,1] \times[0,1]$ contains the circular sector with arc $\frac{\pi}{2}$ and radius 1 , we further have (by applying polar coordinates)

$$
\begin{aligned}
& \geq \frac{16}{100} \frac{\pi}{2} \int_{0}^{1} r^{-14 / 5} r \mathrm{~d} r=\frac{16}{100} \frac{\pi}{2} \int_{0}^{1} r^{-9 / 5} \mathrm{~d} r \\
& =-\left.\frac{16}{100} \frac{\pi}{2} \frac{5}{4} r^{-4 / 5}\right|_{0} ^{1}=+\infty
\end{aligned}
$$

Hence, $f \in H($ rot, $\Omega)$, but $\nu \times\left. f\right|_{\partial \Omega} \notin L^{2}(\partial \Omega)^{3}$. Since

$$
(\nu(\zeta) \times f(\zeta)) \times \nu(\zeta)=\left[\begin{array}{c}
-\frac{4}{10} \zeta_{1}\left(\zeta_{1}^{2}+\zeta_{2}^{2}\right)^{-6 / 5} \\
-\frac{4}{10} \zeta_{2}\left(\zeta_{1}^{2}+\zeta_{2}^{2}\right)^{-6 / 5} \\
0
\end{array}\right] \quad \text { for } \quad \zeta \in[0,1] \times[0,1] \times\{0\}
$$

we also have $\left(\nu \times\left. f\right|_{\partial \Omega}\right) \times \nu \notin L^{2}(\partial \Omega)^{3}$.
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