

# BASIC FUNCTIONAL ANALYSIS FOR THE OPTIMIZATION OF PARTIAL DIFFERENTIAL EQUATIONS

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ABSTRACT. Infinite-dimensional optimization requires – among other things – many results from functional analysis. In this script basics from functional analytic theory is reviewed. The purpose of this work is to give a summary of important facts needed to work in our research group.

## 1. Functional Analysis – Results and Definitions

If  $M$  is a set and  $M_1 \subset M$ , the symbol  $M \setminus M_1$  represents the *complement* of  $M_1$  in  $M$ , i.e.  $M \setminus M_1 = \{x \in M : x \notin M_1\}$ .  $\overline{M}$  will always denote the *closure* of the set  $M$ , which is the smallest closed set containing in  $M$ . The *interior* of the set  $M$ ,  $M^\circ$ , is the largest open set containing in  $M$ . The *boundary* of  $M$  is the set  $\partial M = \overline{M} \setminus M^\circ$ . The set of ordered pairs  $\{(x, y) : x \in M_1, y \in M_2\}$  is called the *Cartesian product* of the sets  $M_1$  and  $M_2$  and it is denoted  $M_1 \times M_2$ .

Let  $f : M \rightarrow M_1$  be a function (or mapping).  $f(M)$  will usually called the *range* of  $f$  and will denoted  $\text{ran}(f)$ . The set  $\{x \in M : f(x) = 0\}$  is said to be the *kernel* of  $f$  and is denoted  $\ker(f)$ . A function  $f$  will be called *injective* if for each  $y \in \text{ran}(f)$  there is at most one  $x \in M$  such that  $f(x) = y$ ;  $f$  is called *surjective* if  $\text{ran}(f) = M_1$ . If  $f$  is both injective and surjective, we will say it is *bijective*. Let  $f : M \rightarrow M_1$  and  $g : M_1 \rightarrow M_2$  be two functions. The *composite mapping*  $r = g \circ f$  is defined by  $r : M \rightarrow M_2, x \mapsto r(x) = g(f(x))$ .

**Definition 1.1.** A (real) linear space is a set,  $V$ , over  $\mathbb{R}$ , whose elements satisfy the following properties

- 1)  $v + w = w + v$  for all  $v, w \in V$ ,
- 2)  $v + (w + u) = (v + w) + u$  for all  $v, w, u \in V$ ,
- 3) There is in  $V$  a unique element, denoted by  $0$  and called the zero element, such that  $v + 0 = v$  for each  $v$ ,
- 4) To each  $v$  in  $V$  corresponds a unique element, denoted by  $-v$ , such that  $v + (-v) = 0$ ,
- 5)  $\alpha(v + w) = \alpha v + \alpha w$  for all  $v, w \in V$  and  $\alpha \in \mathbb{R}$ ,
- 6)  $(\alpha + \beta)v = \alpha v + \beta v$  for all  $v \in V$  and  $\alpha, \beta \in \mathbb{R}$ ,
- 7)  $\alpha(\beta v) = (\alpha\beta)v$  for all  $v \in V$  and  $\alpha, \beta \in \mathbb{R}$ ,
- 8)  $1 \cdot v = v$  for all  $v \in V$ ,
- 9)  $0 \cdot v = 0$  for all  $v \in V$ .

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A (real) normed linear space is a linear space,  $V$ , over  $\mathbb{R}$  and a function,  $\|\cdot\|_V$ , from  $V$  to  $\mathbb{R}$  which satisfies:

- 1)  $\|v\|_V \geq 0$  for all  $v$  in  $V$ ,
- 2)  $\|v\|_V = 0$  if and only if  $v = 0$ ,
- 3)  $\|\alpha v\|_V = |\alpha| \|v\|_V$  for all  $v$  in  $V$  and  $\alpha$  in  $\mathbb{R}$ ,
- 4)  $\|v + w\|_V \leq \|v\|_V + \|w\|_V$  for all  $v$  and  $w$  in  $V$ .

A system of sets  $M_\alpha$ ,  $\alpha \in I$ , is called a covering of the set  $M$  if  $M$  is contained as a subset of the union  $\bigcup_{\alpha \in I} M_\alpha$ . A subset  $M$  of a linear space  $V$  is called compact if every system of open sets of  $V$  which covers  $M$  contains a finite system also covering  $M$ . A subset  $M$  in a linear space  $V$  is precompact, if  $\overline{M}$  is compact in  $V$ . Further we call  $M \subset V$  bounded, if there exists a constant  $K > 0$  such that  $\|v\|_V \leq K$  for all  $v \in M$ .

**Definition 1.2.** A linear operator from a normed linear space  $(V_1, \|\cdot\|_{V_1})$  to a normed linear space  $(V_2, \|\cdot\|_{V_2})$  is a mapping,  $\mathcal{A}$ , from  $V_1$  to  $V_2$  which has the following property:

$$\mathcal{A}(\alpha v + \beta w) = \alpha \mathcal{A}(v) + \beta \mathcal{A}(w) \quad \text{for all } v, w \in V_1 \text{ and } \alpha, \beta \in \mathbb{R}.$$

$\mathcal{A}$  is called a bounded (linear) operator if  $\mathcal{A}$  is linear and there is some  $K > 0$  such that  $\|\mathcal{A}(v)\|_{V_2} \leq K \|v\|_{V_1}$  for all  $v \in V_1$ .

The smallest such  $K$  is called the norm of  $\mathcal{A}$ . By (1.1) we will introduce a notation for the norm of a bounded linear operator. A sequence of elements  $\{v_n\}_{n \in \mathbb{N}}$  of a normed linear space  $V$  is said to converge (strongly) to an element  $v \in V$ ,  $v_n \rightarrow v$ ,  $n \rightarrow \infty$ , if  $\lim_{n \rightarrow \infty} \|v - v_n\|_V = 0$ . The sequence  $\{v_n\}_{n \in \mathbb{N}}$  is called a Cauchy sequence if for all  $\varepsilon > 0$  there exists one  $N \in \mathbb{N}$  such that

$$\|v_m - v_n\|_V < \varepsilon \quad \text{for all } m, n \geq N.$$

A normed linear space in which all Cauchy sequences converge is called complete. A set  $M$  in a normed linear space  $V$  is called dense if every  $v \in V$  is a limit of elements in  $M$ . A function  $f$  from a normed linear space  $(V_1, \|\cdot\|_{V_1})$  to a normed space  $(V_2, \|\cdot\|_{V_2})$  is called continuous at  $v$  if  $\|f(v_n) - f(v)\|_{V_2} \rightarrow 0$  as  $n \rightarrow \infty$  whenever  $\|v_n - v\|_{V_1} \rightarrow 0$  as  $n$  tends to zero. We say  $f$  is Lipschitz-continuous if there exists a constant  $\gamma_f > 0$  such that

$$\|f(v) - f(w)\|_{V_2} \leq \gamma_f \|v - w\|_{V_1} \quad \text{for all } v, w \in V_1.$$

$f$  is called locally Lipschitz continuous if for all open and bounded  $O \subset V_1$  with  $\overline{O} \subset V_1$  there exists  $\gamma_f = \gamma_f(O) > 0$  such that

$$\|f(v) - f(w)\|_{V_2} \leq \gamma_f \|v - w\|_{V_1} \quad \text{for all } v, w \in O.$$

Let  $V_1$  and  $V_2$  be normed linear spaces. A bijection  $f$  from  $V_1$  to  $V_2$  which preserves the norm, i.e.,

$$\|f(v) - f(w)\|_{V_2} = \|v - w\|_{V_1} \quad \text{for all } v, w \in V_1$$

is called an isometry. It is automatically continuous.  $V_1$  and  $V_2$  are said to be isometric if such an isometry exists.

Let  $V$  be a normed linear space. The set  $\{w \in V : \|v - w\|_V < \rho\}$  is called the open ball,  $B(v; \rho)$ , of radius  $\rho$  about the point  $v$ . A set  $U(v) \subset X$  is called a neighborhood of  $v \in U(v)$  if  $B(v; \rho) \subset U(v)$  for some  $\rho > 0$ . Let  $M \subset V$ . A point  $v$  is called a limit point of  $M$ , if for all  $\rho > 0$   $B(v; \rho) \cap (M \setminus \{v\}) \neq \emptyset$ , i.e.,  $x$  is a limit point of  $M$  if  $M$  contains points other than  $v$  arbitrarily near  $v$ .

**Lemma 1.3.** *Let  $\mathcal{A}$  be linear operator between two normed linear spaces. The following properties are equivalent:*

- 1)  $\mathcal{A}$  is continuous at one point.
- 2)  $\mathcal{A}$  is continuous at all points.
- 3)  $\mathcal{A}$  is bounded.

*Proof.* For the proof we refer the reader to Theorem 6.1.1 on page 97 in [15].  $\square$

$V_1$  and  $V_2$  are normed linear spaces. We define

$$L(V_1, V_2) = \{\mathcal{A} : V_1 \rightarrow V_2, \mathcal{A} \text{ is linear and continuous}\}.$$

Due to Lemma 1.3 linear operators in  $L(V_1, V_2)$  are bounded linear operators, which we also call *continuous operators*. Let us introduce the following norm on  $L(V_1, V_2)$ :

$$(1.1) \quad \|\mathcal{A}\|_{L(V_1, V_2)} = \sup_{\|v\|_{V_1} \leq 1} \|\mathcal{A}(v)\|_{V_2} = \sup_{\|v\|_{V_1} = 1} \|\mathcal{A}(v)\|_{V_2} \quad \text{for all } \mathcal{A} \in L(V_1, V_2).$$

**Definition 1.4.** *A complete normed linear space is called a Banach space.*

We mention the *inverse mapping theorem*.

**Theorem 1.5.** *A continuous bijection of one Banach space onto another has a continuous inverse.*

*Proof.* We refer the reader to Theorem III.11 on page 83 in [13].  $\square$

**Lemma 1.6.** *Let  $B_1, B_2$  be two Banach spaces and  $\mathcal{A}$  belong to  $L(B_1, B_2)$ . The passage to the inverse  $\mathcal{A} \rightarrow \mathcal{A}^{-1}$  is continuous (non-linear) mapping of  $L(B_1, B_2)$  into  $L(B_2, B_1)$  for the norm.*

*Proof.* For a proof we refer to Theorem 3 on page 321 in [5].  $\square$

**Remark 1.7.** *If the perturbation  $\mathcal{B}$  of  $\mathcal{A}$  is sufficiently small, i.e.,  $\|\mathcal{B}\|_{L(B_1, B_2)} < \|\mathcal{A}\|_{L(B_1, B_2)}^{-1}$  holds, then  $\mathcal{A} - \mathcal{B}$  is invertible. Let  $\{\mathcal{A}_n\}_{n \in \mathbb{N}}$  be a sequence in  $L(B_1, B_2)$  and  $\mathcal{A} \in L(B_1, B_2)$  such that  $\mathcal{A}^{-1}$  exists and*

$$\lim_{n \rightarrow \infty} \|\mathcal{A}_n - \mathcal{A}\|_{L(B_1, B_2)} = 0.$$

*Thus, there exists  $N \in \mathbb{N}$  with  $\|\mathcal{A}_n - \mathcal{A}\|_{L(B_1, B_2)} < \|\mathcal{A}\|_{L(B_1, B_2)}^{-1}$  for all  $n \geq N$ . This leads to  $\mathcal{A}_n^{-1}$  exists for all  $n \geq N$ .*

**Lemma 1.8.** *Let  $V, V_1$  and  $V_2$  be normed linear spaces. Then  $L(V_1, V_2)$  with the norm  $\|\cdot\|_{L(V_1, V_2)}$  is a normed linear space, and a Banach space, if  $V_2$  is complete. If  $\mathcal{A} \in L(V_1, V_2)$  and  $\mathcal{B} \in L(V_2, V)$  we have  $\mathcal{B} \circ \mathcal{A} \in L(V_1, V)$  and*

$$\|\mathcal{B} \circ \mathcal{A}\|_{L(V_1, V)} \leq \|\mathcal{B}\|_{L(V_2, V)} \|\mathcal{A}\|_{L(V_1, V_2)}.$$

*Proof.* Let us refer to Satz 3.3 on page 102 in [2].  $\square$

We set  $L(V) = L(V, V)$ . The *identity* on  $V$  is the continuous operator  $\mathcal{I}_V : V \rightarrow V$  given by  $\mathcal{I}_V(v) = v$  for all  $v$  in  $V$ .

**Definition 1.9.** *A (real) vector space  $X$  is called (real) inner product space if there is a real-valued function  $\langle \cdot, \cdot \rangle_X$  on  $X \times X$  that satisfies the following four conditions for all  $x, y, z \in X$  and  $\alpha \in \mathbb{R}$ :*

- 1)  $\langle x, x \rangle_X \geq 0$  and  $\langle x, x \rangle_X = 0$  if and only if  $x = 0$ ,
- 2)  $\langle x, y + z \rangle_X = \langle x, y \rangle_X + \langle x, z \rangle_X$ ,

$$3) \langle x, \alpha y \rangle_X = \alpha \langle x, y \rangle_X,$$

$$4) \langle x, y \rangle_X = \langle y, x \rangle_X.$$

The function  $\langle \cdot, \cdot \rangle_X : X \times X \rightarrow \mathbb{R}$  is called (real) inner product.

**Example 1.10.** Let  $\mathbb{R}^n$  denote the set of all  $n$ -tuples of real numbers. We define the inner product

$$\langle \mathbf{x}, \mathbf{y} \rangle_{\mathbb{R}^n} = \mathbf{x}^T \mathbf{y} = \sum_{j=1}^n x_j y_j.$$

for all  $\mathbf{x} = (x_1, \dots, x_n)^T$  and  $\mathbf{y} = (y_1, \dots, y_n)^T$  in  $\mathbb{R}^n$ .

Let  $X$  and  $Y$  be inner product spaces. The mapping  $a : X \times X \rightarrow Y$  with the properties

$$1) a(\alpha x + \beta y, z) = \alpha a(x, z) + \beta a(y, z) \text{ for all } x, y, z \in X \text{ and } \alpha, \beta \in \mathbb{R},$$

$$2) a(x, y) = a(y, x) \text{ for all } x, y \in X,$$

$$3) |a(x, y)| \leq K \|x\|_X \|y\|_X \text{ for some } K > 0 \text{ and for all } x, y \in X$$

is said to be a (real) continuous bilinear form. Two vectors,  $x$  and  $y$ , in an inner product space  $X$  are said to be orthogonal if  $\langle x, y \rangle_X = 0$ . A collection  $\{x_i\}_{i \in \mathbb{N}}$  of vectors in  $X$  is called an orthonormal set if  $\langle x_i, x_i \rangle_X = 1$  for all  $i$ , and  $\langle x_i, x_j \rangle_X = 0$  if  $i \neq j$ .

**Definition 1.11.** A family  $\{x_\lambda\}_{\lambda \in \Lambda}$  ( $\Lambda$  an index set) is said to be total (or complete) in the Hilbert space  $X$  if

$$\langle x, x_\lambda \rangle_X = 0 \text{ for all } \lambda \in \Lambda \implies x = 0.$$

A total orthonormal family is called an orthonormal base.

**Lemma 1.12.** Every inner product space  $X$  is a normed linear space with the norm  $\|x\|_X = \sqrt{\langle x, x \rangle_X}$ .

*Proof.* Let us refer to Theorem II.2 on page 38 in [13].  $\square$

**Definition 1.13.** A complete (real) inner product space is called a (real) Hilbert space.

Suppose  $X$  and  $Y$  are Hilbert spaces. Then the set of pairs  $(x, y)$  with  $x \in X$ ,  $y \in Y$  is a Hilbert space called the direct sum of the spaces  $X$  and  $Y$  and denoted by  $X \oplus Y$ . The natural inner product on  $X \oplus Y$  is given by

$$(1.2) \quad \langle (x_1, y_1), (x_2, y_2) \rangle_{X \oplus Y} = \langle x_1, x_2 \rangle_X + \langle y_1, y_2 \rangle_Y$$

for all  $(x_1, y_1), (x_2, y_2) \in X \oplus Y$ .

Let  $M$  be a closed subspace of a given Hilbert space  $X$ . Under the natural inner product that it inherits as a subspace of  $X$ ,  $M$  is a Hilbert space. We denote by  $M^\perp$  the set of vectors in  $X$  which are orthogonal to  $M$ ;  $M^\perp$  is called the orthogonal complement of  $M$ . It follows from the linearity of the inner product that  $M^\perp$  is a linear subspace of  $X$ . Further, we can prove, that  $M^\perp$  is closed. Thus  $M^\perp$  is also a Hilbert space.  $M$  and  $M^\perp$  have only the zero element in common. The next Theorem 1.14 is usually called the projection theorem.

**Theorem 1.14.** Let  $X$  be a Hilbert space,  $M$  a closed subspace. Then every  $x \in X$  can be uniquely written  $x = z + w$  where  $z \in M$  and  $w \in M^\perp$ .

*Proof.* We refer the reader to Theorem II.3 on page 42 in [13].  $\square$

**Remark 1.15.** *Theorem 1.14 sets up a natural isomorphism between  $M \oplus M^\perp$  and  $X$  given by*

$$(1.3) \quad (z, w) \mapsto z + w.$$

*We will suppress the isomorphism and simply write  $X = M \oplus M^\perp$ . Let us choose  $x_1, x_2 \in X$ ,  $z_1, z_2 \in M$  and  $w_1, w_2 \in M^\perp$  such that  $x_1 = z_1 + w_1$  and  $x_2 = z_2 + w_2$ . By (1.3) we identify  $x_1$  and  $(z_1, w_1)$  respectively  $x_2$  and  $(z_2, w_2)$ . On the Hilbert spaces  $M$  and  $M^\perp$  we use the inner product  $\langle \cdot, \cdot \rangle_X$  that  $M$  and  $M^\perp$  inherit as a subspace of  $X$ . Then it follows*

$$\begin{aligned} \langle x_1, x_2 \rangle_X &= \langle z_1 + w_1, z_2 + w_2 \rangle_X \\ &= \langle z_1, z_2 \rangle_X + \underbrace{\langle w_1, z_2 \rangle_X + \langle z_1, w_2 \rangle_X}_{=0} + \langle w_1, w_2 \rangle_X \\ &\stackrel{(1.2)}{=} \langle (z_1, w_1), (z_2, w_2) \rangle_{M \oplus M^\perp}. \end{aligned}$$

*Thus, the inner product on  $M \oplus M^\perp$  given by (1.2) coincides with  $\langle \cdot, \cdot \rangle_X$ .*

An important class of bounded operators on Hilbert spaces is that of the projections.

**Definition 1.16.** *Let  $X$  be a Hilbert space. A bounded operator  $\mathcal{P}$  into itself is said to be a projection if  $\mathcal{P}^2 \equiv \mathcal{P}$  holds.  $\mathcal{P}$  is called orthogonal if  $\langle x - \mathcal{P}(x), \mathcal{P}(x) \rangle_X = 0$  for all  $x \in X$ .*

The following result is known as the principle of uniform boundedness or the Banach-Steinhaus theorem.

**Theorem 1.17.** *Let  $\{\mathcal{A}\}$  be a set in  $L(B_1, B_2)$  for two Banach spaces  $B_1$  and  $B_2$ . If  $\|\mathcal{A}(x)\|_{B_2}$  is bounded for each fixed  $x \in B_1$ , as  $\mathcal{A}$  ranges over  $\{\mathcal{A}\}$ , then there exists  $K > 0$  such that  $\|\mathcal{A}\|_{L(B_1, B_2)} \leq K$  for all of  $\{\mathcal{A}\}$ .*

*Proof.* Let us refer to Theorem 6.3.1 on page 112 in [15]. □

Now we introduce the dual space of a given Banach space.

**Definition 1.18.** *Let  $B$  be a Banach space. The space  $L(B, \mathbb{R})$  is called the dual space of  $B$  and it is denoted by  $B'$ . The elements of  $B'$  are continuous linear functionals. We write  $\langle f, x \rangle_{B', B}$  for the duality pairing of  $f \in B'$  with an element  $x \in B$ .*

Since  $\mathbb{R}$  is complete,  $B'$  is a Banach space (Lemma 1.8). We define  $\|\cdot\|_{B'} = \|\cdot\|_{L(B, \mathbb{R})}$ . In the applications we will often consider the dual  $X'$  of a Hilbert space  $X$ . Let us recall the Riesz representation theorem. For a proof we refer to Theorem II.4 on page 43 in [13].

**Theorem 1.19.** *Let  $X$  be a Hilbert space with dual  $X'$ . For each  $f \in X'$ , there is a unique  $y_f \in X$  such that  $\langle f, x \rangle_{X', X} = \langle y_f, x \rangle_X$  for all  $x \in X$ . In addition  $\|y_f\|_X = \|f\|_{X'}$ .*

A bounded linear operator from a normed linear space  $V_1$  to a normed linear space  $V_2$  is called an *isomorphism* if it is bijective and continuous and if it possesses a continuous inverse. If it preserves the norm, it is called *isometric isomorphism*. Due to Theorem 1.5 a bounded linear operator possesses a continuous inverse if it is bijective and both  $V_1$  and  $V_2$  are Banach spaces. Obviously, an isometric isomorphism has norm one.

**Remark 1.20.** By Theorem 1.19 we define the Riesz isomorphism  $\mathcal{J}_X$  which maps the Hilbert space  $X$  onto its dual  $X'$  by  $y_f \mapsto f$  and

$$\langle f, x \rangle_{X', X} = \langle \mathcal{J}_X(y_f), x \rangle_{X', X} = \langle y_f, x \rangle_X$$

for all  $x \in X$ . Often there is made no difference between an element  $y_f \in X$  and the corresponding element (the Riesz representant)  $f \in X'$ . We point out that we have  $\|\mathcal{J}_X\|_{L(X, X')} = 1$ . Thus,  $\mathcal{J}_X$  is an isometric isomorphism. The dual space  $X'$  is also a Hilbert space: Due to Theorem 1.19 the natural inner product on  $X'$  is given by

$$(1.4) \quad \langle f, g \rangle_{X'} = \langle \mathcal{J}_X^{-1}(f), \mathcal{J}_X^{-1}(g) \rangle_X$$

for all  $f, g \in X'$ .

If the normed linear space  $V$  has a countable dense subset it is called to be separable.

**Proposition 1.21.** Let  $X$  and  $Y$  be Hilbert spaces. We define

$$\begin{aligned} \mathcal{A} : X' \oplus Y' &\rightarrow (X \oplus Y)' \\ \langle \mathcal{A}(f, g), (x, y) \rangle_{(X \oplus Y)', X \oplus Y} &= \langle f, x \rangle_{X', X} + \langle g, y \rangle_{Y', Y} \end{aligned}$$

for all  $(f, g) \in X' \oplus Y'$  and  $(x, y) \in X \oplus Y$ . Then  $\mathcal{A}$  is an isometric isomorphism.

*Proof.* We have mentioned that  $X \oplus Y$  is a Hilbert space with the inner product (1.2). Due to Remark 1.20 the dual space  $(X \oplus Y)'$  is also a Hilbert space. Since elements of  $X'$  and  $Y'$  are continuous,  $\mathcal{A}$  is also continuous.  $\mathcal{J}_Y$  and  $\mathcal{J}_{X \oplus Y}$  are the Riesz isomorphisms which map  $Y$  onto  $Y'$  respectively  $X \oplus Y$  onto  $(X \oplus Y)'$ . For  $(f, g), (\tilde{f}, \tilde{g}) \in X' \oplus Y'$  and  $\alpha, \beta \in \mathbb{R}$  we obtain

$$\begin{aligned} \langle \mathcal{A}(\alpha f + \beta \tilde{f}, \alpha g + \beta \tilde{g}), (x, y) \rangle_{(X \oplus Y)', X \oplus Y} \\ = \langle \alpha f + \beta \tilde{f}, x \rangle_{X', X} + \langle \alpha g + \beta \tilde{g}, y \rangle_{Y', Y} \\ = \langle \alpha \mathcal{A}(f, g) + \beta \mathcal{A}(\tilde{f}, \tilde{g}), (x, y) \rangle_{(X \oplus Y)', X \oplus Y} \end{aligned}$$

for all  $(x, y) \in X \oplus Y$ . Therefore,  $\mathcal{A}$  is linear and belongs to  $L(X' \oplus Y', (X \oplus Y)')$ . Let us assume  $\mathcal{A}(f, g) = 0$  for  $f$  in  $X'$  and  $g$  in  $Y'$ . Then it follows

$$\langle \mathcal{A}(f, g), (x, y) \rangle_{(X \oplus Y)', X \oplus Y} = 0$$

for all  $(x, y) \in X \oplus Y$ . So,  $\langle f, x \rangle_{X', X} = -\langle g, y \rangle_{Y', Y}$  for all  $(x, y) \in X \oplus Y$ . This is only true if  $(f, g) = (0, 0)$ . Thus,  $\mathcal{A}$  is injective. Let us choose  $r \in (X \oplus Y)'$ .  $\mathcal{A}$  is surjective if and only if there exists  $(f, g) \in X' \oplus Y'$  such that  $\mathcal{A}(f, g) = r$ .  $\tilde{x} \in X$  and  $\tilde{y} \in Y$  are defined by

$$(1.5) \quad (\tilde{x}, \tilde{y}) = \mathcal{J}_{X \oplus Y}^{-1}(r).$$

If we set  $f = \mathcal{J}_X(\tilde{x})$  and  $g = \mathcal{J}_Y(\tilde{y})$ . We achieve

$$\begin{aligned} \langle \mathcal{A}(f, g), (x, y) \rangle_{(X \oplus Y)', X \oplus Y} &= \langle \mathcal{J}_X(\tilde{x}), x \rangle_{X', X} + \langle \mathcal{J}_Y(\tilde{y}), y \rangle_{Y', Y} \\ &= \langle \tilde{x}, x \rangle_X + \langle \tilde{y}, y \rangle_Y \\ &\stackrel{(1.2)}{=} \langle (\tilde{x}, \tilde{y}), (x, y) \rangle_{X \oplus Y} \\ &\stackrel{(1.5)}{=} \langle r, (x, y) \rangle_{(X \oplus Y)', X \oplus Y} \end{aligned}$$

for all  $(x, y) \in X \oplus Y$ . This implies the surjectivity of  $\mathcal{A}$ . Hence,  $\mathcal{A}$  is a bijection. By applying Theorem 1.5  $\mathcal{A}$  is an isomorphism. Finally, we prove the isometry of

**A:** Let  $(f, g) \in X' \oplus Y'$ . For  $(x, y) = \mathcal{J}_{X \oplus Y}^{-1}(\mathcal{A}(f, g))$  the proof of the surjectivity has shown  $f = \mathcal{J}_X(x)$  and  $g = \mathcal{J}_Y(y)$ . Further we have  $\|\mathcal{J}_X(x)\|_{X'} = \|x\|_X$  and  $\|\mathcal{J}_Y(y)\|_{Y'} = \|y\|_Y$  (Theorem 1.19). This yields

$$\begin{aligned}
 \|\mathcal{A}(f, g)\|_{(X \oplus Y)'}^2 &\stackrel{\text{Lem. 1.12}}{=} \langle \mathcal{A}(f, g), \mathcal{A}(f, g) \rangle_{(X \oplus Y)'} \\
 &\stackrel{(1.4)}{=} \langle \mathcal{J}_{X \oplus Y}^{-1} \mathcal{A}(f, g), \mathcal{J}_{X \oplus Y}^{-1} \mathcal{A}(f, g) \rangle_{X \oplus Y} \\
 &= \langle (x, y), (x, y) \rangle_{X \oplus Y} \\
 &\stackrel{(1.2)}{=} \langle x, x \rangle_X + \langle y, y \rangle_Y \\
 &\stackrel{(1.4)}{=} \langle \mathcal{J}_X(x), \mathcal{J}_X(x) \rangle_{X'} + \langle \mathcal{J}_Y(y), \mathcal{J}_Y(y) \rangle_{Y'} \\
 &= \langle f, f \rangle_{X'} + \langle g, g \rangle_{Y'} \\
 &\stackrel{(1.2)}{=} \langle (f, g), (f, g) \rangle_{X' \oplus Y'} \\
 &\stackrel{\text{Lem. 1.12}}{=} \|(f, g)\|_{X' \oplus Y'}^2.
 \end{aligned}$$

Therefore,  $\mathcal{A}$  is isometric, and the proof is complete.  $\square$

**Remark 1.22.** Due to Proposition 1.21 we identify  $X' \oplus Y'$  with  $(X \oplus Y)'$ , so that we use  $X' \oplus Y'$  as the dual space of  $X \oplus Y$ . The natural inner product on  $X' \oplus Y'$  is

$$\begin{aligned}
 \langle (f, g), (\tilde{f}, \tilde{g}) \rangle_{X' \oplus Y'} &= \langle f, \tilde{f} \rangle_{X'} + \langle g, \tilde{g} \rangle_{Y'} \\
 &\stackrel{(1.4)}{=} \langle \mathcal{J}_X^{-1}(f), \mathcal{J}_X^{-1}(\tilde{f}) \rangle_X + \langle \mathcal{J}_Y^{-1}(g), \mathcal{J}_Y^{-1}(\tilde{g}) \rangle_Y
 \end{aligned}$$

for all  $(f, g), (\tilde{f}, \tilde{g}) \in X' \oplus Y'$ .

**Definition 1.23.** A sequence  $\{x_n\}_{n \in \mathbb{N}}$  converges weakly to an element  $x$  of the Hilbert space  $X$ ,  $x_n \rightharpoonup x$ ,  $n \rightarrow \infty$ , if we have

$$\lim_{n \rightarrow \infty} \langle f, x_n \rangle_{X', X} = \langle f, x \rangle_{X', X}$$

for all  $f \in X'$ .

By the Theorem 1.19, a sequence  $\{x_n\}_{n \in \mathbb{N}}$  converges weakly to an element  $x$  of the Hilbert space  $X$  if  $\lim_{n \rightarrow \infty} \langle y, x_n \rangle_X = \langle y, x \rangle_X$  for all  $y \in X$ . The following results are useful in the Hilbert space approach to differential equations.

**Lemma 1.24.** A bounded sequence in a Hilbert space contains a weakly convergent subsequence.

*Proof.* Let us refer to Theorem 5.12 on page 80 in [9].  $\square$

**Lemma 1.25.** Let  $X$  be a Hilbert space. Then the weak convergence of  $\{x_n\}_{n \in \mathbb{N}}$  in  $X$  implies the boundedness of  $\|x_n\|_X$ .

*Proof.* We refer the reader to Korollar 13.3 on page 61 in [10].  $\square$

**Lemma 1.26.** Let  $X$  be a finite dimensional Hilbert space. Then every weak convergent sequence converge strongly in  $X$ .

*Proof.* We choose a sequence  $\{x_n\}_{n \in \mathbb{N}}$  such that  $x_n \rightharpoonup x$ ,  $n \rightarrow \infty$ , for some  $x \in X$ . Further let  $\{\varphi_1, \dots, \varphi_N\}$  be a orthonormal base for the  $N$ -dimensional Hilbert

space  $X$  and  $x = \sum_{i=1}^N x^i \varphi_i$ ,  $x_n = \sum_{i=1}^N x_n^i \varphi_i$ . Hence the weak convergence and the linearity of the inner product imply

$$0 = \lim_{n \rightarrow \infty} |\langle x_n - x, \varphi_j \rangle_X| \leq \lim_{n \rightarrow \infty} \sum_{i=1}^N |x_n^i - x^i \langle \varphi_i, \varphi_j \rangle_X| = \lim_{n \rightarrow \infty} |x_n^j - x^j|$$

for  $j = 1, \dots, N$ . Thus,

$$0 \leq \lim_{n \rightarrow \infty} \|x_n - x\|_X = \lim_{n \rightarrow \infty} \left\| \sum_{i=1}^N (x_n^i - x^i) \varphi_i \right\|_X \leq \sum_{i=1}^N \lim_{n \rightarrow \infty} |x_n^i - x^i| = 0.$$

□

**Definition 1.27.** Let  $V_1$  and  $V_2$  be two normed linear spaces. A linear operator  $\mathcal{K} \in L(V_1, V_2)$  is called **compact** if  $\mathcal{K}$  takes bounded sets in  $V_1$  into precompact sets in  $V_2$ .

An important property of compact operators is given by:

**Lemma 1.28.** A compact operator maps weakly convergent sequences into norm convergent sequences.

*Proof.* Let us refer to Theorem VI.11 on page 199 in [13]. □

The following lemma is important since one can use it to prove that an operator is compact. For a proof we refer to Theorem VI.12 on page 200 in [13].

**Lemma 1.29.** Let  $B_1$ ,  $B_2$  and  $B$  be Banach spaces and  $\mathcal{A} \in L(B_1, B_2)$ . If  $\mathcal{B} \in L(B_2, B)$  and if  $\mathcal{A}$  or  $\mathcal{B}$  is compact, then  $\mathcal{B} \circ \mathcal{A}$  is compact.

Let us recall the *Fredholm alternative*. For a proof we refer to the corollary on page 203 in [13].

**Theorem 1.30.** If  $\mathcal{K}$  is a compact operator from a Hilbert space  $X$  into itself, then either  $(\mathcal{I}_X - \mathcal{K})^{-1}$  exists or  $\mathcal{K}(v) = v$  has a solution.

A further important theorem about compact operators is given by *Riesz-Schauder theorem*. Therefore we need some more definitions.

**Definition 1.31.** Let  $B$  be a Banach space and  $\mathcal{A} \in L(B)$ . A complex number  $\lambda$  is said to be in the resolvent set  $\rho(\mathcal{A})$  of  $\mathcal{A}$  if  $\lambda \mathcal{I}_B - \mathcal{A}$  is a bijection with a bounded inverse. If  $\lambda \notin \rho(\mathcal{A})$ , then  $\lambda$  is said to be in the spectrum  $\sigma(\mathcal{A})$  of  $\mathcal{A}$ .

We note that by Theorem 1.5,  $\lambda \mathcal{I}_B - \mathcal{A}$  automatically has a bounded inverse if it is bijective. We consider a subset of the spectrum.

**Definition 1.32.** Let  $B$  be a Banach space and  $\mathcal{A} \in L(B)$ . An element  $v \neq 0$  which satisfies  $\mathcal{A}(v) = \lambda v$  for some  $\lambda \in \mathbb{C}$  is called an **eigenvector** (or **eigenfunction**) of  $\mathcal{A}$ ;  $\lambda$  is called the corresponding **eigenvalue**.

**Theorem 1.33.** Let  $\mathcal{K}$  be a compact operator on a Hilbert space  $X$ , then  $\sigma(\mathcal{K})$  is a discrete set having no limit points except perhaps  $\lambda = 0$ . Further any nonzero  $\lambda \in \sigma(\mathcal{K})$  is an eigenvalue of finite multiplicity (i.e., the corresponding space of eigenvectors is finite dimensional).

*Proof.* We refer to Theorem VI.15 on page 203 in [13]. □



If  $x_n \rightharpoonup x$ ,  $n \rightarrow \infty$ , we do not have  $\lim_{n \rightarrow \infty} \|x_n\|_X = \|x\|_X$  in general. But we get the following result.

**Lemma 1.34.** *Let  $X$  be a Hilbert space. If  $x_n \rightharpoonup x$ ,  $n \rightarrow \infty$ , then  $\|x\|_X \leq \liminf_{n \rightarrow \infty} \|x_n\|_X$ .*

*Proof.* Let us refer to Lemma 13.2.1 on page 351 in [15]. □

**Definition 1.35.** *Let  $X$  be a Hilbert space.  $f : X \rightarrow \mathbb{R}$  is called weakly lower semicontinuous if*

$$f(x) \leq \liminf_{n \rightarrow \infty} f(x_n)$$

as  $x_n \rightharpoonup x$ ,  $n \rightarrow \infty$ .

**Remark 1.36.** *From Lemma 1.34 we infer that the norm is weakly lower semicontinuous.*

**Definition 1.37.** *Let  $X$  and  $Y$  be two Hilbert spaces with duals  $X'$  and  $Y'$  respectively. Then we associate with every bounded operator  $\mathcal{A} : X \rightarrow Y'$  defined on all of  $X$ , an adjoint, denoted by  $\mathcal{A}^*$  and defined by*

$$\langle \mathcal{A}^*(y), x \rangle_{X', X} = \langle \mathcal{A}(x), y \rangle_{Y', Y} \quad \text{for all } x \in X \text{ and } y \in Y.$$

**Lemma 1.38.** *Let  $X$  and  $Y$  be two Hilbert spaces with duals  $X'$  and  $Y'$ , respectively, and  $\mathcal{A} \in L(X, Y')$  be a bounded operator; its adjoint  $\mathcal{A}^*$  has the following properties:*

- 1)  $\mathcal{A}^* \in L(Y, X')$ ;
- 2)  $\|\mathcal{A}^*\|_{L(Y, X')} = \|\mathcal{A}\|_{L(X, Y')}$ : the mapping  $\mathcal{A} \mapsto \mathcal{A}^*$  is thus an isometry of  $L(X, Y')$  into  $L(Y, X')$ ;
- 3) We have

$$\begin{aligned} \overline{\text{ran } (\mathcal{J}_X^{-1} \mathcal{A}^*)} &= (\ker (\mathcal{J}_Y^{-1} \mathcal{A}))^\perp, \\ \overline{\text{ran } (\mathcal{J}_Y^{-1} \mathcal{A})} &= \ker (\mathcal{J}_X^{-1} \mathcal{A}^*)^\perp, \\ \ker (\mathcal{J}_X^{-1} \mathcal{A}^*) &= \ker (\mathcal{J}_Y^{-1} \mathcal{A} \mathcal{J}_X^{-1} \mathcal{A}^*), \\ \overline{\text{ran } (\mathcal{J}_Y^{-1} \mathcal{A})} &= \overline{\text{ran } \mathcal{J}_Y^{-1} \mathcal{A} \mathcal{J}_X^{-1} \mathcal{A}^*}. \end{aligned}$$

- 4) *If any of the two subspaces  $\text{ran } (\mathcal{J}_Y^{-1} \mathcal{A})$ ,  $\text{ran } (\mathcal{J}_Y^{-1} \mathcal{A} \mathcal{J}_X^{-1} \mathcal{A}^*)$  is closed, then so the other.*

*Proof.* We refer the reader to Theorem 4 on page 322 in [5], Theorem 8.4 on page 232 and Theorem 11.2 on page 244 in [14]. □

**Lemma 1.39.** *Let  $X$  and  $Y$  be two Hilbert spaces with duals  $X'$  and  $Y'$  respectively. Then we have for every  $\mathcal{A} \in L(X, Y')$ :*

$$\mathcal{A} \text{ is surjective} \implies \mathcal{A}^* \text{ is injective.}$$

*Proof.* Since  $\mathcal{A}$  is surjective, we get  $\text{ran } (\mathcal{A}) = Y'$  and  $\text{ran } (\mathcal{J}_Y^{-1} \mathcal{A}) = Y$ . Due to Lemma 1.38 we obtain  $\ker (\mathcal{J}_X^{-1} \mathcal{A}^*)^\perp = Y$ . This leads to  $\ker (\mathcal{J}_X^{-1} \mathcal{A}^*) = \{0\}$ , and therefore we have  $\ker (\mathcal{A}^*) = \{0\}$ . Hence,  $\mathcal{A}^*$  is injective. □

**Remark 1.40.** *If  $\mathcal{A}^*$  is injective then  $\overline{\text{ran } (\mathcal{J}_Y^{-1} \mathcal{A})} = Y$  by Lemma 1.38. From this we derive  $\overline{\text{ran } (\mathcal{A})} = Y'$ . Therefore,  $\text{ran } (\mathcal{A})$  is only dense in  $Y'$ . If  $\text{ran } (\mathcal{A})$  is closed,  $\mathcal{A}$  is surjective. In Theorem 1.42 we will give equivalent assertions for surjectivity of  $\mathcal{A}$ .*

**Example 1.41.** We shall give an example for an injective operator that is not surjective. Let us choose  $\Omega \subset \mathbb{R}^n$  and  $\mathcal{A}^* = (-\Delta)^{-1} : L^2(\Omega) \rightarrow L^2(\Omega)$ <sup>i</sup>,  $u = \mathcal{A}^*(f)$  is a weak solution of

$$\begin{cases} -\Delta u = f & \text{in } \Omega, \text{ii} \\ u = 0 & \text{on } \partial\Omega, \end{cases}$$

i.e.,  $u \in H_0^1(\Omega)$  solves

$$(1.6) \quad \sum_{i=1}^n \int_{\Omega} u_{x_i} \varphi_{x_i} dx = \int_{\Omega} f \varphi dx \quad \text{for all } \varphi \in H_0^1(\Omega).$$

It can be proved that (1.6) has a unique solution  $u \in H_0^1(\Omega)$ . It is clear, that  $\mathcal{A}^*$  is linear. Further, if  $u = \mathcal{A}^*(f)$  and  $v = \mathcal{A}^*(g)$  then  $w = u - v \in H_0^1(\Omega)$  is a solution of

$$\sum_{i=1}^n \int_{\Omega} w_{x_i} \varphi_{x_i} dx = \int_{\Omega} (f - g) \varphi dx \quad \text{for all } \varphi \in H_0^1(\Omega).$$

$f \neq g$  leads to  $w \neq 0$ . Thus,  $\mathcal{A}^*$  is injective. Further the linear operator  $-\Delta : L^2(\Omega) \rightarrow L^2(\Omega)$  is selfadjoint. This yields

$$\begin{aligned} \langle \mathcal{A}^*(f), g \rangle_{L^2(\Omega)} &\stackrel{(4.1)}{=} - \int_{\Omega} \mathcal{A}^*(f) g dx = - \int_{\Omega} \mathcal{A}^*(f) \underbrace{\Delta(\mathcal{A}^*(g))}_{=g} dx \\ &\stackrel{\text{Int. by parts}}{=} - \int_{\Omega} \underbrace{\Delta(\mathcal{A}^*(f))}_{=f} \mathcal{A}^*(g) dx = \int_{\Omega} \mathcal{A}^*(g) f dx \\ &= \langle \mathcal{A}^*(g), f \rangle_{L^2(\Omega)} \end{aligned}$$

for all  $f, g \in L^2(\Omega)$ . Therefore,  $\mathcal{A}^* \equiv \mathcal{A}$ . Since  $\text{ran}(\mathcal{A}) \subset H_0^1(\Omega) \subsetneq L^2(\Omega)$  we conclude that  $\mathcal{A} = (-\Delta)^{-1}$  is not surjective.

Let  $V_1$  and  $V_2$  be two Hilbert spaces and  $f : M \subset V_1 \rightarrow V_2$ .  $f$  is said to be closed if its graph,  $\{(x, f(x)) : x \in M\}$ , is a closed set in  $V_1 \times V_2$ . By  $D(f)$  we denote the domain of  $f$ .

**Theorem 1.42.** Let  $\mathcal{A}$  be a closed densely defined linear operator from a Hilbert space  $X$  into the dual  $Y'$  of a Hilbert space  $Y$ . Then the following assertions are equivalent:

- 1)  $\mathcal{A}$  is surjective, i.e.,  $\text{ran}(\mathcal{A}) = Y'$ .
- 2) There exists a constant  $K \geq 0$  such that

$$\|v\|_Y \leq K \|\mathcal{A}^*(v)\|_{X'} \quad \text{for all } v \in D(\mathcal{A}^*) \subset Y.$$

- 3)  $\ker(\mathcal{A}^*) = \{0\}$  and  $\text{ran}(\mathcal{A}^*)$  is closed.

*Proof.* Let us refer to Theorem II.19 on page 29 in [3]. □

**Lemma 1.43.** Let  $X$  and  $Y$  be two Hilbert spaces with duals  $X'$  and  $Y'$ , respectively, and  $\mathcal{A} \in L(X, Y')$  be surjective. Then the linear operator  $\mathcal{B} = \mathcal{A} \circ \mathcal{J}_X^{-1} \circ \mathcal{A}^* : Y \rightarrow Y'$  is bounded and invertible, and  $\mathcal{B}^{-1} \in L(Y', Y)$  holds.

<sup>i</sup>We introduce the Hilbert spaces  $L^2(\Omega)$  and  $H_0^1(\Omega)$  in Section 4.

<sup>ii</sup>The Laplacian is given by  $\Delta = \frac{\partial^2}{\partial x_1^2} + \dots + \frac{\partial^2}{\partial x_n^2}$ .

*Proof.* Since  $\mathcal{A}$  is surjective we have  $\ker(\mathcal{A}^*) = \{0\}$  by Lemma 1.39. Due to Lemma 1.38 we get  $\ker(\mathcal{A}\mathcal{J}_X^{-1}\mathcal{A}^*) = \{0\}$ . Thus,  $\mathcal{B}$  is injective. The surjectivity of  $\mathcal{A}$  leads to the surjectivity of  $\mathcal{J}_Y^{-1}\mathcal{A}$  and the closedness of  $\text{ran}(\mathcal{J}_Y^{-1}\mathcal{A})$ . Applying Lemma 1.38 again we obtain  $\text{ran}(\mathcal{J}_Y^{-1}\mathcal{B})$  is closed, too. This leads to

$$Y = \text{ran}(\mathcal{J}_Y^{-1}\mathcal{A}) = \overline{\text{ran}(\mathcal{J}_Y^{-1}\mathcal{A})} \stackrel{\text{Lem. 1.38}}{=} \overline{\text{ran}(\mathcal{J}_Y^{-1}\mathcal{B})} = \text{ran}(\mathcal{J}_Y^{-1}\mathcal{B}).$$

Thus,  $\mathcal{B}$  is surjective and therefore bijective. Since  $\mathcal{A} \in L(X, Y')$  holds  $\mathcal{A}^*$  is bounded (Lemma 1.38). By Lemma 1.8 the operator  $\mathcal{B}$  is continuous. Due to Theorem 1.5 we get  $\mathcal{B}^{-1}$  exists and is continuous.  $\square$

## 2. Local theory of optimization

In this section we recall optimality conditions for infinite dimensional optimization problems.

**Definition 2.1.** Let  $B_1$  and  $B_2$  be Banach spaces and  $f : B_1 \rightarrow B_2$ . If there exists  $\mathcal{A} \in L(B_1, B_2)$ , such that at some point  $x \in B_1$

$$\lim_{\|y\|_{B_1} \rightarrow 0} \frac{\|f(x+y) - f(x) - \mathcal{A}(y)\|_{B_2}}{\|y\|_{B_1}} = 0,$$

then  $\mathcal{A}(y)$  is called the Fréchet-differential of  $f(x)$  at  $x$ , written  $\delta f(x; y)$ . The operator  $\mathcal{A}$  is called the Fréchet-derivative of  $f(x)$  at  $x$ , and we write  $\mathcal{A} = f'(x)$  and  $\delta f(x; y) = f'(x)y$ .

If  $f : B_1 \rightarrow B_2$  has a Fréchet-derivative at  $x$ , it is unique (see Proposition 1 on page 172 in [11]) and  $f$  is continuous at  $x$  (see Proposition 3 on page 173 in [11]). Further we have

$$(2.1) \quad \langle f'(x), y \rangle_{B_2', B_2} = \lim_{t \rightarrow 0} \frac{1}{t} (f(x+ty) - f(x))$$

for each  $y \in B_2$  (see Proposition 2 on page 173 in [11]). If the correspondence  $x \mapsto f'(x)$  is continuous at the point  $x_0$ , we say that the Fréchet derivative of  $f$  is continuous at  $x_0$ . If the derivative of  $f$  is continuous on some open set  $O$ , we say that  $f$  is continuously Fréchet-differentiable on  $O$ .

Let  $X$  be a Hilbert space and  $J : X \rightarrow \mathbb{R}$  be a cost functional. We consider the constrained minimizing problem

$$(2.2) \quad \text{minimize } J(x) \quad \text{subject to } e(x) = 0,$$

where  $e$  is a constraint function from  $X$  into the dual space  $Y'$  of a Hilbert space  $Y$ . The Lagrange functional associated with (2.2) is denoted

$$\mathcal{L}(x, \lambda) = J(x) + \langle e(x), \lambda \rangle_{Y', Y}$$

and the Lagrange multiplier is some specific  $\lambda \in Y$ . Partial derivatives with respect to the variable  $x \in X$  will be denoted by primes.

**Definition 2.2.** We consider the constrained minimizing problem (2.2). If  $x_0 \in X$  is such that  $e'(x_0)$  maps  $X$  onto  $Y'$ , the point  $x_0$  is said to be a regular point of  $e$ .

**Theorem 2.3.** Let  $x_0$  be a regular point of the continuously Fréchet-differentiable function  $e$  mapping the Hilbert space  $X$  into the dual space  $Y'$  of a Hilbert space  $Y$ . Then there is a neighborhood  $U(x_0) \subset X$  of the point  $x_0$ , such that  $e'(x)$  is surjective for all  $x \in U(x_0)$ .

*Proof.* We define the linear mapping

$$(2.3) \quad \mathcal{B}(x) = e'(x) \circ \mathcal{J}_X^{-1} \circ e'(x)^*$$

for all  $x \in X$ . From Lemma 1.43 and Definition 2.2 it follows that  $\mathcal{B}(x_0)$  is invertible. Since the mapping  $x \mapsto e(x)$  is continuously Fréchet-differentiable the function  $x \mapsto \mathcal{B}(x)$  is continuous (Lemma 1.8). Due to Lemma 1.6 and Remark 1.7 there exists a neighborhood  $U(x_0)$  of  $x_0$  such that  $\mathcal{B}(x)$  is invertible. By Lemma 1.38 this yields  $e'(x)$  is surjective for all  $x \in U(x_0)$ .  $\square$

**Remark 2.4.** *Due to Lemma 1.39 the operator  $e'(x)^*$  is injective in the same neighborhood  $U(x_0)$  of  $x_0$ .*

Our aim is to give necessary conditions for an extremum of  $J$  subject to  $e(x) = 0$  where  $J$  is a real-valued functional on a Hilbert space  $X$  and  $e$  is a mapping from  $X$  into the dual space  $Y'$  of a Hilbert space  $Y$ .

**Theorem 2.5.** *If the continuously Fréchet-differentiable functional  $J$  has a local extremum under the constraint  $e(x) = 0$  at a regular point  $x^*$ , then there exists an element  $\lambda^* \in Y$  such that the Lagrangian functional is stationary at  $x^*$ , i.e.,*

$$(2.4) \quad \mathcal{L}'(x^*, \lambda^*) = J'(x^*) + e'(x^*)^* \lambda^* = 0.$$

*Proof.* Let us refer to Theorem 1 on page 243 in [11].  $\square$

**Remark 2.6.** *The equations (2.4) and  $e(x^*) = 0$  are called the first-order necessary optimality condition for a local extremum of  $J$  at the point  $x^*$  under the constraint  $e(x) = 0$ .*

Now we mention the *second-order sufficient optimality condition* for a local minimum. We refer for a proof to Theorem 5.6 in [12].

**Theorem 2.7.** *The twice continuously Fréchet-differentiable functional  $J$  has a local minimum at the point  $x^*$  under the constraint  $e(x) = 0$ , if there exists  $\lambda^* \in Y$  and  $\kappa > 0$  such that*

$$\langle \mathcal{L}''(x^*, \lambda^*)v, v \rangle_{X', X} \geq \kappa \|v\|_X^2 \quad \text{for all } v \in \ker(e'(x^*)).$$

### 3. Analysis of an abstract variational problem

In optimization theory saddle-point problems arise very often. For that purpose we recall basic results for an abstract saddlepoint problem.

Let  $X$  and  $Y$  be two (real) Hilbert spaces with dual spaces  $X'$  and  $Y'$  respectively. The following bilinear forms are given:

$$a : X \times X \rightarrow \mathbb{R}, \quad b : X \times Y \rightarrow \mathbb{R}$$

with norms

$$\|a\| = \sup_{\substack{\varphi, \phi \in X \\ \varphi, \phi \neq 0}} \frac{a(\varphi, \phi)}{\|\varphi\|_X \|\phi\|_X}, \quad \|b\| = \sup_{\substack{\varphi \in X, \psi \in Y \\ \varphi \neq 0, \psi \neq 0}} \frac{b(\varphi, \psi)}{\|\varphi\|_X \|\psi\|_Y}.$$

Let  $l \in X'$  and  $g \in Y'$  be given. We seek  $(x, \lambda) \in X \times Y$  such that

$$(3.1) \quad \begin{cases} a(x, \varphi) + b(\varphi, \lambda) = \langle l, \varphi \rangle_{X', X} & \text{for all } \varphi \in X, \\ b(x, \psi) = \langle g, \psi \rangle_{Y', Y} & \text{for all } \psi \in Y. \end{cases}$$

With the two bilinear forms  $a$  and  $b$  we define linear operators  $\mathcal{A} \in L(X, X')$  and  $\mathcal{B} \in L(X, Y')$ :

$$\begin{aligned} \langle \mathcal{A}(\varphi), \phi \rangle_{X', X} &= a(\varphi, \phi) & \text{for all } \varphi, \phi \in X, \\ \langle \mathcal{B}(\varphi), \psi \rangle_{Y', Y} &= b(\varphi, \psi) & \text{for all } \varphi \in X, \psi \in Y. \end{aligned}$$

The adjoint operator  $\mathcal{B}^*$  of  $\mathcal{B}$  is defined by Definition 1.37. It can be shown that

$$\|\mathcal{A}\|_{L(X, X')} = \|a\|, \quad \|\mathcal{B}\|_{L(X, Y')} = \|b\|.$$

Using these operators (3.1) yields to

$$(3.2) \quad \begin{cases} \mathcal{A}(x) + \mathcal{B}^*(\lambda) &= l & \text{in } X' \\ \mathcal{B}(x) &= g & \text{in } Y'. \end{cases}$$

Let  $V = \ker(\mathcal{B})$  and  $V(g) = \{\varphi \in X : \mathcal{B}(\varphi) = g\}$ . It follows

$$\begin{cases} V(g) = \{\varphi \in X : b(\varphi, \psi) = \langle g, \psi \rangle_{Y', Y} \text{ for all } \psi \in Y\}, \\ V = V(0). \end{cases}$$

Since  $\mathcal{B}$  is continuous,  $V$  is a closed subspace of  $X$ .

Now, we associate with (3.1) the following problem: Find  $u \in V(g)$  such that

$$(3.3) \quad a(x, \varphi) = \langle l, \varphi \rangle_{X', X} \quad \text{for all } \varphi \in V.$$

It is clear, if  $(x, \lambda) \in X \times Y$  is a solution to (3.1), then  $u \in V(g)$  is a solution to (3.3).

We want to find conditions, which ensure that the converse is true. Therefore, we define the polar set  $V^0$  of  $V$  by

$$V^0 = \{f \in X' : \langle f, \varphi \rangle_{X', X} = 0 \text{ for all } \varphi \in V\}.$$

**Lemma 3.1.** *The following properties are equivalent:*

- 1) there exists a constant  $\beta > 0$  such that

$$(3.4) \quad \inf_{\psi \in Y} \sup_{\varphi \in X} \frac{b(\varphi, \psi)}{\|\varphi\|_X \|\psi\|_Y} \geq \beta;$$

- 2) the operator  $\mathcal{B}^*$  is an isomorphism from  $Y$  onto  $V^0$  and

$$\|\mathcal{B}^*(\psi)\|_{X'} \geq \beta \|\psi\|_Y \quad \text{for all } \psi \in Y;$$

- 3) the operator  $\mathcal{B}$  is an isomorphism from  $V^\perp$  onto  $Y'$  and

$$\|\mathcal{B}(\varphi)\|_{Y'} \geq \beta \|\varphi\|_X \quad \text{for all } \varphi \in V^\perp.$$

*Proof.* Let us refer the reader to Lemma 4.1 on page 58 in [8]. □

The condition (3.4) is called *inf-sup condition* or *Babuška-Brezzi condition*. To formulate the next theorem, we introduce the linear continuous (restriction) operator  $\mathcal{F} \in L(X', V')$ :

$$\langle \mathcal{F}(f), \varphi \rangle_{V', V} = \langle f, \varphi \rangle_{X', X} \quad \text{for all } f \in X', \varphi \in V.$$

Obviously, we derive

$$\|\mathcal{F}(h)\|_{V'} \leq \|h\|_{X'}.$$

**Theorem 3.2.** *(3.1) and (3.2) are well-posed (i.e., there is a unique solution) if and only if the following conditions hold:*

- 1) the operator  $\mathcal{F} \circ \mathcal{A}$  is an isomorphism from  $V$  onto  $V'$ ,
- 2) the bilinear form  $b(\cdot, \cdot)$  satisfies the inf-sup condition (3.4).

*Proof.* We refer to Theorem 4.1 on page 59 in [8].  $\square$

The next corollary is an important application of Theorem 3.2.

**Corollary 3.3.** *Let  $a$  be  $V$ -elliptic, i.e., there exists a constant  $\kappa_0 > 0$  such that*

$$a(\varphi, \varphi) \geq \kappa_0 \|\varphi\|_X^2 \quad \text{for all } \varphi \in V.$$

*Then, (3.1) respectively (3.2) is well-posed if and only if the bilinear form  $b(\cdot, \cdot)$  satisfies the inf-sup condition (3.4). Let  $(x, \lambda) \in X \times Y$  be the unique solution to (3.1). Then we have the estimates*

$$\begin{aligned} \|x\|_X &\leq \frac{1}{\kappa_0} \|l\|_{X'} + \left( \frac{1}{\beta} + \frac{\|a\|}{\kappa_0 \beta} \right) \|g\|_{Y'}, \\ \|\lambda\|_Y &\leq \left( \frac{1}{\beta} + \frac{\|a\|}{\kappa_0 \beta} \right) \|l\|_{X'} + \left( \frac{\|a\|}{\beta^2} + \frac{\|a\|^2}{\kappa_0 \beta^2} \right) \|g\|_{Y'}. \end{aligned}$$

*Proof.* Let us refer to Corollary 4.1 on page 61 in [8] and Theorem 1.1 on page 42 in [4].  $\square$

Now we turn to the discretization of (3.1). Let  $X_N$  and  $Y_M$  be two finite dimensional spaces such that

$$X_N \subset X, \quad \dim X_N = N, \quad Y_M \subset Y, \quad \dim Y_M = M.$$

Let  $X'_N$  and  $Y'_M$  denote their dual spaces with the dual norms:

$$(3.5) \quad \|l_N\|_{X'_N} = \sup_{\varphi_N \in X_N} \frac{\langle l_N, \varphi_N \rangle_{X', X}}{\|\varphi_N\|_X}, \quad \|g_M\|_{Y'_M} = \sup_{\psi_M \in Y_M} \frac{\langle g_M, \psi_M \rangle_{Y', Y}}{\|\psi_M\|_Y}.$$

Clearly,

$$\|l\|_{X'_N} \leq \|l\|_{X'}, \quad \|g\|_{Y'_M} \leq \|g\|_{Y'} \quad \text{for all } (l, g) \in X' \times Y'.$$

Like in the continuous case, we associate with  $a(\cdot, \cdot)$  and  $b(\cdot, \cdot)$  the operators  $\mathcal{A}_N \in L(X, X'_N)$ ,  $\mathcal{B}_M \in L(X, Y'_M)$  and  $\mathcal{B}_N^* \in L(Y, X'_N)$  defined by

$$\begin{aligned} \langle \mathcal{A}_N(\phi), \varphi_N \rangle_{X', X} &= a(\phi, \varphi_N) \quad \text{for all } \varphi_N \in X_N, \quad \text{for all } \phi \in X, \\ \langle \mathcal{B}_M(\varphi), \psi_M \rangle_{Y', Y} &= b(\varphi, \psi_M) \quad \text{for all } \psi_M \in Y_M, \quad \text{for all } \varphi \in X, \\ \langle \mathcal{B}_N^*(\psi), \varphi_N \rangle_{X', X} &= b(\varphi_N, \psi) \quad \text{for all } \varphi_N \in X_N, \quad \text{for all } \psi \in Y. \end{aligned}$$

$\mathcal{B}_N^*$  is not the dual operator of  $\mathcal{B}_M$  but if  $\mathcal{B}_M$  is restricted to  $X_N$  and  $\mathcal{B}_N^*$  to  $Y_M$ , then  $\mathcal{B}_M$  and  $\mathcal{B}_N^*$  are indeed dual operators. Moreover, we have:

$$\|\mathcal{B}_M(\varphi)\|_{Y'_M} \leq \|\mathcal{B}(\varphi)\|_{Y'} \quad \text{for all } \varphi \in X$$

with similar inequalities for  $\|\mathcal{A}_N\|_{X'_N}$  and  $\|\mathcal{B}_N^*(\psi)\|_{X'_N}$ . For each  $g \in Y'$ , we define the finite-dimensional analogue of  $\tilde{V}(g)$ :

$$V_{NM}(g) = \{\varphi_N \in X_N : b(\varphi_N, \psi_M) = \langle g, \psi_M \rangle_{Y', Y} \text{ for all } \psi_M \in Y_M\}$$

and we set

$$V_{NM} = V_{NM}(0) = \ker(\mathcal{B}_M) \cap X_N = \{\varphi_N \in X_N : b(\varphi_N, \psi_M) = 0 \text{ for all } \psi_M \in Y_M\}.$$

Right away we remark that generally  $V_{NM} \not\subset V$  and  $V_{NM}(g) \not\subset V(g)$  because  $Y_M$  is a proper subspace of  $Y$ . Now we approximate (3.1) by

$$(3.6) \quad \begin{cases} a(x_N, \varphi_N) + b(\varphi_N, \lambda_M) &= \langle l, \varphi_N \rangle_{X', X} & \text{for all } \varphi_N \in X_N, \\ b(x_N, \psi_M) &= \langle g, \psi_M \rangle_{Y', Y} & \text{for all } \psi_M \in Y_M \end{cases}$$

and we associate with (3.6) the following problem:

$$(3.7) \quad a(x_N, \varphi_N) = \langle l, \varphi_N \rangle_{X', X} \quad \text{for all } \varphi_N \in V_{NM}.$$

Here again, the first component  $x_N$  of any solution  $(x_N, \lambda_M)$  of (3.6) is also a solution of (3.7). The converse is true due to the next theorem.

**Theorem 3.4.** *Let  $(x, \lambda)$  be a solution to (3.1).*

1) *Assume that the following conditions hold:*

- (a)  $V_{NM}(g) \neq \emptyset$ ;
- (b) *there exists a constant  $\kappa_{NM} > 0$  such that:*

$$(3.8) \quad a(\varphi_N, \varphi_N) \geq \kappa_{NM} \|\varphi_N\|_X^2 \quad \text{for all } \varphi_N \in V_{NM}.$$

*Then (3.7) has a unique solution  $x_N \in V_{NM}(g)$  and the “error bound” holds*

$$(3.9) \quad \begin{aligned} & \|x - x_N\|_X \\ & \leq \left(1 + \frac{\|a\|}{\kappa_{NM}}\right) \inf_{\varphi_N \in V_{NM}(g)} \|x - \varphi_N\|_X + \frac{\|b\|}{\kappa_{NM}} \inf_{\psi_M \in Y_M} \|\lambda - \psi_M\|_Y. \end{aligned}$$

2) *Assume that hypothesis holds and, in addition, that:*

- (c) *there exists a constant  $\beta_{NM} > 0$  such that*

$$(3.10) \quad \sup_{\varphi_N \in X} \frac{b(\varphi_N, \psi_M)}{\|\varphi_N\|_X} \geq \beta_{NM} \|\psi_M\|_Y \quad \text{for all } \psi_M \in Y_M.$$

*Then  $V_{NM}(g) \neq \emptyset$  and there exists a unique  $\lambda_M$  in  $Y_M$  such that  $(x_N, \lambda_M)$  is the only solution of (3.6). Moreover, We get the estimate*

$$(3.11) \quad \begin{aligned} & \|x - x_N\|_X + \|\lambda - \lambda_M\|_Y \\ & \leq K_{NM} \inf_{\varphi_N \in X_N} \|x - \varphi_N\|_X + C_{NM} \inf_{\psi_M \in Y_M} \|\lambda - \psi_M\|_Y, \end{aligned}$$

*where the constants are given by*

$$K_{NM} = \left(1 + \frac{1}{\beta_{NM}} + \frac{\|a\|}{\kappa_{NM}}\right) \left(1 + \frac{\|b\|}{\beta_{NM}}\right), \quad C_{NM}^2 = \beta_{NM} + \|b\| + \frac{\|b\|}{\kappa_{NM}}.$$

*Proof.* For the proof we refer the reader to Theorem 1.1 on page 114 in [8] and Proposition 2.4 on page 54 in [4].  $\square$

**Remark 3.5.** 1) *It can be shown that*

$$\inf_{\varphi_N \in V_{NM}(g)} \|x - \varphi_N\|_X \leq \left(1 + \frac{1}{\beta_{NM}}\right) \inf_{\phi_N \in X_N} \|x - \phi_N\|_X$$

*holds. The condition*

$$(3.12) \quad \sup_{\varphi_N \in X} \frac{b(\varphi_N, \psi_M)}{\|\varphi_N\|_X} \geq \beta^* \|\psi_M\|_Y \quad \text{for all } \psi_M \in Y_M.$$

*is called the inf-sup condition.*

2) *If there exists two positive constants  $\kappa^* > 0$  and  $\beta^* > 0$  such that  $\kappa_{NM} \geq \kappa^*$  and  $\beta_{NM} \geq \beta^*$ , then both  $K_{NM}$  and  $C_{NM}$  are independent of  $N$  and  $M$ :*

$$K = \left(1 + \frac{1}{\beta^*} + \frac{\|a\|}{\kappa^*}\right) \left(1 + \frac{\|b\|}{\beta^*}\right), \quad C = \beta^* + \|b\| + \frac{\|b\|}{\kappa^*}.$$

- 3) The bound (3.9) can be slightly improved without making use of condition (3.10). Indeed, it can be proved (see Remark 1.1 on page 116 in [8])

$$\begin{aligned} & \|x - x_N\|_X \\ & \leq \left(1 + \frac{\|a\|}{\kappa^*}\right) \inf_{\varphi_N \in V_{NM}(g)} \|x - \varphi_N\|_X + \frac{1}{\kappa^*} \inf_{\psi_M \in Y_M} \sup_{\varphi_N \in V_{NM}} \frac{b(\varphi_N, \lambda - \psi_M)}{\|\varphi_N\|_X}. \end{aligned}$$

Note that the expression

$$\inf_{\psi_M \in Y_M} \sup_{\varphi_N \in V_{NM}} \frac{b(\varphi_N, \lambda - \psi_M)}{\|\varphi_N\|_X}$$

takes into account the fact that  $V_{NM} \not\subset V$ . It vanishes when  $V_{NM} \subset V$ .

- (1) Observe that the bilinear form  $a$  is  $V_{NM}$ -elliptic as soon as  $a(\varphi_N, \varphi_N) > 0$  for all  $\varphi_N \neq 0$ . Analogously, the bilinear form  $b$  satisfies the discrete inf-sup condition (3.8) provided  $\ker(\mathcal{B}_M) \cap Y_M = \{0\}$ . But in the general case both assumptions have to be checked.

The following lemma established a useful criterion for (3.12)

**Lemma 3.6.** *The inf-sup condition (3.12) holds with a constant  $\beta^* > 0$  independent of  $N, M$  if and only if there is a restriction operator  $r^N \in L(X, X_N)$  satisfying:*

$$b(\varphi - r^N(\varphi), \psi_M) = 0 \quad \text{for all } (\psi_M, \varphi) \in Y_M \times X$$

and

$$\|r^N(\varphi)\|_X \leq K \|\varphi\|_X \quad \text{for all } \varphi \in X$$

with a constant  $K > 0$  independent of  $N$ .

*Proof.* Let us refer to Lemma 1.1 on page 117 in [8]. □

#### 4. Function spaces

Since we are interested in optimal control of partial differential equations, we require basic definitions of function spaces and associated results.

The term *domain* and the symbol  $\Omega$  shall be reserved for an open set in the  $n$ -dimensional, real Euclidian space  $\mathbb{R}^n$ . A typical point of  $\mathbb{R}^n$  is denoted by  $x = (x_1, \dots, x_n)$ ; its norm  $|x|_2 = (\sum_{j=1}^n x_j^2)^{\frac{1}{2}}$ .

**Definition 4.1.** *Let  $\Omega$  be an open set of  $\mathbb{R}^n$  with boundary  $\Gamma$ . If boundary  $\Gamma$  is a  $(N - 1)$ -dimensional manifold of class  $C^r$  ( $r \geq 1$  which must be specified) and  $\Omega$  is locally located on one side of  $\Gamma$ , we will say,  $\Omega$  is of class  $C^r$ . The boundary  $\Gamma$  is locally Lipschitz if for any  $x \in \Gamma$ , there is a neighborhood such that  $\Gamma$  admits a representation as a hypersurface  $x_n = \theta(x_1, \dots, x_{n-1})$ , where  $\theta$  is Lipschitz continuous and  $x_1, \dots, x_{n-1}$  are rectangular coordinates in  $\mathbb{R}^n$  in a basis that may be different from the canonical basis.*

**Remark 4.2.** *If  $\Omega$  is of class  $C^1$ , then  $\Omega$  is locally Lipschitz.*

**Definition 4.3.**  $\Omega \subset \mathbb{R}^n$  is said to be disconnected if there exists two nonempty subsets  $\Omega_1, \Omega_2 \subset \Omega$  such that

$$\Omega_1 \cap \Omega_2 = \emptyset, \quad \Omega = \Omega_1 \cup \Omega_2.$$

If  $\Omega$  is not disconnected,  $\Omega$  is called connected.



If  $\mathbf{a} = (a_1, \dots, a_n)$  is an  $n$ -tuple of nonnegative integers  $a_j$ , we call  $\mathbf{a}$  a *multi-index* and denote by  $x^{\mathbf{a}}$  the monomial  $x_1^{a_1} \cdots x_n^{a_n}$ , which has degree  $|\mathbf{a}| = \sum_{j=1}^n a_j$ . Similarly

$$D^{\mathbf{a}} = \frac{\partial^{a_1}}{\partial x_1^{a_1}} \cdots \frac{\partial^{a_n}}{\partial x_n^{a_n}}$$

denotes a differential operator of order  $|\mathbf{a}|$ .  $D^{(0, \dots, 0)}\varphi = \varphi$  for a function  $\varphi$  on  $\Omega$ . We shall write  $G \subset \subset \Omega$  provided  $\overline{G} \subset \Omega$  and  $\overline{G}$  is a compact subset of  $\mathbb{R}^n$ . If  $\varphi$  is a function defined on  $G$ , we define the *support* of  $\varphi$  as

$$\text{supp } \varphi = \overline{\{x \in G : \varphi(x) \neq 0\}}.$$

We say that  $\varphi$  has *compact support* in  $\Omega$  if  $\text{supp } \varphi \subset \subset \Omega$ . We shall denote by  $\partial G$  the boundary of  $G$  in  $\mathbb{R}^n$ . For any nonnegative integer  $m$  let  $C^m(\Omega)$  be the vector space consisting of all functions  $\varphi$  which, together with all their partial derivatives  $D^{\mathbf{a}}\varphi$  of order  $|\mathbf{a}| \leq m$ , are continuous on  $\Omega$ . We abbreviate  $C^0(\Omega) = C(\Omega)$ . Let  $C^\infty(\Omega) = \bigcap_{m=1}^\infty C^m(\Omega)$ . The subspaces  $C_0(\Omega)$  and  $C_0^\infty(\Omega)$  consist of all those functions in  $C(\Omega)$  and  $C^\infty(\Omega)$ , respectively, which have compact support in  $\Omega$ .

Since  $\Omega$  is open, functions in  $C^m(\Omega)$  need not be bounded on  $\Omega$ . If  $\varphi \in C(\Omega)$  is bounded and uniformly continuous on  $\Omega$ , then it possesses a unique, bounded, continuous extension to the closure  $\overline{\Omega}$  of  $\Omega$ . Accordingly, we define the vector space  $C^m(\overline{\Omega})$  to consist of all those functions  $\varphi \in C^m(\Omega)$  for which  $D^{\mathbf{a}}\varphi$  is bounded and uniformly continuous on  $\Omega$  for  $0 \leq |\mathbf{a}| \leq m$ .  $C^m(\overline{\Omega})$  is a Banach space with norm given by

$$\|\varphi\|_{C^m(\overline{\Omega})} = \max_{0 \leq |\mathbf{a}| \leq m} \sup_{x \in \Omega} |D^{\mathbf{a}}\varphi(x)|.$$

If  $0 < s \leq 1$ , we define  $C^{m,s}(\overline{\Omega})$  to be the subspace of  $C^m(\overline{\Omega})$  consisting of those functions  $\varphi$ , for which,  $0 \leq |\mathbf{a}| \leq m$ ,  $D^{\mathbf{a}}\varphi$  satisfies in  $\Omega$  a *Hölder condition of exponent  $s$* , that is, there exists a constant  $K > 0$  such that

$$|D^{\mathbf{a}}\varphi(x) - D^{\mathbf{a}}\varphi(y)| \leq K |x - y|^s, \quad x, y \in \Omega.$$

$C^{m,s}(\overline{\Omega})$  is a Banach space with norm given by

$$\|\varphi\|_{C^{m,s}(\overline{\Omega})} = \|\varphi\|_{C^m(\overline{\Omega})} + \max_{0 \leq |\mathbf{a}| \leq m} \sup_{x, y \in \Omega} \frac{|D^{\mathbf{a}}\varphi(x) - D^{\mathbf{a}}\varphi(y)|}{|x - y|^s}.$$

Where no confusion of domains may occur we will write  $\|\cdot\|_{C^m}$  in place of  $\|\cdot\|_{C^m(\overline{\Omega})}$  and  $\|\cdot\|_{C^{m,s}}$  instead of  $\|\cdot\|_{C^{m,s}(\overline{\Omega})}$ . It should be noted that for  $0 < r < s \leq 1$ ,

$$C^{m+1}(\overline{\Omega}) \subsetneq C^{m,s}(\overline{\Omega}) \subsetneq C^{m,r}(\overline{\Omega}) \subsetneq C^m(\overline{\Omega}).$$

We denote by  $L^p(\Omega)$  the class of all measurable functions  $\varphi$ , defined on  $\Omega$ , for which

$$\int_{\Omega} |\varphi|^p dx < \infty.$$

The function  $\|\cdot\|_{L^p(\Omega)}$  defined by

$$\|\varphi\|_{L^p(\Omega)} = \left( \int_{\Omega} |\varphi|^p dx \right)^{\frac{1}{p}}$$

is a norm on  $L^p(\Omega)$  provided  $1 \leq p < \infty$ .  $L^2(\Omega)$  is a separable Hilbert space with the inner product

$$(4.1) \quad \langle \varphi, \psi \rangle_{L^2(\Omega)} = \int_{\Omega} \varphi \psi dx.$$

$$\|\varphi\|_{L^\infty(\Omega)} = \operatorname{ess\,sup}_{x \in \Omega} |\varphi(x)|$$

is a norm on  $L^\infty(\Omega)$ . In situations where no confusion of domains may occur we shall write  $\|\cdot\|_{L^p}$  in place of  $\|\cdot\|_{L^p(\Omega)}$  and  $\langle \cdot, \cdot \rangle_{L^2}$  instead of  $\langle \cdot, \cdot \rangle_{L^2(\Omega)}$ . The next theorem is known as *Fischer-Riesz theorem*. For a proof we refer the reader to Theorem 2.10 on page 26 and Corollary 2.11 on page 27 in [1].

**Theorem 4.4.**  *$L^p(\Omega)$  is a Banach space if  $1 \leq p \leq \infty$ . Every convergent sequence in  $L^p(\Omega)$  has a subsequence converging pointwise a.e. on  $\Omega$ .*

We shall have occasion to use a generalization of Hölders inequality.

**Proposition 4.5.** *Let  $\Omega$  be a bounded domain in  $\mathbb{R}^n$ . Assume  $p_1, \dots, p_m \in [1, \infty]$  and  $\sum_{j=1}^m p_j^{-1} = 1$  (with  $\infty^{-1} = 0$ ). If  $\varphi_i \in L^{p_i}(\Omega)$  for  $i = 1, \dots, m$  then the function  $\prod_{j=1}^m \varphi_j$  belongs to  $L^1(\Omega)$  and we have the estimate*

$$(4.2) \quad \int_{\Omega} \left| \prod_{j=1}^m \varphi_j \right| dx \leq \prod_{j=1}^m \|\varphi_j\|_{L^{p_j}}.$$

*Proof.* The proof follows by an induction argument.

- 1) The case  $m = 2$  follows directly from Hölders inequality (Theorem 2.3 on page 23 and Remark 2.5 on page 24 in [1]).
- 2) Induction hypothesis: We assume that the statement is proved for  $m \geq 2$ .
- 3) Now let  $p_1, \dots, p_{m+1} \in [1, \infty]$ ,  $m \geq 2$  and  $\sum_{j=1}^{m+1} p_j^{-1} = 1$ . For  $\varphi_i \in L^{p_i}(\Omega)$  we know

$$\int_{\Omega} \left| \prod_{j=1}^{m+1} \varphi_j \right| dx \stackrel{\text{Hölder's ineq.}}{\leq} \left\| \prod_{j=1}^m \varphi_j \right\|_{L^q} \|\varphi_{m+1}\|_{L^{p_{m+1}}}$$

with  $q^{-1} + p_{m+1}^{-1} = 1$ . Since  $\varphi_i \in L^{p_i}(\Omega)$  we have  $\varphi_j^q \in L^{\frac{p_i}{q}}(\Omega)$  and  $\sum_{j=1}^m \left(\frac{p_i}{q}\right)^{-1} = 1$ . By applying the induction hypothesis we obtain:

$$\left\| \prod_{j=1}^m \varphi_j \right\|_{L^q} = \left( \int_{\Omega} \left| \prod_{j=1}^m \varphi_j^q \right| dx \right)^{\frac{1}{q}} \leq \prod_{j=1}^m \|\varphi_j^q\|_{L^{\frac{p_i}{q}}}^{\frac{1}{q}} = \prod_{j=1}^m \|\varphi_j\|_{L^{p_i}}.$$

Therefore, we get  $\prod_{j=1}^{m+1} \varphi_j \in L^1(\Omega)$ , and the formula (4.2) is proved.  $\square$

Next we introduce Sobolev spaces of integer. These spaces are defined over an arbitrary domain  $\Omega \subset \mathbb{R}^n$  and are vector subspaces of various spaces  $L^p(\Omega)$ . We define a function  $\|\cdot\|_{W^{m,p}(\Omega)}$ , where  $m$  is a nonnegative integer and  $1 \leq p \leq \infty$ , as follows:

$$(4.3) \quad \|\varphi\|_{W^{m,p}(\Omega)} = \left( \sum_{0 \leq |\mathbf{a}| \leq m} \|D^{\mathbf{a}}\varphi\|_{L^p(\Omega)}^p \right)^{\frac{1}{p}} \quad \text{if } 1 \leq p < \infty,$$

$$(4.4) \quad \|\varphi\|_{W^{m,\infty}(\Omega)} = \max_{0 \leq |\mathbf{a}| \leq m} \|D^{\mathbf{a}}\varphi\|_{L^\infty(\Omega)}$$

for any function  $\varphi$  for which the right side makes sense. In situations where no confusion of domains may occur we shall write  $\|\cdot\|_{W^{m,p}}$  in place of  $\|\cdot\|_{W^{m,p}(\Omega)}$ . (4.3) and (4.4) are norms on any linear space on which the right side takes finite

values provided functions are identified in the space if they are equal a.e. in  $\Omega$ . We define to any given values of  $m$  and  $p$  the Sobolev spaces:

$$W^{m,p}(\Omega) = \text{the completion of } \{\varphi \in C^m(\Omega) : \|\varphi\|_{W^{m,p}} < \infty\} \text{ with respect to the norm } \|\cdot\|_{W^{m,p}(\Omega)},$$

$$(4.5) \quad W_0^{m,p}(\Omega) = \text{the closure of } C_0^\infty(\Omega) \text{ in the space } W^{m,p}(\Omega).$$

Equipped with the appropriate norms (4.3) and (4.4), these are called *Sobolev spaces* over  $\Omega$ . Clearly,  $W^{0,p}(\Omega) = L^p(\Omega)$ , and if  $1 \leq p < \infty$ ,  $W_0^{0,p}(\Omega) = L^p(\Omega)$  by Theorem 2.19 on page 31 in [1].

**Lemma 4.6.**  $W^{m,p}(\Omega)$  is a Banach space.

*Proof.* Let us refer the reader to Theorem 3.2 on page 45 in [1]. □

A function  $\varphi$  defined a.e. on  $\Omega$  is said to be *locally integrable* on  $\Omega$  provided  $\varphi \in L^1(M)$  for every measurable  $M \subset\subset \Omega$ . In this case we write  $\varphi \in L_{\text{loc}}^1(\Omega)$ . We now define the concept of weak derivative of a locally integrable function  $\varphi \in L_{\text{loc}}^1(\Omega)$ . If there exists a  $\phi \in L_{\text{loc}}^1(\Omega)$ , such that

$$\int_{\Omega} \varphi D^{\mathbf{a}}\psi \, dx = (-1)^{|\mathbf{a}|} \int_{\Omega} \phi \psi \, dx \quad \text{for all } \psi \in C_0^\infty(\Omega),$$

it is unique up to sets of measure zero and it is called the *weak* or *distributional partial derivative* of  $\varphi$  and is denoted by  $D^{\mathbf{a}}\varphi$ . If  $\varphi$  is sufficiently smooth to have continuous partial derivative  $D^{\mathbf{a}}\varphi$  in the usual (classical) sense, then  $D^{\mathbf{a}}\varphi$  is also a distributional derivative of  $\varphi$ . For example a function  $\varphi$ , continuous on  $\mathbb{R}$ , which has a bounded derivative  $\varphi'$  except at finitely many points, has a derivative in the distributional sense.

We can introduce  $W^{m,p}(\Omega)$  for any  $m \geq 0$  and  $1 \leq p < \infty$  in a different way:

$$W^{m,p}(\Omega) = \{\varphi \in L^p(\Omega) : D^{\mathbf{a}}\varphi \in L^p(\Omega) \text{ for } 0 \leq |\mathbf{a}| \leq m, \\ D^{\mathbf{a}}\varphi \text{ is the weak partial derivative}\},$$

(Theorem 3.16 on page 52 in [1]).

**Lemma 4.7.**  $W^{m,p}(\Omega)$  is separable if  $1 \leq p < \infty$ . In particular,  $W^{m,2}(\Omega)$  is a separable Hilbert space with inner product

$$\langle \varphi, \psi \rangle_{W^{m,2}} = \sum_{0 \leq |\mathbf{a}| \leq m} \langle D^{\mathbf{a}}\varphi, D^{\mathbf{a}}\psi \rangle_{L^2} = \sum_{0 \leq |\mathbf{a}| \leq m} \int_{\Omega} D^{\mathbf{a}}\varphi D^{\mathbf{a}}\psi \, dx.$$

*Proof.* We refer to Theorem 3.5 on page 47 in [1]. □

We write  $H^m(\Omega)$  in place of  $W^{m,2}(\Omega)$  and  $H_0^m(\Omega)$  instead of  $W_0^{m,2}(\Omega)$ .

**Lemma 4.8.** Let  $\Omega \subset \mathbb{R}^n$  be open and bounded with Lipschitz-continuous boundary  $\Gamma$ . For  $m \geq 1$  and real  $p$  with  $1 \leq p < \infty$  there exists a continuous linear extension operator  $\mathcal{F} : W^{m,p}(\Omega) \rightarrow W^{m,p}(\mathbb{R}^n)$  such that

$$\mathcal{F}(\varphi)|_{\Omega} = \varphi \quad \text{for all } \varphi \in W^{m,p}(\Omega).$$

*Proof.* Let us refer the reader to Theorem 1.2 on page 5 in [8]. □

If the boundary  $\Gamma$  is Lipschitz continuous, one can show that there exists an operator  $\tau_\Gamma : H^1(\Omega) \rightarrow L^2(\Gamma)$ , linear and continuous, such that

$$\tau_\Gamma \varphi = \text{trace of } \varphi \text{ on } \Gamma \text{ for every } \varphi \in C^1(\overline{\Omega}).$$

It then seems natural to call  $\tau_\Gamma \varphi$  *the trace of  $\varphi$  on  $\Gamma$* , and denote it by  $\varphi|_\Gamma$  even if  $\varphi$  is a general function in  $H^1(\Omega)$ . A deeper analysis shows that by taking all the traces of all functions of  $H^1(\Omega)$  one does not obtain the whole space  $L^2(\Gamma)$  but only a subset of it. Further, such a subspace contains  $H^1(\Gamma)$ . Hence we have,

$$H^1(\Gamma) \subsetneq \tau_\Gamma(H^1(\Omega)) \subsetneq L^2(\Gamma) \equiv H^0(\Gamma).$$

Therefore we introduce the space

$$H^{1/2}(\Gamma) = \tau_\Gamma(H^1(\Omega))$$

with

$$\|g\|_{H^{1/2}} = \inf_{\substack{\varphi \in H^1(\Omega) \\ \tau_\Gamma \varphi = g}} \|\varphi\|_{H^1}.$$

In a similar way one can see that the traces of functions in  $H^2(\Omega)$  belong to a space  $H^{3/2}(\Gamma)$ . We define

$$H^{3/2}(\Gamma) = \tau_\Gamma(H^2(\Omega))$$

and

$$\|g\|_{H^{3/2}} = \inf_{\substack{\varphi \in H^2(\Omega) \\ \tau_\Gamma \varphi = g}} \|\varphi\|_{H^2}.$$

We shall need a special form of the *Sobolev embedding theorem*. The normed space  $V_1$  is said to be *continuous embedded* in the normed space  $V_2$ , and write  $V_1 \hookrightarrow V_2$  to designate this embedding, provided

- 1)  $V_1$  is a vector subspace of  $V_2$ , and
- 2) the identity operator  $\mathcal{I}_{V_1, V_2}$  defined on  $V_1$  into  $V_2$  by  $\mathcal{I}_{V_1, V_2}(v) = v$  for all  $v \in V_1$  is continuous.

We say,  $V_1$  is *compact embedded* in  $V_2$ ,  $V_1 \hookrightarrow\hookrightarrow V_2$ , if the embedding operator  $\mathcal{I}_{V_1, V_2}$  is compact.

**Definition 4.9.** *The bounded domain  $\Omega \subset \mathbb{R}^n$  has a locally Lipschitz boundary, if each point  $x$  on the boundary  $\partial\Omega$  should have a neighborhood  $U(x)$  such that  $\partial\Omega \cap U(x)$  is the graph of a Lipschitz-continuous function.*

**Lemma 4.10.** *Let  $j \geq 0$  and  $\Omega$  be a bounded domain in  $\mathbb{R}^n$ , such that  $\partial\Omega$  is locally Lipschitz continuous. Suppose  $2m > n > 2(m-1)$ . Then*

$$H^{j+m}(\Omega) \hookrightarrow C^{j,s}(\overline{\Omega})$$

for  $s \in (0, m - n/2]$ .

*Proof.* We choose  $p = 2$  and apply Theorem 5.4, Part II, on page 98 in [1].  $\square$

**Remark 4.11.** *If the assumptions of Lemma 4.10 are satisfied, there exists a constant  $K > 0$  such that*

$$\|\varphi\|_{C^{j,s}} \leq K \|\varphi\|_{H^{j+m}} \quad \text{for all } \varphi \in H^{j+m}(\Omega).$$

Now we mention a special case of the *Rellich-Kondrachov theorem*. For a proof we refer the reader to Theorem 6.2 on page 144 in [1].

**Lemma 4.12.** *Let  $\Omega \subset \mathbb{R}^n$  be a bounded domain with locally Lipschitz boundary,  $j, m$  be integers,  $j \geq 0, m \geq 1$ , and let  $1 \leq p < \infty$ . Then:*

$$\begin{aligned} W^{j+m,p}(\Omega) &\hookrightarrow\hookrightarrow W^{j,q}(\Omega) && \text{if } n > mp \text{ and } 1 \leq q < \frac{np}{n-mp}, \\ W^{j+m,p}(\Omega) &\hookrightarrow\hookrightarrow W^{j,q}(\Omega) && \text{if } n = mp \text{ and } 1 \leq q < \infty, \\ W^{j+m,p}(\Omega) &\hookrightarrow\hookrightarrow W^{j,q}(\Omega) && \text{if } n < mp \text{ and } 1 \leq q \leq \infty. \end{aligned}$$

**Remark 4.13.** *If we choose  $n \leq 3, j = 0, m = 1, p = 2, q = 4$  we get  $H^1(\Omega)$  is compact embedded in  $L^4(\Omega)$  from Lemma 4.12. By Lemma 1.28 the embedding operator maps weakly convergent sequences in  $H_0^1(\Omega)$  into norm convergent sequences in  $L^4(\Omega)$ . In particular, if  $n = 1$  holds, we have  $H^1(\Omega) \hookrightarrow\hookrightarrow L^q(\Omega)$  for  $1 \leq q \leq \infty$ .*

The following lemma characterize weak convergence in  $H^1(0, 1)$ .

**Lemma 4.14.** *For every  $f \in H^1(0, 1)'$  there exists  $\varphi_1, \varphi_2 \in L^2(0, 1)$  such that*

$$\langle f, y \rangle_{(H^1)', H^1} = \int_0^1 y' \varphi_1 + y \varphi_2 dx$$

for all  $y \in H^1(0, 1)$ .

*Proof.* Let us refer the reader to Theorem 3.8 on page 48 in [1]. □

The next lemma gives a useful application of Green's formula. For a proof we refer to Lemma 1.4 on page 10 in [8].

**Lemma 4.15.** *Let  $\Omega$  be a bounded open subset of  $\mathbb{R}^n$  with Lipschitz-continuous boundary  $\Gamma$  and  $\mathbf{n} = (n_1, \dots, n_n)$  the outward unit normal.*

1) *For  $u, v \in H^1(\Omega)$  and  $1 \leq i \leq n$  we have*

$$\int_{\Omega} uv_{x_i} dx = - \int_{\Omega} u_{x_i} v dx + \int_{\Gamma} \tau_{\Gamma}(uv) n_i ds.$$

2) *If in addition  $u \in H^2(\Omega)$  we derive*

$$(4.6) \quad \sum_{i=1}^n \int_{\Omega} u_{x_i} v_{x_i} dx = - \sum_{i=1}^n \int_{\Omega} u_{x_i x_i} v dx + \sum_{i=1}^n \int_{\Gamma} \tau_{\Gamma}(vu_{x_i}) n_i ds.$$

**Remark 4.16.** *By using*

$$\Delta u = \sum_{i=1}^n u_{x_i x_i}, \quad \nabla u = (u_{x_1}, \dots, u_{x_n}), \quad \mathbf{n} = (n_1, \dots, n_n)$$

we conclude from (4.6)

$$\int_{\Omega} \nabla u \cdot \nabla v dx = - \int_{\Omega} \Delta u v dx + \int_{\Gamma} \tau_{\Gamma}(v \nabla u) \cdot \mathbf{n} ds.$$

In the remaining of this section we recall the Fourier transform. By  $\iota = \sqrt{-1}$  we denote the *imaginary unit* in  $\mathbb{C}$ .

**Definition 4.17.** *For  $u \in L^1(\mathbb{R})$  we define the Fourier transform of  $u$ , denoted by  $\hat{u}$ :*

$$\hat{u}(y) = \int_{\mathbb{R}} \exp(-\iota xy) u(x) dx \quad \text{for all } y \in \mathbb{R}.$$

**Remark 4.18.** The mapping  $u \mapsto \hat{u}$  defined by Definition 4.17 is obviously linear. From the inequality

$$|\hat{u}(y)| \leq \|u\|_{L^1(\mathbb{R})} \quad \text{for all } y \in \mathbb{R}$$

we deduce:

$$(4.7) \quad \begin{cases} \text{if } u \in L^1(\mathbb{R}), \hat{u} \text{ is a bounded continuous function on } \mathbb{R} \text{ with} \\ \|\hat{u}\|_{L^\infty(\mathbb{R})} \leq \|u\|_{L^1(\mathbb{R})}. \end{cases}$$

In addition we have the Riemann-Lebesgue theorem (see also Theorem IX.7 on page 327 in [13]):

$$(4.8) \quad \hat{u}(y) \rightarrow 0 \quad \text{in } \mathbb{C} \quad \text{when } |y| \rightarrow \infty.$$

Let  $v \in C_0^\infty(\mathbb{R})$ . Then we derive, as a result of an integration by parts:

$$\hat{v}(y) = \frac{1}{iy} \int_{\mathbb{R}} \exp(-ixy) v'(x) dx,$$

from which we have

$$(4.9) \quad |\hat{v}(y)| \leq \frac{1}{|y|} \|v'\|_{L^1(\mathbb{R})} \rightarrow 0, \quad \text{when } |y| \rightarrow \infty.$$

Now, as  $C_0^\infty(\mathbb{R})$  is dense in  $L^1(\mathbb{R})$ , if  $u \in L^1(\mathbb{R})$ , for all  $\varepsilon > 0$ , we find  $v \in C_0^\infty(\mathbb{R})$ , such that  $\|v - u\|_{L^1(\mathbb{R})} \leq \varepsilon/2$ . Thus from

$$\hat{u}(y) = \widehat{(u - v)}(y) + \hat{v}(y),$$

we derive using (4.9)

$$|\hat{u}(y)| \leq \|u - v\|_{L^1(\mathbb{R})} + \frac{1}{|y|} \|v'\|_{L^1(\mathbb{R})} < \frac{\varepsilon}{2} + \frac{\varepsilon}{2} = \varepsilon$$

for  $|y|$  sufficiently large (see (4.9)), from which (4.8) follows.

We put

$$\mathcal{S}(\mathbb{R}) = \mathcal{S} = \{u \in C^\infty(\mathbb{R}) : \text{for all } \alpha, l \in \mathbb{N}, x^\alpha u^{(l)}(x) \rightarrow 0 \text{ as } |x| \rightarrow \infty\}.$$

$\mathcal{S}$  is the space of functions of class  $C^\infty$  of rapid decay at infinity, which is not a normed space, but of which the topology can be defined by the (denumerable) sequence of semi-norms

$$u \mapsto \sup_{x \in \mathbb{R}} |x^k u^{(l)}(x)| = d_{kl}(u)$$

which yields a complete metrisable space:

$$d(u, v) = \sum_{k, l \in \mathbb{N}} a_{kl} \frac{d_{kl}(u - v)}{1 + d_{kl}(u - v)} \quad \text{for all } u, v \in \mathcal{S},$$

where the coefficients  $a_{kl}$  are chosen to be such that  $\sum_{k, l \in \mathbb{N}} a_{kl} = 1$ , is a distance on  $\mathcal{S}$ . We should notice that if  $u \in \mathcal{S}$ , then  $xu^{(l)}(x) \in L^p(\mathbb{R})$  for all  $p \geq 1$  and for all  $k, l \in \mathbb{N}$ . Further,  $\mathcal{S}$  is dense in  $L^p(\mathbb{R})$  for all  $p$  with  $1 \leq p < \infty$  (on the contrary  $\mathcal{S}$  is not dense in  $L^\infty(\mathbb{R})$ ). For  $u \in \mathcal{S}$  we can thus define its Fourier transform by Definition 4.17, as well as the Fourier transform of  $x^k u^{(l)}(x)$  for all  $k, l \in \mathbb{N}$ . Hence, we also have  $\hat{u} \in \mathcal{S}$ . Further we have the inversion formula

$$u(x) = \frac{1}{2\pi} \int_{\mathbb{R}} \exp(ixy) \hat{u}(y) dy \quad \text{for all } x \in \mathbb{R}.$$

Let  $u, v \in \mathcal{S}$ . Then the following properties are valid:

1)  $\widehat{u^{(k)}} = (\iota y)^k \hat{u}$ .

2) Parseval's formula:

$$\int_{\mathbb{R}} u(x) \overline{v(x)} dx = \int_{\mathbb{R}} \hat{u}(y) \overline{\hat{v}(y)} dy.$$

3) Plancharel's formula:

$$\int_{\mathbb{R}} |u(x)|^2 dx = \frac{1}{\sqrt{2\pi}} \int_{\mathbb{R}} |\hat{u}(y)|^2 dy.$$

### 5. Evolution problems: variational methods

In the optimal control of parabolic systems we shall need to make frequent use of the notion of integral of a Banach space-valued function  $\varphi$  defined on an interval of  $\mathbb{R}$ . We begin therefore with a brief discussion of the Bochner integral, referring the reader to [16], for instance, for further details and proofs of our assertions.

Let  $B$  be a Banach space with norm denoted by  $\|\cdot\|_B$  and  $\{M_1, \dots, M_m\}$  be a finite collection of mutually disjoint, measurable subsets of  $\mathbb{R}$ , each having finite measure, and let  $\{b_1, \dots, b_m\}$  be a corresponding collection of points of  $B$ . The function  $\varphi$  on  $\mathbb{R}$  defined by

$$\varphi(t) = \sum_{j=1}^m \chi_{M_j}(t) b_j, \text{ iii}$$

is called a *simple function*. For simple functions we define

$$\int_{\mathbb{R}} \varphi(t) dt = \sum_{j=1}^m \mu(M_j) b_j$$

where  $\mu(M)$  denotes the (Lebesgue) measure of  $M$ . Let  $M$  be a measurable set in  $\mathbb{R}$  and  $\varphi$  an arbitrary function defined a.e. on  $M$  into  $B$ . The function  $\varphi$  is called (*strongly*) *measurable* on  $M$  if there exists a sequence  $\{\varphi_n\}_{n \in \mathbb{N}}$  of simple functions with supports in  $M$  such that

$$(5.1) \quad \lim_{n \rightarrow \infty} \|\varphi_n(t) - \varphi(t)\|_B = 0 \quad \text{a.e. in } M.$$

It can be shown that any function  $\varphi$  whose range is separable is measurable provided the scalar-valued function  $\langle f, \varphi(\cdot) \rangle_{B', B}$  is measurable on  $M$  for each  $f \in B'$ . We suppose that a sequence of simple functions  $\varphi_n$  satisfying (5.1) can be chosen in such a way that

$$\lim_{n \rightarrow \infty} \int_M \|\varphi_n(t) - \varphi(t)\|_B dt = 0.$$

Then  $\varphi$  is called (*Bochner*) *integrable* on  $M$  and we define

$$(5.2) \quad \int_M \varphi(t) dt = \lim_{n \rightarrow \infty} \int_{\mathbb{R}} \varphi_n(t) dt.$$

The integrals on the right side of (5.2) do converge in (the norm topology of)  $B$  to a limit which is independent of the choice of approximating sequence  $\{\varphi_n\}_{n \in \mathbb{N}}$ .

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iii  $\chi_M$  denotes the characteristic function of  $M$ :  $\chi_M(t) = 1$  if  $t \in M$  and  $\chi_M(t) = 0$  if  $t \notin M$ .

A measurable function  $\varphi$  is integrable on  $M$  if and only if  $\|\varphi(\cdot)\|_B$  is (Lebesgue) integrable on  $M$ :

$$\left\| \int_M \varphi(t) dt \right\|_B \leq \int_M \|\varphi\|_B dt.$$

**Definition 5.1.** Let  $-\infty \leq a < b \leq \infty$ . We denote by  $L^p(a, b; B)$  <sup>iv</sup> the linear space of (equivalence classes of) functions  $\varphi$  measurable on  $(a, b)$  into  $B$  such that

- 1)  $\varphi$  is measurable for  $dt$ ,
- 2)

$$(5.3) \quad \begin{cases} \|\varphi\|_{L^p(a, b; B)} = \left( \int_a^b \|\varphi(t)\|_B dt \right)^{\frac{1}{p}} < \infty & \text{if } 1 \leq p < \infty, \\ \|\varphi\|_{L^\infty(a, b; B)} = \operatorname{ess\,sup}_{t \in (a, b)} \|\varphi(t)\|_B < \infty & \text{if } p = \infty. \end{cases}$$

If  $\varphi \in L^p(c, d; B)$  for every  $c, d$  with  $a < c < d < b$ , then we write  $\varphi \in L^p_{\text{loc}}(a, b; B)$ , and, if  $p = 1$ , call  $\varphi$  locally integrable.

**Proposition 5.2.** For  $1 \leq p \leq \infty$ ,  $L^p(a, b; B)$  is a Banach space.

*Proof.* We refer to Proposition 1 on page 469 in [7]. □

**Proposition 5.3.** If  $B$  is a Banach space,  $a$  and  $b$  are finite,  $f \in B'$  and  $\varphi \in L^p(a, b; B)$  for  $p \geq 1$  we have

$$\left\langle f, \int_a^b \varphi(t) dt \right\rangle_{B', B} = \int_a^b \langle f, \varphi(t) \rangle_{B', B} dt.$$

*Proof.* Let us refer the reader to Corollary 2 on page 470 in [7]. □

**Definition 5.4.** By  $\mathcal{D}(a, b)$  we denote the linear space  $C_0^\infty(a, b)$ . We say, the sequence  $\{\varphi_n\}_{n \in \mathbb{N}}$  tends to zero in  $\mathcal{D}(a, b)$  if there is a closed subspace  $M \subset (a, b)$  such that  $\varphi_n^{(i)}(t) = 0$  for all  $i \in \mathbb{N}$ , for all  $n \in \mathbb{N}$  and for all  $t \in (a, b) \setminus M$  and it follows  $\|\varphi_n^{(i)}\|_{L^\infty} \rightarrow 0$  as  $n \rightarrow \infty$  for all  $i \in \mathbb{N}$ . We call every continuous linear mapping of  $\mathcal{D}(a, b)$  into a Banach space  $B$  a vectorial distribution over  $(a, b)$  with values into  $B$ , and we write  $\mathcal{D}'(a, b; B) = L(\mathcal{D}(a, b), B)$ .

**Remark 5.5.** Let  $(B_1, B_2)$  be a pair of Banach spaces with  $B_1 \hookrightarrow B_2$ . Then we derive

$$L^p(a, b; B_1) \hookrightarrow L^p(a, b; B_2) \quad \text{for } 1 \leq p \leq \infty.$$

**Proposition 5.6.** Let  $B$  be a Banach space and  $u \in L^1_{\text{loc}}(a, b; B)$ . Then the mapping

$$\varphi \mapsto \int_a^b \varphi(t) u(t) dt$$

is a distribution over  $(a, b)$  with values in  $B$ .

*Proof.* We refer the reader to Proposition 4 on page 470 in [7]. □

**Remark 5.7.** We identify the function  $u$  with the distribution with which it is associated.

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<sup>iv</sup>When there is no risk of confusion, we shall write the simplified notation  $L^p(B)$ .



**Proposition 5.8.** *Let  $B$  be a Banach space. The functions  $u, v \in L^1_{\text{loc}}(a, b; B)$  define the same distributions if and only if  $u$  and  $v$  are equal (in a scalar sense) a.e.*

*Proof.* Let us refer to Proposition 5 on page 471 in [7]. □

**Remark 5.9.** *The Proposition 5.8 means that for all  $f \in B'$  the functions:  $t \rightarrow \langle f, u(t) \rangle_{B', B}$  and  $t \rightarrow \langle f, v(t) \rangle_{B', B}$  are equal almost everywhere. If  $B$  is separable, this implies  $u = v$  a.e.*

**Definition 5.10.** *Let  $B$  be a Banach space,  $f \in \mathcal{D}'(a, b; B)$  and  $m$  a nonnegative integer. Then the mapping  $\varphi \mapsto (-1)^m f\left(\frac{d^m \varphi}{dt^m}\right)$ ,  $\varphi \in \mathcal{D}(a, b)$ , is a distribution — the distributional derivative — that we denote by  $\frac{d^m f}{dt^m}$ . We have:*

$$\frac{d^m f}{dt^m}(\varphi) = (-1)^m f\left(\frac{d^m \varphi}{dt^m}\right) \quad \text{for all } \varphi \in \mathcal{D}(a, b).$$

**Remark 5.11.** *Let  $B_1$  and  $B_2$  be two separable Banach spaces. If  $u \in L^1_{\text{loc}}(a, b; B)$  and if  $B$  is a space of functions of the variable  $x$ , for instance  $B = L^p(\Omega)$ , then  $u$  is identified with a function  $u(t, x)$ .  $u(t)$  denotes the mapping  $x \mapsto u(t, x)$  for almost all  $t$ . The distributional derivative  $\frac{du}{dt}$  is identified with the derivative  $\frac{\partial u}{\partial t}$  in  $\mathcal{D}'(a, b; B)$ . We use the following notation for the derivative of  $u$  with respect to  $t$ :*

$$\frac{du}{dt} \quad \text{or} \quad u' \quad \text{or} \quad u_t.$$

**Definition 5.12.** *Let  $B$  be a Banach space and  $u \in L^2(a, b; B)$ . Then for all  $\varphi \in \mathcal{D}(a, b)$ :*

$$\frac{du}{dt}(\varphi) = - \int_a^b u(t)\varphi'(t) dt.$$

*We say that  $u' = \frac{du}{dt} \in L^2(a, b; B)$  if there exists  $v \in L^2(a, b; B)$  such that:*

$$\begin{cases} \text{for all } \varphi \in \mathcal{D}(a, b), & v(\varphi) = -u(\varphi'), \\ \text{i.e.: } \int_a^b v(t)\varphi(t) dt = \int_a^b u(t)\varphi'(t) dt. \end{cases}$$

The space we shall introduce next is of fundamental importance. We consider two real, separable Hilbert spaces  $V, H$ . It is supposed that  $V$  is dense in  $H$  so that, by identifying  $H$  and its dual  $H'$ , we have

$$(5.4) \quad V \hookrightarrow H \equiv H' \hookrightarrow V',$$

each space being dense in the following.

**Definition 5.13.** *Let  $a, b \in \mathbb{R} \cup \{-\infty, +\infty\}$ ,  $a < b$ . Moreover,  $V$  and  $H$  are Hilbert spaces satisfying (5.4). The space  $W(a, b; V)$  is given by*

$$W(a, b; V) = \left\{ \varphi : \varphi \in L^2(a, b; V), \frac{d\varphi}{dt} \in L^2(a, b; V') \right\}^v$$

**Proposition 5.14.** *The space  $W(a, b; V)$  endowed with the norm*

$$\|\varphi\|_{W(V)} = \left( \|\varphi\|_{L^2(V)}^2 + \left\| \frac{d\varphi}{dt} \right\|_{L^2(V')}^2 \right)^{\frac{1}{2}} = \left( \int_a^b \|\varphi(t)\|_V^2 + \left\| \frac{d\varphi(t)}{dt} \right\|_{V'}^2 dt \right)^{\frac{1}{2}}$$

*is a Hilbert space.*

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<sup>v</sup>When there is no risk of confusion, we shall write the simplified notation  $W(V)$ .

*Proof.* We refer the reader to Proposition 6 on page 473 in [7].  $\square$

We are now interested in regularity properties of elements belonging to  $W(V)$ . For a proof of the following lemma we refer to Theorem 1 on page 472 in [7].

**Lemma 5.15.** *For  $a, b \in \mathbb{R}$ , every  $\varphi \in W(V)$  is almost everywhere equal to a continuous function of  $(a, b)$  in  $H$ . Further, we have:*

$$W(V) \hookrightarrow C([a, b]; H),$$

*the space  $C([a, b]; H)$  being equipped with the norm of uniform convergence.*

**Remark 5.16.** *By Lemma 5.15 it makes sense to speak of the traces  $\varphi(a), \varphi(b) \in H$  for  $\varphi \in W(V)$  with  $[a, b] \subset \mathbb{R}$ . Moreover, we can show, that the mapping  $\varphi \mapsto \varphi(a)$  from  $W(V)$  is surjective (see Remark 5 on page 477 in [7]).*

**Lemma 5.17.** *Let  $[a, b] \subset \mathbb{R}$  and  $\varphi, \phi \in W(V)$ . Then*

$$\int_a^b \left\langle \frac{d\varphi(t)}{dt}, \phi(t) \right\rangle_{V', V} dt + \int_a^b \left\langle \frac{d\phi(t)}{dt}, \varphi(t) \right\rangle_{V', V} dt = \langle \varphi(b), \phi(b) \rangle_H - \langle \varphi(a), \phi(a) \rangle_H$$

*Proof.* We refer to Theorem 2 on page 477 in [7].  $\square$

A very useful property is

**Proposition 5.18.** *For  $\varphi \in W(V)$  and  $\psi \in V$  we obtain:*

$$\left\langle \frac{d\varphi(\cdot)}{dt}, \psi \right\rangle_{V', V} = \frac{d}{dt} \langle \varphi(\cdot), \psi \rangle_H$$

*in the distributional sense.*

*Proof.* Let us refer to Proposition 7 on page 477 in [7].  $\square$

Let  $X$  and  $Y$  be two separable Hilbert spaces with  $X \hookrightarrow Y$  and  $X$  being dense in  $Y$ . We now define the space  $W(a, b; X, Y)$  by

$$W(a, b; X, Y) = \left\{ \varphi : \varphi \in L^2(a, b; X), \frac{d\varphi}{dt} \in L^2(a, b; Y) \right\}$$

equipped with the norm

$$\|\varphi\|_{W(X, Y)} = \left( \|\varphi\|_{L^2(X)}^2 + \left\| \frac{d\varphi}{dt} \right\|_{L^2(Y)}^2 \right)^{\frac{1}{2}} = \left( \int_a^b \|\varphi(t)\|_X^2 + \left\| \frac{d\varphi(t)}{dt} \right\|_Y^2 dt \right)^{\frac{1}{2}}$$

It can be shown that  $W(a, b; X, Y)$  is a Hilbert space and that

$$(5.5) \quad \begin{cases} \text{i) } X \text{ is dense in } [X, Y]_\theta, \theta \in [0, 1]^{\text{vi}}, \\ \text{ii) } W(a, b; X, Y) \hookrightarrow C([a, b]; [X, Y]_{1/2}) \end{cases}$$

(see (1.61) on page 480 in [7]).

We are given two real, separable Hilbert spaces  $V$  and  $H$ .  $V$  is supposed to be dense in  $H$  and we identify  $H$  with its dual  $H'$ . Moreover:  $V \hookrightarrow H \hookrightarrow V'$ . We denote by  $W(V)$  the space  $W(0, T; V)$  with  $0 < T < \infty$ . By Lemma 5.15 we derive

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<sup>vi</sup>With  $[X, Y]_1 = Y$  and  $[X, Y]_0 = X$ , the space  $[X, Y]_\theta$  is called the *holomorphic interpolant* of the spaces  $X$  and  $Y$  (see Chapter VIII, §3 in [6]).

$W(V) \hookrightarrow C([0, T]; H)$ . For each  $t \in [0, T]$  we are given a continuous bilinear form over  $V \times V$  and we make the hypothesis:

$$(5.6) \quad \begin{cases} \text{for every } \varphi, \psi \in V, \text{ the function } t \mapsto a(t; \varphi, \psi) \text{ is measurable and} \\ \text{there exists a constant } K = K(T) > 0 \text{ (independent of } t \in (0, T), \\ \varphi, \psi \in V) \text{ such that} \\ |a(t; \varphi, \psi)| \leq K \|\varphi\|_V \|\psi\|_V \quad \text{for all } \varphi, \psi \in V. \end{cases}$$

Therefore, for each  $t \in [0, T]$  the bilinear form  $a(t; \varphi, \psi)$  defines a continuous linear operator  $\mathcal{A}(t)$  from  $V$  into  $V'$  with

$$\sup_{t \in (0, T)} \|\mathcal{A}(t)\|_{L(V, V')} \leq K.$$

We make the following assumption (of coercivity over  $V$  with respect to  $H$ ):

$$(5.7) \quad \begin{cases} \text{there exists } \lambda, \alpha \text{ constants, } \alpha > 0 \text{ such that} \\ a(t; \varphi, \varphi) + \lambda \|\varphi\|_H^2 \geq \alpha \|\varphi\|_V^2 \quad \text{for all } t \in [0, T] \text{ and } \varphi \in V. \end{cases}$$

We give some examples of bilinear forms  $a(t; \varphi, \psi)$ .

**Example 5.19.** 1) We take  $V = H_0^1(\Omega)$ ,  $H = L^2(\Omega)$ .

$$a(t; \varphi, \psi) = a(\varphi, \psi) = \langle \psi, \varphi \rangle_{H_0^1}.$$

Then (5.6) and (5.7) holds with  $\alpha = 1$  and  $\lambda = 0$ .

2) We take  $V = H^1(\Omega)$ ,  $H = L^2(\Omega)$ .

$$a(t; \varphi, \psi) = \langle \varphi, \psi \rangle_{H^1}.$$

Then (5.6) and (5.7) holds with  $\alpha = 1$  and  $\lambda = 0$ .

3) Now let  $V$  be a closed subspace of  $H^1(\Omega)$  with

$$H_0^1(\Omega) \hookrightarrow V \hookrightarrow H^1(\Omega) \text{ and } H = L^2(\Omega).$$

We set  $Q = \Omega \times (0, T)$  and

$$a(t; \varphi, \psi) = \sum_{i,j=1}^n \int_{\Omega} a_{ij}(x, t) \frac{\partial \varphi}{\partial x_i} \frac{\partial \psi}{\partial x_j} dx + \int_{\Omega} a_0(x, t) \varphi \psi dx$$

where  $a_{ij}, a_0 \in L^\infty(Q)$ ,  $1 \leq i, j \leq n$  and

$$\sum_{i,j=1}^n a_{ij}(x, t) \xi_i \xi_j \geq \alpha \sum_{i=1}^n \xi_i^2$$

for a constant  $\alpha > 0$  and for  $\xi_i \in \mathbb{R}$  a.e. in  $Q$ . Then, for  $\lambda$  large enough we derive for all  $\varphi \in H^1(\Omega)$ :

$$a(t; \varphi, \varphi) + \lambda \|\varphi\|_{L^2}^2 \geq \alpha \|\varphi\|_V^2.$$

Let

$$(5.8) \quad u_0 \in H, \quad f \in L^2(V')$$

be given. We are looking for

$$(5.9) \quad u \in W(V),$$

such that

$$(5.10) \quad \begin{cases} \frac{d}{dt} \langle u(\cdot), \varphi \rangle_H + a(\cdot; u(\cdot), \varphi) = \langle f, \varphi \rangle_{V', V} \\ \text{in the sense of } \mathcal{D}'(a, b) \text{ for all } \varphi \in V, \end{cases}$$

and

$$(5.11) \quad u(0) = u_0.$$

**Remark 5.20.** 1) From Lemma 5.15 the initial condition (5.11) is sensible.

2) Due to Proposition 5.18 we have

$$\frac{d}{dt} \langle u(\cdot), \varphi \rangle_H = \left\langle \frac{du(\cdot)}{dt}, \varphi \right\rangle_{V', V} \quad \text{for all } \varphi \in V.$$

**Remark 5.21.** If we set  $u = we^{kt}$ ,  $k \in \mathbb{R}$ ,  $w$  satisfies

$$\left\langle \frac{dw(\cdot)}{dt}, \varphi \right\rangle_{V', V} + a(\cdot; w(\cdot), \varphi) + k \langle w(\cdot), \varphi \rangle_H = \langle e^{-kt} f(\cdot), \varphi \rangle_{V', V}$$

and

$$w(0) = u_0$$

by changing  $u$  to  $ue^{kt}$  and choosing  $k$ , we can assume that (5.7) holds with  $\lambda = 0$  (that has no consequences since  $T$  is finite). In the following we, we shall therefore make the hypothesis:

$$(5.12) \quad a(\cdot; \varphi, \varphi) \geq \alpha \|\varphi\|_V^2 \quad \text{for all } t \in [0, T]^{\text{vii}} \text{ and } \varphi \in V.$$

**Theorem 5.22.** We suppose  $V$ ,  $H$  are given and satisfy  $V \hookrightarrow H \hookrightarrow V'$  and  $a(\cdot; u, \varphi)$  satisfies (5.6), (5.12).  $u_0$  and  $f$  are given and satisfy (5.8). Then there exists a unique solution of (5.9)-(5.11).

*Proof.* Let us refer the reader to Theorem 1 on page 512 and Theorem 2 on page 513 in [7].  $\square$

If  $u$  is the solution of (5.10)-(5.11), we derive

$$\frac{1}{2} \|u(t)\|_H^2 + \int_0^t a(s; u(s), u(s)) ds = \frac{1}{2} \|u_0\|_H^2 + \int_0^t \langle f(s), u(s) \rangle_{V', V} ds,$$

the so called *energy equality*, as the quantity

$$E(t) = \frac{1}{2} \|u_0\|_H^2 + \int_0^t \langle f(s), u(s) \rangle_{V', V} ds$$

represents the energy of the system.

**Theorem 5.23.** Let  $a(t; \varphi, \psi)$  satisfy (5.6) and (5.12),  $(u_0, f)$ ,  $(u_0^*, f^*) \in H \times L^2(V')$  and let  $u$  and  $u^*$  be the corresponding solutions of (5.9)-(5.11). Then

$$\begin{aligned} \|u - u^*\|_{L^1(H)} &\leq \left( \|u_0 - u_0^*\|_H^2 + \frac{1}{\alpha} \|f - f^*\|_{L^2(V')}^2 \right)^{\frac{1}{2}} \\ \|u - u^*\|_{L^2(V)} &\leq \frac{1}{\sqrt{\alpha}} \left( \|u_0 - u_0^*\|_H^2 + \frac{1}{\alpha} \|f - f^*\|_{L^2(V')}^2 \right)^{\frac{1}{2}} \end{aligned}$$

*Proof.* We refer to Theorem 3 on page 520 in [7].  $\square$

<sup>vii</sup>Or likewise  $t \in [0, T]$  a.e.

**Remark 5.24.** We assume  $f \in L^1(H)$ . Then it can be proved that (5.9)-(5.11) has a unique solution in the space

$$W^*(V) = \{\varphi : \varphi \in L^2(V), \varphi' \in L^2(V') + L^1(H)\}.$$

The same estimates as in Theorem 5.23 hold (see Remark 6 on page 521 and Theorem 4 on page 522 in [7]). Thus problem (5.9)-(5.11) can be considered with

$$f = f_1 + f_2, \quad f_1 \in L^1(H), \quad f_2 \in L^2(V')$$

and there exists a unique solution to this problem, the assumptions being those of the beginning of this section.

### 6. Principal Notations

In the following list of symbols we give the symbol and a descriptive name or phrase for an explanation. The number at the indicate the pagenummer on which the symbols are introduced.

$M \setminus M_1$	complement of $M_1$ in $M$	1
$\overline{M}$	closure of a set $M$	1
$M^\circ$	interior of $M$	1
$\partial M$	boundary of a set $M$	1
$M_1 \times M_2$	Cartesian product of the sets $M_1$ and $M_2$	1
$\text{ran}(f)$	range of the function $f$	1
$\ker(f)$	kernel of the function $f$	1
$g \circ f$	compositive mapping given by $x \mapsto g(f(x))$	1
$\ \cdot\ _V$	norm on a (real) normed linear space $V$	2
$v_n \xrightarrow{n \rightarrow \infty} v$	(strong) convergence of the sequence $\{v_n\}_{n \in \mathbb{N}}$ to $v$	2
$B(v; \rho)$	open ball of radius $\rho$ about the point $v$	2
$U(v)$	neighborhood of $v$	2
$L(V_1, V_2)$	set of bounded linear operators $\mathcal{A} : V_1 \rightarrow V_2$	3
$\ \mathcal{A}\ _{L(V_1, V_2)}$	norm of $\mathcal{A} \in L(V_1, V_2)$	3
$\mathcal{A}^{-1}$	inverse of a bounded linear operator $\mathcal{A}$	3
$L(V)$	set of bounded linear operators $\mathcal{A} : V \rightarrow V$	3
$\mathcal{I}_V$	identity on a normed linear space $V$	3
$\langle \cdot, \cdot \rangle_X$	(real) inner product on a Hilbert space $X$	3
$a(\cdot, \cdot)$	continuous bilinear form	4
$X \oplus Y$	direct sum of the Hilbert spaces $X$ and $Y$	4
$M^\perp$	orthogonal complement of a closed space $M \subset X$	4
$\mathcal{P}_M$	linear projection onto the subset $M$	5
$B'$	dual space of a Banach space $B$	5
$\langle \cdot, \cdot \rangle_{B', B}$	duality pairing of $B'$ with its Banach space $B$	5
$\mathcal{J}_X$	Riesz isomorphism of which maps a Hilbert space $X$ onto its dual $X'$	6
$x_n \xrightarrow{n \rightarrow \infty} x$	weak convergence of a sequence $\{x_n\}_{n \in \mathbb{N}}$ to $x$	7
$\mathcal{K}$	compact operator	8
$\rho(\mathcal{A})$	resolvent set of a bounded linear operator $\mathcal{A}$	8
$\sigma(\mathcal{A})$	spectrum of a bounded linear operator $\mathcal{A}$	8
$\mathcal{A}^*$	adjoint of a bounded linear operator $\mathcal{A}$	9
$\delta f(x; y)$	Fréchet differential of $f$ at $x$ in the direction $y$	11
$f'(x)$	Fréchet derivative of $f$ at the point $x$	11

$J$	cost functional	11
$e$	constraint function	11
$\mathcal{L}$	Lagrange functional	11
$x^*$	(local) optimal solution of a constrained minimizing problem	12
$\lambda^*$	Lagrange multiplier of a constrained minimizing problem	12
$ \cdot _2$	Euclidian norm in $\mathbb{R}^n$	16
$\Omega$	open set of $\mathbb{R}^n$	16
$\mathbf{a}$	multi index	17
$\text{supp}(f)$	support of $f$	17
$C^m(\Omega)$	linear space consisting of $m$ -times continuously differentiable functions	17
$C^m(\overline{\Omega})$	Banach space which is subspace of $C^m(\Omega)$	17
$\ \cdot\ _{C^m}$	norm on $C^m(\overline{\Omega})$	17
$C_0(\Omega), C_0^\infty(\Omega)$	subspaces of $C(\Omega)$ respectively $C^m(\Omega)$	17
$C^{m,s}(\overline{\Omega})$	Hölder spaces	17
$\ \cdot\ _{C^{m,s}}$	norm on $C^{m,s}(\overline{\Omega})$	17
$L^p(\Omega)$	$L^p$ -spaces	17
$\ \cdot\ _{L^p}$	norm on $L^p(\Omega)$	17
$L_{\text{loc}}^1(\Omega)$	linear space consisting locally integrable functions	19
$W^{m,p}(\Omega)$	Sobolev spaces	19
$\ \cdot\ _{W^{m,p}}$	norm on $W^{m,p}(\Omega)$	18
$W_0^{m,p}(\Omega)$	Sobolev spaces	19
$\ \cdot\ _{W_0^{m,p}}$	norm on $W_0^{m,p}(\Omega)$	18
$H^m(\Omega), H_0^m(\Omega)$	Hilbert spaces consisting of weak differentiable functions	19
$\ \cdot\ _{H^m}, \ \cdot\ _{H_0^m}$	norm on $H^m(\Omega)$ respectively $H_0^m(\Omega)$	19
$V_1 \hookrightarrow V_2$	continuous embedding of a normed linear space $V_1$ into a normed linear space $V_2$	20
$V_1 \hookrightarrow\hookrightarrow V_2$	compact embedding of a normed linear space $V_1$ into a normed linear space $V_2$	20
$\chi_M$	characteristic function of a set $M$	23
$L^p(0, T; B)$	Banach space consisting of functions with vector values in a Banach space $B$	24
$\ \cdot\ _{L^p(a,b;B)}$	norm on the Banach space $L^p(0, T; B)$	24
$W(V)$	Hilbert space consisting of vector-valued functions	25
$\ \cdot\ _{W(V)}$	norm on $W(V)$	25

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