Functional Analysis

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LECTURE 1: INTRODUCTION TO MEASURE THEORY

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The course contents will be mostly on measure theory with some basic functional analysis notions about Banach and Hilbert spaces.

1.1 Measure spaces

Def. (σ -algebra)

Let X be an arbitrary set and Σ be a collection of subsets of X. Σ is called a σ -algebra if it is closed w.r. to complement and countable unions:

i
$$A \in \Sigma \implies A^c = X \setminus A \in \Sigma$$
.

ii
$$A_1, A_2, \ldots \in \Sigma \implies \bigcup_{n \in \mathbb{N}} A_i \in \Sigma$$
.

Remark The smallest σ -algebra is $\{\emptyset, X\}$, whereas the biggest possible σ -algebra is the family of all subsets of X, $\mathscr{P}(X) = \{A : A \subseteq X\}$.

Remark In principle we need to consider σ -algebras which are strictly smaller than $\mathscr{P}(X)$, since a coherent definition of measure is very difficult to define on it.

Def. (σ -algebra generated by a set)

If $C \subseteq \mathscr{P}(X)$ is a family of subsets of X, then the σ -algebra generated by C is the smallest σ -algebra on X which contains C,

$$\sigma(C) = \bigcap_{i:C \subset \mathcal{F}} \mathcal{F}, \quad \mathcal{F} \text{ is a } \sigma\text{-algebra of } X.$$

Example (Borel σ -algebra)

If we consider $X = \mathbb{R}$, then the Borelian σ -algebra \mathcal{B} is the σ -algebra generated by all open intervals (a, b), $a, b \in \mathbb{R}$. That is, if $C = \{(a, b) : a < b, a, b \in \mathbb{R}\}$, then

$$\mathcal{B} = \sigma(C)$$
.

In any set X where we can define a topology, that is, a notion of open and closed sets, we can also define its associated Borel σ -algebra as the smallest σ -algebra that contains all open sets.

Under the above-defined σ -algebra, all intervals of the form

$$[a, b], [a, b), (a, b], (a, \infty), [a, \infty), (-\infty, b], (-\infty, b)$$

are all contained in \mathcal{B} . This is easy to prove by using the property of closeness w.r. to countable unions and complements. For instance,

$$(a, +\infty) = \bigcup_{i=1}^{\infty} (a, a+i) \in \mathcal{B}$$

$$\mathbb{R} \setminus (a, +\infty) = (-\infty, a] \in \mathcal{B}.$$

Moreover, \mathcal{B} can be equivalently defined as the σ -algebra generated by the following sets:

$$\sigma\big(\{(a,b) : a < b\}\big)$$

$$\sigma(\{(a,b] : a < b\})$$

$$\sigma(\{[a,b] : a < b\})$$

$$\sigma\big(\{[a,b) : a < b\}\big)$$

Remark The obvious question is: why would we need a notion of Borel σ -algebra when defining a measure on \mathbb{R} (and, by extension, on \mathbb{R}^d)? The reason is that we can construct some subsets A of \mathbb{R} such that $A \notin \mathcal{B}$, for example the Vitali set. On those pathological sets it is not possible to define a function such that it follows the properties that we expect from a *measure*.

Def. (Measure)

Let (X, Σ) be a measurable space, then we say that a function $\mu : \Sigma \to [0, +\infty]$ is a **measure** if μ is such that

$$\rightarrow \mu(\emptyset) = 0$$

$$\rightarrow (\sigma\text{-}additivity): \forall A_i \in \Sigma, A_i \cap A_j = \emptyset, \text{ then } \mu(\bigcup_{i=1}^{\infty} A_i) = \sum_{i=1}^{\infty} \mu(A_i).$$

In this case we say that the triple (X, Σ, μ) is a **measure space**.

Finiteness If μ is such that it assigns finite mass to X, i.e. $\mu(X) < \infty$, then μ is said to be a *finite measure*. A finite measure with $\mu(X) = 1$ is called a *probability measure*.

σ-finiteness If X can be written as a countable union of sets, $X = \bigcup_{i=1}^{\infty} A_i$ such that $\mu(A_i) < \infty$ for all i, then μ is σ-finite.

General properties of a measure Given a measure space (X, Σ, μ) , we can prove some general properties starting from the definition of μ :

- a) (Monotonicity): if $A, B \in \Sigma$ are such that $A \subseteq B$, then $\mu(A) \le \mu(B)$.
- b) $(\sigma$ -subadditivity): if $(A_i)_{i\in\mathbb{N}}$ is a sequence of elements of Σ , then

$$\mu\Big(\bigcup_{i=1}^{\infty} A_i\Big) \le \sum_{i=1}^{\infty} \mu(A_i).$$

c) (*Continuity*): if $(A_i)_{i\in\mathbb{N}}$ is a monotone increasing sequence of elements of Σ such that $A_1\subseteq A_2\subseteq A_3\subseteq\ldots$, then

$$\mu\Big(\bigcup_{i=1}^{\infty} A_i\Big) = \lim_{i \to \infty} \mu(A_i).$$

d) If $A_{i+1} \subseteq A_i$ at each i we have

$$\mu(\bigcap_{i} A_i) = \lim_{i \to \infty} \mu(A_i).$$

e) (Continuity ii): if $(A_i)_{i\in\mathbb{N}}$ is a monotone decreasing sequence of elements of Σ such that $A_{i+1}\subseteq A_i$ and $\mu(A_{i_0})<\infty$ for some i_0 , then

$$\mu\Big(\bigcap_{i=1}^{\infty} A_i\Big) = \lim_{i \to \infty} \mu(A_i).$$

1.2 Characterization of Borel measures

We try to characterize all measures on the Borel σ -algebra of \mathbb{R} , which are called *Borel measures*. To do so, let $F : \mathbb{R} \longrightarrow \mathbb{R}$ be an increasing function which is right-continuous, that is

$$\lim_{x \to a^{+}} F(x) = F(a).$$

Def. (Measure induced by F)

Let F be a function defined as above, we define the **measure induced by** F as the function μ_F such that for a set (a, b] we have

$$\mu_F(a, b] = F(b) - F(a).$$

For simplicity we define $\mu_F(\emptyset) = 0$.

σ-additivity When restricted to $C = \{(a, b] : a < b\}$ the measure μ_F is non-negative, additive, continuous w.r. to increasing sequences of sets. Therefore, we have that μ_F is also σ -additive on the σ -algebra \mathcal{B} .

Example (Measure of some particular sets)

We note that for some particular choices of sets,

$$\mu_F(\mathbb{R}) = \mu_F\left(\bigcup_n (a-n, b+n)\right) = \sup_x F(x) - \inf_x F(x)$$

$$\mu_F(\{a\}) = \mu_F(c, a) - \mu_F(c, a) = F(a) - \lim_{x \to a^-} F(x)$$

To complete the characterization of Borel measures we have the following theorem, which roughly states that defining a measure on \mathcal{B} can be done simply by defining a measure on \mathcal{C} .

Theorem 1 (Carathéodory existence theorem)

Let (X, \mathcal{B}) be a measurable space with Borel σ -algebra \mathcal{B} . Then,

- i. There exists a unique Borel measure $\bar{\mu}_F$ which coincides with μ_F on the intervals (a, b] and is σ -finite $\iff \sup_x F(x) \inf_x F(x) < \infty$.
- ii. Given a Borel measure μ on \mathbb{R} there exists a monotone increasing and right-continuous function F defined as

$$F(x) = \begin{cases} \mu(0, x] & \text{if } x > 0\\ -\mu(x, 0] & \text{if } x < 0 \end{cases}$$

such that μ is equal to the measure induced by F,

$$\mu = \mu_F$$
.

Proof.

No.

We now have the following relationship, which completely characterizes the Borel measures on \mathbb{R} .

Increasing right-continuous
$$F \iff$$
 Borel measure μ_F on \mathbb{R}

From this result we can now define the usual notion of measure (i.e. size) of sets in \mathbb{R} , which will be the building block for the formalization of the notion of integration over general measure spaces.

Def. (Lebesgue measure)

Let F(x) = x for all $x \in \mathbb{R}$, then if we define $\mu_F = \mu(a, b] = b - a \bar{\mu}_F$ its completion via the Carathéodory is called the **Lebesgue measure** on \mathbb{R} and we indicate it by \mathcal{L} .

Properties of the Lebesgue measure

- i.
- ii.
- iii.
- iv.
- v.

1.3 Decomposition of measures

Before introducing the Lebesgue integral, we now describe a useful characterization of the relationship between measures, which is useful to generalize the notion of a random variable from the dichotomy discrete/continuous to a more fundamental description.

Def. (Absolutely continuous measure)

Let ν, μ be measures defined on (X, Σ) , then we say that ν is **absolutely continuous with** respect to μ and we write $\nu \ll \mu$ if for every $A \in \Sigma$ it holds that

$$\mu(A) = 0 \implies \nu(A) = 0.$$

Prop. 1 (Absolutely continuous measure induced by a density)

Let $f \geq 0$ be a measurable function such that for all m > 0, $\int_{-m}^{m} f(x) dx < \infty$. Then, the function ν_f defined as

$$\nu_f(A) = \int_A f(x) \, \mathrm{d}x$$

is a measure on $(\mathbb{R}^n, \mathcal{M})$ which is both σ -finite and absolutely continuous w.r. to \mathcal{L} . if $f \in L^1(\mathbb{R}^n)$, then the measure is also finite.

Def. (Singular measure)

Let ν, μ be measures defined on (X, Σ) , then we say that ν is **singular with respect to** μ and we write $\nu \perp \mu$ if there exist $A, B \in \Sigma$ such that

$$A \cap B = \emptyset$$
 (disjoint)

$$A \cup B = X$$
 (partition)

$$\nu(A) = 0 = \mu(B)$$
 (measures are orthogonal)

Example (Dirac measure is singular)

Consider the Dirac measure δ_{x_0} centered on a point x_0 , then we can write

$$\mathbb{R} = \left(\underbrace{\mathbb{R} \setminus \{x_0\}}_{A} \cup \underbrace{\{x_0\}}_{B}\right),$$

and we have that $\mathcal{L}(B) = 0 = \delta_{x_0}(A)$.

We now state a fundamental theorem which completely characterizes the relationship of continuity and singularity between two measures.

Theorem 2 (Lebesgue decomposition)

Let ν, μ be two σ -finite measures on a measurable space (X, Σ) , then there exist two unique measures η (absolutely continuous part) and ρ (singular part) such that

$$\nu = \eta + \rho$$

$$\eta \ll \mu$$
,

$$\rho \perp \mu$$
.

LECTURE 2: LEBESGUE INTEGRATION

2021-10-26

We are now ready to define the notion of Lebesgue integration, which greatly extends integrability beyond the simpler Riemann integral.

2.1 Measurable functions

Def. (Measurable function)

Let (X, Σ) and (Y, \mathcal{E}) be two measurable spaces and let $f: X \longrightarrow Y$ be a function. We say that f is **measurable** with respect to \mathcal{E} and Σ if for all $E \in \mathcal{E}$ we have that

$$f^{-1}(E) \in \Sigma$$
.

Equivalent property For the special case of $f: X \to \mathbb{R}$ we are interested in the following equivalent condition:

$$f^{-1}(t, +\infty) = \{x \in X : f(x) > t\}$$
 is measurable.

Starting from the notion of measurable function, we are now ready to define what it means for a sequence of functions to converge.

Def. (Convergence in measure)

Let $(f_n)_{n\in\mathbb{N}}$ be a sequence of measurable functions and f be a measurable function, all defined on the measure space (X, Σ, μ) . Then, f_n converge to f in measure if for every $\varepsilon > 0$ we have

$$\lim_{n \to \infty} \mu(\lbrace x \in X : |f_n(x) - f(x)| \ge \varepsilon \rbrace) = 0.$$

Convergence in probability If $(\Omega, \mathcal{F}, \mathbb{P})$ is a probability space, this convergence is called *convergence in probability* since it reads

$$\lim_{n \to \infty} \mathbb{P}(\{\omega \in \Omega : |X_n(\omega) - X(\omega)| \ge \varepsilon\}) = 0.$$

2.2 Lebesgue integration

Both the definition of Lebesgue integration and the discussion of its properties are always developed in three stages. First, the definition of some property that holds for a simple class of functions; second, the simple functions are used to approximate some positive measurable function f of interest; finally, we can extend the property to any measurable function f by considering it as a sum of its positive and negative parts, f^+ and f^- .

Def. (Simple function)

Let A_1, A_2, \ldots, A_K be a finite family of disjoint sets and $c_1, c_2, \ldots, c_K > 0$ positive constants. We say that φ is a *simple function* if

$$\varphi(x) = \sum_{i=1}^{K} x_i \mathbb{1}_{A_i}(x).$$

Measurability of φ A simple function φ is a step function with a finite number of jumps, which can be immediately proven to be measurable: for a single indicator function,

$$t \ge 1 \longrightarrow \{x : \mathbb{1}_A(x) > t\} = \emptyset \in \mathcal{B}$$

$$t < 1 \longrightarrow \{x : \mathbb{1}_A(x) > t\} = A \in \mathcal{B}.$$

Then, for a simple function as defined above we can extend this argument by considering the indicator function of the union $\bigcup_{i=1}^{K} A_i$.

When considering a simple function we can give an intuitive definition of the Lebesgue integral, which represents the total area under the graph of the function.

Def. (Lebesgue integral of a simple function)

Let $\varphi(x) = \sum_{i=1}^{K} c_i \mathbb{1}_{A_i}(x)$, then we define the **Lebesgue integral of** φ as the linear functional

$$\int_{\mathbb{R}^N} \varphi(x) \, \mathrm{d}x = \sum_{i=1}^K c_i \mathcal{L}(A_i).$$

For a general measurable function f, we can extend the notion of the Lebesgue integral by approximating it from below with simple functions and taking the best approximation over all possible simple functions.

Def. (Lebesgue integral of a positive function)

Let $f: \mathbb{R}^n \longrightarrow R$ be a measurable function such that $f(x) \geq 0$, then we define its **Lebesgue** integral as

$$\int_{\mathbb{R}^n} f(x) \, \mathrm{d}x = \sup \Big\{ \int_{\mathbb{R}^n} \varphi(x) \, \mathrm{d}x : \varphi \text{ simple function and } \varphi < f \Big\}.$$

In the general case where f is not positive we define instead the integral in terms of its positive and negative parts, when such an operation is well-defined.

Def. (Lebesgue integral of a function)

Let $f: \mathbb{R}^n \longrightarrow R$ be a measurable function, possibly negative. Let now its **positive part** be $f^+(x) = \max\{0, f(x)\}$ and its **negative part** be $f^-(x) = \max\{0, -f(x)\}$. Then, we define the **Lebesgue integral of f** as

$$\int_{\mathbb{R}^n} f(x) dx = \int_{\mathbb{R}^n} f(x)^+ dx - \int_{\mathbb{R}^n} f(x)^- dx,$$

whenever this subtraction is well-defined.

Properties of the Lebesgue integral

f If f = 0 almost everywhere, then $\int_{\mathbb{R}^n} f(x) dx = 0$. Conversely, if $f \geq 0$ is measurable and $\int_{\mathbb{R}^n} f(x) dx = 0$ then f = 0 almost everywhere.

-) If f and g are measurable functions such that f=g almost everywhere, then $\int_{\mathbb{R}^n} f(x) \, \mathrm{d}x = \int_{\mathbb{R}^n} g(x) \, \mathrm{d}x$
- \rightarrow If $f,g\in L^1(\mathbb{R}^n)$, i.e. are absolutely integrable, then

$$\int_{\mathbb{R}^n} \left(\alpha f(x) + \beta g(x) \right) dx = \alpha \int_{\mathbb{R}^n} f(x) dx + \beta \int_{\mathbb{R}^n} g(x) dx.$$

 \rightarrow If $f,g\in L^1(\mathbb{R}^n)$ and $f\leq g$ almost everywhere, then

$$\int_{\mathbb{R}^n} f(x) \, \mathrm{d}x \le \int_{\mathbb{R}^n} g(x) \, \mathrm{d}x.$$

This definition of the integral is especially useful when dealing with sequences of functions for which we want to establish some sort of convergence of integrals.

Theorem 3 (Monotone convergence)

Let $f_k : \mathbb{R}^n \to \mathbb{R}$ be a sequence of measurable functions such that $f_k \geq 0$ and $f_k(x) \leq f_{k+1}(x)$ for all x and for all k, then

$$\lim_{k \to \infty} \int_{\mathbb{R}^n} f_k(x) \, \mathrm{d}x = \int_{\mathbb{R}^n} \lim_{k \to \infty} f_k(x) \, \mathrm{d}x.$$

Proof.

No.

Prop. 2 (Repartition function)

Let $f: \mathbb{R}^n \to \mathbb{R}$ be a measurable positive function, and let for every t > 0 F(t) be the repartition function of f,

$$F(t) = \mathcal{L}\{x : f(x) > t\}.$$

Then, it holds that

$$\int_{\mathbb{R}^n} f(x) \, \mathrm{d}x = \int_0^\infty F(t) \, \mathrm{d}t.$$

Proof.

No.

LECTURE 3: BANACH SPACES

2021-11-09

We start by defining some useful spaces in functional analysis, namely L^p and M^p spaces.

3.1 L^p spaces

Def. (L^p spaces)

For $p \in [1, \infty)$ we define the following vectorial space

$$L^p(A) = \Big\{ f : \mathbb{R}^n \to \mathbb{R} : f \text{ is measurable and } \int_A |f(x)|^p \, \mathrm{d}x < \infty \Big\}.$$

For $p = \infty$, we define

$$L^{\infty}(A) = \Big\{ f: R^n \to \mathbb{R}: f \text{ is measurable and } |f(x)| \le c \text{ for almost every } x \in A \Big\}.$$

Def. $(M^p \text{ spaces})$

When we consider $A = \Omega$, we have the space of random variables with finite p^{th} moment,

$$M^p = \Big\{ X: \Omega \to \mathbb{R}: X \text{ is measurable and } \mathbb{E}[|X|^p] < \infty \Big\},$$

with an analogous definition for $p = \infty$,

$$M^{\infty} = \Big\{ X : \Omega \to \mathbb{R} : X \text{ is measurable and } |X(\omega)| \le c \text{ for almost every } \omega \in \Omega \Big\}.$$

Convergence With these definitions, it's immediate to define a notion of convergence in terms of p spaces for general functions (L^p) and for random variables (M^p) :

 \rightarrow convergence in p-space for a sequence of functions $(f_n)_{n\in\mathbb{N}}$,

$$f_n \xrightarrow{L^p} f \iff \int_{\mathbb{R}} |f_n(x) - f(x)|^p dx \xrightarrow{n \to \infty} 0,$$

 \rightarrow convergence in p-mean for a sequence of random variables $(X_n)_{n\in\mathbb{N}}$

$$X_n \xrightarrow{L^p} X \iff \mathbb{E}[|X_n - X|^p] \xrightarrow{n \to \infty} 0.$$

These spaces are particularly important for functional analysis since they are examples of Banach spaces, whose structure we are going to study more generally.

We now state some inequalities which are useful for studying L^p and M^p spaces in their generality.

Def. (Conjugate exponent)

Let p > 1, then the **conjugate exponent** of p is q > 1 such that

$$\frac{1}{p} + \frac{1}{q} = 1 \iff q = \frac{p}{p-1}.$$

Moreover, if p=1 we say that its conjugate exponent is $q=+\infty$ and vice versa.

Theorem 4 (Young's inequality)

Let p, q be conjugate exponents, then for all x, y > 0 we have that

$$xy \le \frac{x^p}{p} + \frac{y^q}{q}.$$

Remark This property is a generalization of the classic identity $(a-b)^2 > 0 \implies ab < \frac{a^2}{2} + \frac{b^2}{2}$.

Theorem 5 (Hölder's inequality)

Let $p, q \in [1, \infty]$ be conjugate exponents, i.e. $\frac{1}{p} + \frac{1}{q} = 1$, then for all measurable functions f on $O \subseteq \mathbb{R}^n$ we have

$$||fg||_1 \le ||f||_p ||g||_q,$$

which if expanded becomes

$$\int_{O} |f(x)g(x)| \, \mathrm{d}x \le \left(\int_{O} |f(x)|^{p} \right)^{\frac{1}{p}} \left(\int_{O} |g(x)|^{q} \right)^{\frac{1}{q}}.$$

Moreover, if $f, g \in L^p(O)$ then $fg \in L^1(O)$ and this becomes an equality $\iff |f|^p$ and $|g|^q$ are linearly dependent in $L^1(O)$.

Corollary 1 (Minkowski's inequality)

Let $f,g \in L^p(O)$, then we have that $\|\cdot\|_p$ satisfies the triangle inequality, i.e.

$$||f + g||_p \le ||f||_p + ||g||_p,$$

which if expanded becomes

$$\left(\int_{O} |f(x) + g(x)|^{p} \right)^{\frac{1}{p}} \leq \left(\int_{O} |f(x)|^{p} \right)^{\frac{1}{p}} + \left(\int_{O} |g(x)|^{p} \right)^{\frac{1}{p}}.$$

 $\|\cdot\|_p$ is a norm This corollary states that $\|\cdot\|_p$ is indeed a norm on L^p , which is therefore a normed vector space.

3.2 Banach spaces

Def. (Distance induced by a norm)

Let $(X, \|\cdot\|)$ be a normed vectorial space, then we say that the *distance induced by the* **norm** $\|\cdot\|$ is the function

$$d(x, y) = ||x - y||, \quad x, y \in X.$$

Distance $d(\cdot, \cdot)$ as defined above is indeed a distance between two elements of X.

Topology This distance lets us define an induced topology on X – i.e. a notion of open and closed sets – firstly by defining the **balls of radius r centered in x_0** as

$$B(x_0, r) = \{x \in X : d(x_0, x) < r\},\,$$

and then defining a set $A \subseteq X$ as **open** if for all $x \in A$ there exists r > 0 such that $B(x,r) \subset A$.

Def. (Banach space)

Let $(x_n)_{n\in\mathbb{N}}$ be a Cauchy sequence, i.e. for all $\varepsilon>0$ there exists a $N\in\mathbb{N}$ such that $\|x_n-x_m\|<\varepsilon$ for all n,m>N. If all such sequences have limit in X, then X is called a **Banach space**.

Completeness This property is called *completeness*, and the definition states that "a Banach space is a complete normed vector space".

Theorem 6 (L^p are Banach spaces)

The spaces $L^p(\mathbb{R})$, $p \in [1, \infty]$ are Banach spaces w.r. to the distance induced by the norm

$$||f||_p = \left(\int_{\mathbb{R}} |f(x)|^p dx\right)^{\frac{1}{p}}, \quad p \in [1, \infty),$$

and

$$||f||_{\infty} = \sup_{x} |f(x)|.$$

 M^p spaces The theorem can be stated in terms of M^p being Banach spaces, when endowed with the norms

$$||X||_p = \mathbb{E}[|X|^p]^{\frac{1}{p}}, \quad p \in [1, \infty),$$

$$||X||_{\infty} = \sup_{\omega} |X(\omega)|, \quad p = \infty.$$

Proof.

Lecture 4: L^p spaces and bounded linear operators

2021-11-10

Recall that

$$M^p = \{X : \Omega \to \mathbb{R}, \text{ such that } \mathbb{E}[|X^p|] < \infty\}$$

$$L^p = \{ f : \mathbb{R} \to \mathbb{R}, \text{ such that } \int_{\mathbb{R}} |f(x)|^p \, \mathrm{d}x < \infty \}$$

These are actually the same spaces if you consider $M^p = L^p(\Omega)$, indeed you can see that

$$\mathbb{E}[|X|^p] = \int_{\Omega} |x|^p \, \mathrm{d}\mathbb{P}(\omega)$$

Prop. 3 (L^p spaces inclusion)

 $M^1 \supset M^2 \supset \ldots$, and in general $M^n \subset M^k$ if $n \geq k$. In particular, if X is a random variable such that $\mathbb{E}[|X|^k] < \infty$, then $\mathbb{E}[|X|^n] < \infty$ for all $n \leq k$.

Proof.

By Jensen's inequality, if $f: \mathbb{R} \longrightarrow \mathbb{R}$ convex (\Longrightarrow meas. and cont.) and X is a random variable, then the random variable f(X) is such that

$$\mathbb{E}[f(X)] \ge f(\mathbb{E}[X]).$$

We fix $n \ge k \ge 1$ and $X \in M^n \implies \mathbb{E}[|X|^n] < \infty$. We want to prove that also $X \in M^k$: we therefore fix $f(x) = |x|^{\frac{n}{k}}$, and we see that

$$\frac{n}{k} \ge 1 \implies f$$
 is a convex function.

Applying Jensen's inequality,

$$\mathbb{E}[f(Y)] \ge f(\mathbb{E}[Y])$$

we use the functions f and r.v. $Y = |X|^k$ to see that

$$\begin{cases} f(Y) = \left(|X|^k\right)^{\frac{n}{k}} = |X|^n \\ \mathbb{E}[f(Y)] = \mathbb{E}[|X|^n] < \infty & \text{by assumption} \end{cases}$$

and so we obtain $\infty > \mathbb{E}[|X|^n] \geq \mathbb{E}[|X|^k]^{\frac{n}{k}}$, by taking square roots we have

$$\infty > \mathbb{E}[|X|^n]^{\frac{1}{n}} \ge \mathbb{E}[|X|^k]^{\frac{1}{k}},$$

from which we conclude that $X \in M^n \implies X \in M^k$ for all $1 \le k \le n$ and moreover $||X||_k \le ||X||_n$.

Finally, we observe that this implies a relationship w.r. to convergence for all $1 \le k \le n$:

$$X_n \xrightarrow{M_n} X \implies X_n \xrightarrow{M_k} X$$
 for all $1 \le k \le n$.

 L^p spaces inclusion This proof is analogous when considering more general $L^p(A)$ spaces such that $\mathcal{L}(A) < \infty$, since M^p is a special case of a L^p space when choosing $A = \Omega$. Another proof can be obtained by applying Hölder's inequality:

Proof.

Consider $f \in L^n(A)$, we want to prove that $f \in L^k(A)$ for $1 \le k \le n$. Since $f \in L^n(A)$,

$$\int_{\mathbb{R}} |f(x)|^n \mathbb{1}_A(x) \, \mathrm{d}x < \infty,$$

which means that $f(x)\mathbb{1}_A(x) \in L^n(\mathbb{R})$. Moreover, its' trivial to see that $\mathbb{1}_A(x) \in L^p(\mathbb{R})$ for all p if $\mathcal{L}(A) < \infty$, since

$$\int_{\mathbb{R}} |\mathbb{1}_A(x)|^p \, \mathrm{d}x = \mathcal{L}(A).$$

Now, let $q = \frac{n}{n-1}$ be the conjugate exponent of n, we have that

$$\begin{cases} f \mathbb{1}_A \in L^n(\mathbb{R}) \\ \mathbb{1}_A \in L^{\frac{n}{n-1}}(\mathbb{R}) \end{cases}$$

then by Hölder's inequality, $(f \mathbb{1}_A) \cdot \mathbb{1}_A \in L^1(\mathbb{R})$ and we see that

$$||f||_{L^{1}(A)} = \int_{\mathbb{R}} |f(x)| \mathbb{1}_{A}(x) \, \mathrm{d}x \le ||f||_{L^{n}(A)} \left(\int_{\mathbb{R}} \mathbb{1}_{A}(x)^{\frac{n}{n-1}} \, \mathrm{d}x \right)^{\frac{n-1}{n}}$$
$$= ||f||_{L^{n}(A)} \cdot \left(\mathcal{L}(A) \right)^{\frac{n-1}{n}}$$

Therefore, we conclude that $L^n(A) \subseteq L^1(A)$ for any $n \ge 1$. With a similar argument, by accurately choosing the exponent k we also have that $f \in L^k(A)$ for all $1 \le k \le n$,

$$f \in L^k(A) \iff |f|^k \in L^1(A).$$

We have that

$$|f|^k \in L^{\frac{n}{k}}(A) \quad \frac{n}{k} > 1$$

since $\left(|f|^k\right)^{\frac{n}{k}}=|f|^n$ and by choosing the conjugate exponent of $\frac{n}{k},$

$$q = \frac{\frac{n}{k}}{\frac{n}{k} - 1}$$

Remark

- Recall that if $L(A) = \mathbb{R}$ then $\mathcal{L}(A) = \infty$ and the inclusion relationship is not true, since we proved before that $L^2(\mathbb{R}) \not\subseteq L^1(\mathbb{R})$.
- \rightarrow This holds in general for random variables since they induce a bounded measure, $\mathbb{P}(\Omega) = 1$.
- > We have that $L^{\infty}(A) \subseteq \bigcap_{k \ge 1} L^k(A)$ but it is not strictly equal to it.

We also have a way of computing the L^k norm in terms of the L^n norm:

$$\int_{A} |f|^{k} dx \le ||f||_{n}^{k} \mathcal{L}(A)^{\frac{n-k}{n}} \implies ||f||_{k} \le ||f||_{n} \cdot \mathcal{L}(A)^{\frac{n-k}{nk}}.$$

4.1 Operators between Banach spaces

Since we are working with spaces whose elements are functions (or random variables), we need to define a notion of *functional*, i.e. a transformation between two functions. The notion of a *linear functional* generalizes the linear transformation of standard algebra.

Def. (Linear operator)

Let X, Y be Banach spaces, then we say that a function $T: X \longrightarrow Y$ is a *linear operator* if T maintains the vectorial structure, i.e.

$$T(\alpha x_1 + \beta x_2) = \alpha T(x_1) + \beta T(x_2)$$
 for all $x_1, x_2 \in X$ and $\alpha_1, \alpha_2 \in \mathbb{R}$.

Def. (Continuity)

We say that $T: X \to Y$ is a **continuous** operator if for every sequence $x_n \in X$ such that $x_n \to x \in X$, then $Tx_n \to Tx$, i.e. converging sequences are mapped in converging sequences.

Def. (Boundedness)

We say that $T: X \to Y$ is **bounded** if there exists a constant C > 0 such that

$$||Tx||_Y < C||x||_X$$
 for all $x \in X$.

Remark Differently from standard functions $f : \mathbb{R} \to \mathbb{R}$ we don't require the image to be bounded, since such a condition turns out to be too restrictive for linear operators between Banach spaces.

Prop. 4 (Continuity implies boundedness)

If X, Y are Banach spaces and $T: X \longrightarrow Y$ is a linear operator, then T is continuous if and only if T is bounded.

Equivalence Because of this theorem, we will always talk about bounded operators instead of continuous operators, since they are the same.

Proof.

Bounded \Longrightarrow continuous: Let $x_n \in X$ such that $x_n \xrightarrow{n \to \infty} x \in X$, which by definition means that $||x_n - x||_X \xrightarrow{n \to \infty} 0$. From this, we want to prove that $Tx_n \xrightarrow{n \to \infty} Tx$.

Since the operator is bounded, there exists a C such that for all $x \in X$

$$||Tx||_Y \le C||x||_X,$$

and applying this to $y = x_n - x \in X$ we have

$$0 \le ||T(x_n - x)||_Y \le C||x_n - x||_X \xrightarrow{n \to \infty} 0,$$

and by linearity we have that $||Tx_n - Tx||_Y \to 0$.

Continuous \implies bounded: We actually prove another fact which implies this property, i.e. T not bounded \implies T not continuous, by constructing a sequence which is not converging under the map T. Assume that T is not bounded, then for any C>0 we can always find at least a point $x\in X$ such that

$$||Tx||_Y > C||x||_X,$$

In particular, let C = n, then we can always find a $x_n \in X$ such that

$$||Tx_n||_Y > n||x_n||_X$$
 for $||x_n||_X \neq 0$

Consider now the sequence of points $y_n = \frac{x_n}{n||x_n||_X}$, for which we have that $y_n \in X$ since X is a vectorial space. Then,

$$||y_n||_X = \left\|\frac{x_n}{n||x_n||_X}\right\| \stackrel{\text{norm.}}{=} \frac{1}{n}.$$

In particular, this means that $||y_n - 0||_X = \frac{1}{n} \xrightarrow{n \to \infty} 0$ and if T is continuous, then we should see that $Ty_n \to T0 = 0$ in Y since it is a linear operator. Now,

$$Ty_n = T\left(\frac{x_n}{n\|x_n\|_X}\right) \stackrel{\text{lin.}}{=} \frac{1}{n\|x_n\|_X} T(x_n),$$

and

$$||Ty_n||_Y = \left\| \frac{1}{n||x_n||_X} T(x_n) \right\|_Y \stackrel{\text{lin.}}{=} \frac{1}{n||x_n||_X} ||Tx_n||_Y \stackrel{\text{Hp.}}{>} \frac{1}{n||x_n||_X} n||x_n||_X = 1,$$

therefore Ty_n does not converge to 0 in Y and therefore we prove the desired property.

Theorem 7 (Space of linear operators)

Let X, Y be Banach spaces, then we have that the set of linear operators between X and Y,

$$\mathcal{B}(X,Y) = \{T : X \to Y \text{ linear bounded operators}\},\$$

is a Banach space with norm given by

$$||T|| = \sup_{||X||_X \le 1} ||Tx|| \stackrel{lin.}{=} \sup_{x \ne 0} \frac{||Tx||_Y}{||x||_X}.$$

Proof.

 \rightarrow Vectorial space: $\mathcal{B}(X,Y)$ is naturally a vectorial space, since $\alpha T + \beta S$ is a linear operator defined by

$$(\alpha T + \beta S)(x) := \alpha T(x) + \beta S(x).$$

Moreover, if T, S are bounded then $\alpha T + \beta S$ is also bounded by

$$\begin{split} \|\alpha Tx + \beta Sx\|_Y &\leq \alpha \|Tx\|_Y + \beta \|Sx\|_Y \leq \underbrace{\lfloor \alpha |\cdot C\|x\|_X}_{T \text{ bounded}} + \underbrace{\lfloor \beta |\cdot D\|x\|_X}_{S \text{ bounded}} \\ &= \left(|\alpha |C + |\beta |D\right) \|x\|_X. \end{split}$$

- \rightarrow Norm: ||T|| as defined in the theorem above is indeed a norm since it satisfies the three properties,
 - 1. $||T|| \ge 0$ for all T.
 - $2. \ \|\alpha T\| = \sup_{\|X\|_X \le 1} \|\alpha Tx\| = |\alpha| \cdot \sup_{\|x\| \le 1} \|Tx\| = |\alpha| \cdot \|T\|.$
 - $3. \ \|T+S\| = \sup_{\|x\| \le 1} \|Tx+Sx\|_Y \le \sup_{\|x\| \le 1} \|Tx\|_Y + \sup_{\|x\| \le 1} \|Sx\|_Y.$
- \rightarrow Completeness: We are not going to prove that $\mathcal{B}(X,Y)$ is a complete space w.r. to convergence induced by the norm, since it is a bit complex.

Another important result for linear bounded operators is the following, which is a theorem that can be used as an intermediate step in the proof of the completeness of $\mathcal{B}(X,Y)$.

Theorem 8 (Banach-Steinhaus)

If T_n is a sequence of bounded linear operators from X to Y and for every $x \in X$ there exists $T_n x$ in Y then the operator defined by

$$Tx := \lim_{n \to \infty} T_n x$$

is a bounded linear operator.

Example (Linear bounded operators)

Consider for a fixed $X \in M^p$ and $q = \frac{p}{p-1}$ the operator defined by

$$T: M^q \longrightarrow \mathbb{R}$$
$$Y \longmapsto \mathbb{E}[X \cdot Y].$$

We prove that for a fixed X, T is a linear bounded operator.

Linearity is immediate because of linearity of the $\mathbb{E}[\cdot]$ operator. As for boundedness we have to show that $\exists C > 0$ such that $\|Tx\|_{\mathcal{Y}} \leq C\|x\|_{\mathcal{X}}$, where $\mathcal{X} = M^q$ and $\mathcal{Y} = \mathbb{R}$. For all $Y \in M^q$ we therefore want to check that

$$|\mathbb{E}[X \cdot Y]| \stackrel{?}{\leq} C \cdot \mathbb{E}[\|Y\|^q]^{\frac{1}{q}}.$$

We can do so by applying Hölder's inequality, which allows us to write

$$|\mathbb{E}[X \cdot Y]| \overset{\mathrm{Jens.}}{\leq} \mathbb{E}[|X \cdot Y|] \overset{\mathrm{H\"old}}{\leq} \underbrace{\mathbb{E}[|X|^p]^{\frac{1}{p}}}_{C} \cdot \mathbb{E}[|Y|^q]^{\frac{1}{q}},$$

and the operator is bounded the constant $C = \mathbb{E}[|X|^p]^{\frac{1}{p}}$ (recall that X is fixed).

Moreover, we can actually prove (exercise) that the norm of T is equal to

$$||T|| = \mathbb{E}[|X|^p]^{\frac{1}{p}}$$

Example (Set of matrix operators)

If we consider $T: \mathbb{R}^n \to \mathbb{R}^m$, then the set of bounded linear operators is the set of operators defined by the $M_{m \times n}(\mathbb{R})$ matrices,

$$x \mapsto Tx = Ax, \quad A \in M_{m \times n}(\mathbb{R}).$$

Moreover, the norm ||T|| is connected to the norm of the matrix A.

Linear algebra From the example above, we can interpret the space $\mathcal{B}(X,Y)$ as the infinite-dimensional generalization of the space of matrices,

Space of
$$m \times n$$
 matrices $\xrightarrow{\text{Infinite dim.}} \mathcal{B}(X,Y)$,

and the results we are going to prove for infinite-dimensional linear operators are similar to those of standard finite-dimensional vector spaces.

LECTURE 5: HILBERT SPACES

2021-11-23

We now discuss Hilbert spaces, which are spaces where it is possible to define the notions of length and orthogonality. With these definitions, we are able to work with the elements of the space geometrically, as if they were vectors in Euclidean space.

5.1 Subspaces and spaces

Def. (Closed subspace)

X Banach space and $V \subseteq X$ is a **closed subspace of X** if

- 1. $V \subseteq X$ and V is a vectorial space.
- 2. If $f_n \in V$ for al n such that $f_n \to f$ in X, then $f \in V$.

V contains all the limit points of its converging sequences.

Example (M^2)

Let $V = \{X \in M^2 : \mathbb{E}[X] = 0\}.$

- 1. V is a subspace of M^2 , indeed if $X_1, X_2 \in V$ then also $\alpha X_1 + \beta X_2 \in V$.
- 2. Check that every converging sequence in V we have a limit in V itself.

 $X_n \xrightarrow{L^2} X$ is such that $\mathbb{E}[X_n] = 0$, since $X_n \in V$ and $\mathbb{E}[(X_n - X)^2] \to 0$ since $V \subseteq X$. We want to deduce that $\mathbb{E}[X] = 0$.

$$\mathbb{E}[|X_n - X|] \stackrel{prop. 3}{\leq} \sqrt{\mathbb{E}[(X_n - X)^2]},$$

therefore $\mathbb{E}[|X_n - X|] \to 0$, but now apply the expected value operator to all members of the following inequality to get convergence of X_n to X in V,

$$-|X_n - X| \le X_n - X \le |X_n - X|$$

Def. (Hilbert space)

Let X be a Banach space, then X is a **Hilbert space** if on X it is also defined a scalar product $\langle \cdot, \cdot \rangle : X \times X \longrightarrow \mathbb{R}$.

Properties of a scalar product

- 1. Positive: $\langle x, y \rangle \ge 0$ for all $x, y \in X$.
- 2. Symmetric: $\langle x, y \rangle = \langle y, x \rangle$ for all $x, y \in E$.
- 3. Bilinear: $\langle ax, y \rangle = a \langle x, y \rangle$ for all $x, y \in E$ and $a \in \mathbb{R}$..

Example (Scalar product on M^2)

If we consider M^2 , then $||X||_2 = \sqrt{\mathbb{E}[|X|^2]}$ and $X \cdot Y \in L^1$ by Hölder's inequality. Is the norm and the scalar product can be defined as

$$\langle X, Y \rangle = \mathbb{E}[X \cdot Y] \in \mathbb{R}.$$

Attention In general when $k \neq 2$, $X \in M^k$ and $Y \in M^k \implies XY \in M^1$.

Example (L^2)

Analogously, for L^2 we can define for $f, g \in L^2$ since $fg \in L^1$ the scalar product defined by

$$\langle f, g \rangle = \int_{\mathbb{R}} f(x)g(x) \, \mathrm{d}x.$$

Prop. 5 (Scalar product associated with a norm)

A norm $\|\cdot\|$ on a Banach space is associated with a scalar product $\langle\cdot,\cdot\rangle$, i.e. $\langle x,x\rangle=\|x\|^2\iff it\ satisfies\ the\ parallelogram\ identity,$

$$||x + y||^2 + ||x - y||^2 = 2||x||^2 + 2||y||^2, \quad \forall x, y \in X.$$

Theorem 9 (Cauchy-Schwartz inequality)

Let $x, y \in X$, then $|\langle x, y \rangle| \le ||x|| \cdot ||y||$.

Proof.

Consider for $a \in \mathbb{R}$ the norm

$$||x - ay||^2 = ||x||^2 + a^2 ||y||^2 - 2a\langle x, y \rangle,$$

and this is a parabola in a. Minimizing in a and substituting the minimum yields the required inequality.

Theorem 10 (Bicontinuity of $\langle \cdot, \cdot \rangle$)

Let $(x_n)_n$ and $(y_n)_n$ be sequences in a Hilbert space H such that $x_n \xrightarrow{n \to \infty} x$ and $y_n \xrightarrow{n \to \infty} y$. Then,

$$\lim_{m \to \infty} \langle x_n, y_m \rangle = \langle x, y \rangle.$$

Proof.

Consider $h_1, h_2 \in H$ and the scalar product

$$\begin{aligned} |\langle x+h_1,y+h_2\rangle - \langle x,y\rangle| &\stackrel{\text{bilin.}}{=} |\langle x,h_2\rangle + \langle h_1,y\rangle + \langle h_1,h_2\rangle| \\ &\leq \|x\| \|h_2\| + \|h_1\| \|y\| + \|h_1\| \|h_2\| \qquad \text{(Cauchy-Schwarz)} \\ &\frac{h_1,h_2 \to 0}{} 0. \end{aligned}$$

5.2 Orthogonality

With Hilbert spaces we can now define a notion of orthogonality between vectors $x, y \in H$.

Def. (Orthogonality)

 $x, y \in H$ are **orthogonal** if $\langle x, y \rangle = 0$.

Example $(L^2(-1,1))$

Consider $L^2(-1,1) = \{ f : \mathbb{R} \to \mathbb{R} \text{ such that } \int_{-1}^1 |f(x)|^2 dx < \infty \}.$

The functions $f(x) = \sin x$ and g(x) = 1 are orthogonal, since

$$\langle f, g \rangle = \int_{-1}^{1} \sin x \, \mathrm{d}x = 0.$$

Def. (Orthogonal complement)

Let $S \subseteq H$, we define the **orthogonal complement** S^{\perp} of S as the set of orthogonal elements to S,

$$S^{\perp} = \{ h \in H : \langle h, s \rangle = 0, \ \forall s \in S \}.$$

Subspace The orthogonal complement S^{\perp} can be shown to be a subspace of S since if $h_1, h_2 \in S^{\perp} \implies h_1 + h_2 \in S^{\perp}$ by bilinearity of the scalar product.

Closedness Moreover, S^{\perp} is also closed since by the Cauchy-Schwartz inequality (theorem 9) the scalar product is continuous w.r. to convergence of elements:

$$|\langle x_n, y \rangle - \langle x, y \rangle| = |\langle x_n - x, y \rangle| \le ||x_n - x|| \cdot ||y|| \xrightarrow{x_n \to x} 0.$$

Double complement We have that $(S^{\perp})^{\perp} = \{h \in H : \langle h, v \rangle, \ \forall v \in S^{\perp}\}$, therefore $S \subseteq (S^{\perp})^{\perp}$. However, $(S^{\perp})^{\perp}$ is a closed subspace of H, and if S is itself a closed subspace of H then $S = (S^{\perp})^{\perp}$. In general, $(S^{\perp})^{\perp}$ is the smallest closed subspace of H which contains S.

Def. (Direct sum)

A vector space H is said to be the **direct sum** of two subspaces V and W, and we write $H = V \bigoplus W$, if each $h \in H$ has a unique representation

$$h = v + w, \quad v \in V, w \in W.$$

Theorem 11 (Orthogonal projection)

Let H be a Hilbert space and let $V \subseteq H$ be a closed subspace of H. Then, $H = V \odot V^{\perp}$, in the sense that for each $h \in H$ there exists an unique $v \in V$ and an unique $w \in V^{\perp}$ such that

$$h = v + w$$
.

Orthogonal projection This theorem states that v is the *orthogonal projection* of h onto V, which is the element in V which has **minimal distance** from h,

$$v = \operatorname*{argmin}_{z \in V} \|h - z\|.$$

Proof.

We want to solve the minimization problem assuming $h \notin V$, otherwise v = h,

$$\min_{z \in V} \|z - h\|.$$

We define

$$\delta := \inf_{z \in V} \|h - z\|,$$

and we take a sequence $v_n \in V$ such that $\delta \leq ||h - v_n|| \leq \delta + \frac{1}{n}$. Since we are in a Hilbert space, using the parallelogram identity we have that

$$||v_n - v_m|| \xrightarrow{n,m \to \infty} 0,$$

therefore $v_n \xrightarrow{n \to \infty} v \in V$ since V is closed and

$$\delta \le \|h - v\| \le \delta.$$

Def. (Orthonogonal projection operator)

For all $V \subseteq H$ closed subset of H we define the **orthogonal projection operator** P_V as the operator that projects h = v + w onto the space V,

$$h \longmapsto P_V(h) = v.$$

Theorem 12 (Projection principle)

The orthogonal projection $v = P_V(h)$ is characterized by the tow following properties:

- 1. (Projection is in the subspace): $v \in V$.
- 2. $\langle v h, z \rangle = 0$ for all $z \in V$.

5.3 Random variables

Of particular interest is the case in which $(X, \mathcal{X}, \mu) = (\Omega, \mathcal{F}, P)$ with P probability measure, indeed $L^2(\Omega)$ is the space of (equivalence classes of) random variables such that

$$\mathbb{E}[|X|^2] < \infty,$$

with norm $\|X\| = \mathbb{E}[|X|^2]^{\frac{1}{2}}$, inner product $\langle X,Y \rangle = \mathbb{E}[XY]$ and distance between random variables given by $d(X,Y) = \mathbb{E}[|X-Y|^2]^{\frac{1}{2}}$. Since $L^2(\Omega)$ is a Hilbert space, in particular it is also complete and therefore if $(X_n)_n$ is a Cauchy sequence of random variables, i.e. such that

$$\lim_{m,n\to\infty} \mathbb{E}[|X_m - X_n|^2] = 0,$$

then $(X_n)_n$ converges in quadratic mean to $X \in L^2(\Omega)$.

LECTURE 6: HILBERT SPACE (CONT.)

2021-11-24

We are interested in finding a solution to the problem of computing the projection $P_V(x)$ of an element x of the Hilbert space X, where V is a closed subset of X.

6.1 Basis of a Hilbert space

Example (Conditional distribution)

Consider for instance the problem of finding the minimizer over

$$\mathbb{E}[(X - \mathbb{E}[X|V])^2] = \min_{z \in V} \mathbb{E}[(X - Z)^2],$$

therefore $\mathbb{E}[X|V]$ is the orthogonal projection of X over the subspace.

Example (Linear least squares estimation)

Given $X, Y \in M^2$ we want to find $a, b \in \mathbb{R}$ such that

$$\underset{a,b \in \mathbb{R}}{\operatorname{argmin}} \, \mathbb{E}[(Y - a - bX)^2].$$

We can rewrite this problem as an orthogonal projection over a Hilbert space M_2 given the closed subset $V = \{Y \in M_2 : Y = a + bX, \text{ for some } a, b \in \mathbb{R}\}$. Given $Y \in M_2$ we define the orthogonal projection of Y

$$P_V(Y) = \mathbb{E}[Y|V],$$

as the element in V with minimal distance from Y,

$$P_V(Y) = a + bX$$
, for some $a, b \in \mathbb{R}$,

such that

$$\mathbb{E}[(Y - P_V(Y))^2] = \min_{Z \in V} \mathbb{E}[(Y - Z)^2] = \min_{\alpha, \beta \in \mathbb{R}} \mathbb{E}[(Y - \alpha - \beta X)^2].$$

Therefore, the problem is how to compute the orthogonal projection on V.

We want a method to compute the orthogonal projection of elements in a Hilbert space when given a subspace V. In general, we can use two approaches in order to perform the required operation:

- 1. Write down the objective function and minimize it, as in the examples above.
- 2. Use the concept of a basis as in standard linear algebra.

Def. (Orthonormal set)

Let H be a Hilbert space, then a set $\{e_i, i \in \mathbb{N}\}$ is called an **orthonormal set** if

$$\langle e_i, e_j \rangle = \delta_{ij} = \begin{cases} 0 & \text{if } i \neq j \\ 1 & \text{if } i = j \end{cases}$$
 for any $i, j \in I$.

Example (\mathbb{R}^n)

In \mathbb{R}^n an orthonormal set is $\{(1,0,\ldots,0),(0,1,\ldots,0),\ldots,(0,0,\ldots,1)\}.$

Theorem 13 (Orthonormal basis)

Let H be a Hilbert space and $\{e_i, i \in I\}$ be an orthonormal set of elements of H. Then, the following conditions are equivalent

- 1. If for some $h \in H$ we have $\langle h, e_i \rangle = 0$ for all $i \in I$, then h = 0.
- 2. For every $h \in H$, then $h = \sum_{i \in I} \langle h, e_i \rangle e_i$.

Orthonormal basis If 1. or 2. hold, then we say that $\{e_i, i \in I\}$ is an *orthonormal basis* of H.

Def. (Fourier coefficients)

The elements $\langle h, e_i \rangle$ are called **Fourier coefficients** of x w.r. to the orthonormal sequence $(e_i)_{i \in I}$.

Lemma 1 (Fourier coefficients)

Any h in an inner product space H can have at most countably many nonzero Fourier coefficients $\langle h, e_i \rangle$ w.r. to a family of orthonormal elements $\{e_i, i \in I\}$. Hence, we can arrange the elements such that

$$h = \sum_{k=1}^{\infty} \langle h, e_k \rangle e_k,$$

and this sum does not depend on the order of the elements.

Approximation We have a direct way of approximating any element $h \in H$ in terms of the basis $\{e_i, i \in I\}$ by using the second property of the theorem. Indeed, for N fixed we have

$$\lim_{N \to \infty} \|h - \sum_{i=1}^{N} \langle h, e_i \rangle e_i \|,$$

and we have that the norm of h is equal to

$$||h||^2 = \langle h, h \rangle = \langle \sum_i \langle h, e_i \rangle e_i, \sum_j \langle h, e_j \rangle e_j \rangle = \sum_{i,j} \langle h, e_i \rangle \langle h, e_j \rangle \langle e_i, e_j \rangle = \sum_{i=j} \langle h, e_i \rangle \langle h, e_j \rangle.$$

Theorem 14 (Separability)

Every Hilbert space admits an orthonormal basis. Moreover, if the orthonormal basis is countable, $|I| = |\mathbb{N}|$, then H is called a **separable** Hilbert space.

Remark The sum of an uncountable set can be interpreted first by saying that $I = \mathbb{R}$ means

$$h = \sum_{i \in I} \langle h, e_i \rangle e_i \iff \langle h, e_i \rangle \neq 0 \text{ only for a countable } J \subseteq \mathbb{N} \text{ and } h = \sum_{j \in J} \langle h, e_i \rangle e_i.$$

Let $V \subseteq H$ be a subspace of H, then V is itself a Hilbert space with the same scalar product of H and it is closed. Then, we can find an orthonormal basis of V given by the orthonormal set of vectors $\{v_i : i \in I\}$ and for every $h \in H$ we have the orthogonal projection

$$P_V(h) = \sum_i \langle h, v_i \rangle v_i.$$

Example (Linear least squares)

back to our example, $V = \{\alpha + \beta X, \alpha, \beta \in \mathbb{R}\}$. Then, a basis of V is simply $\{1, X\}$, which however is not orthonormal. Firstly, we want them to be orthogonal and

$$\langle 1, X \rangle = \mathbb{E}[X] \neq 0$$
 in general,

therefore we replace X by $X - \mathbb{E}[X]$ and $\{1, X - \mathbb{E}[X]\}$ is still a basis of V. Finally, we normalize each vector by their norm

$$\|1\|_2 = \sqrt{\mathbb{E}[(1)^2]} = 1$$

$$\|X - \mathbb{E}[X]\|_2 = \sqrt{\mathbb{E}[(X - \mathbb{E}[X])^2]} = \sqrt{\mathbb{V}[X]}.$$

Therefore we finally obtain the orthonormal basis $\left\{1, \frac{X - \mathbb{E}[X]}{\sqrt{\|\mathbb{V}[X]\|}}\right\}$, and the orthogonal projection $P_V(Y)$ of Y onto V is simply

$$\begin{split} P_V(Y) &= \sum_{i=1}^2 \langle Y, v_i \rangle v_i \\ &= \langle Y, 1 \rangle \cdot 1 + \langle Y, \frac{Y - \mathbb{E}[X]}{\mathbb{V}[X]} \rangle \cdot \frac{X - \mathbb{E}[X]}{\mathbb{V}[X]} \\ &= \mathbb{E}[Y] + \frac{\mathbb{E}[XY] - \mathbb{E}[X]}{\sqrt{\mathbb{V}[X]}} \cdot \frac{X - \mathbb{E}[X]}{\sqrt{\mathbb{V}[X]}} \\ &= \mathbb{E}[Y] + \frac{\text{Cov}(X, Y)}{\mathbb{V}[X]} (X - \mathbb{E}[X]). \end{split}$$

6.2 Properties of linear operators

We now generalize to the infinite-dimensional case the properties of linear operators seen in linear algebra.

$$T: \mathbb{R}^n \longrightarrow \mathbb{R}^n \qquad \xrightarrow{n \to \infty} \qquad T: H \longrightarrow H$$
$$x \longmapsto Ax \qquad \qquad f \longmapsto Tf$$

Def. (Adjoint operator)

Let $T: H \longrightarrow H$ be a linear operator on a Hilbert space H, we say that $T^*: H \longrightarrow H$ is the *adjoint* operator of T if T^* is such that

$$\langle Tx, y \rangle = \langle x, T^*y \rangle$$
, for all $x, y \in H$.

Transpose This definition generalizes the transpose of a matrix A associated to a linear transformation, since for any matrix $A \in \mathcal{M}_{m \times n}(\mathbb{R})$,

$$\langle Ax, y \rangle = (Ax)^{\top} y = x^{\top} A^{\top} y = \langle x, A^{\top} y \rangle.$$

We consider the operator

$$T: L^2(0,1) \longrightarrow L^2(0,1)$$

$$f \longmapsto Tf(x) = \int_0^x f(t) dt.$$

Example

Firstly we have to prove that if $f \in L^2(0,1)$ then $Tf \in L^2(0,1)$ since

$$\int_{0}^{1} |Tf(t)|^{2} dt = \int_{0}^{1} \left| \int_{0}^{t} f(s) ds \right|^{2} dt$$

$$\leq \int_{0}^{1} t \int_{0}^{t} |f(s)|^{2} ds dt \qquad \text{(Jensen)}$$

$$\leq \int_{0}^{1} t \int_{0}^{1} |f(s)|^{2} ds dt \qquad \text{(0 } \leq s \leq t \leq 1\text{)}$$

$$\leq \int_{0}^{1} t ||f||_{2}^{2} dt$$

$$= \frac{1}{2} ||f||_{2}^{2}.$$

Therefore, $Tf \in L^2$ and t is linear (integral is linear) and bounded since by the computation above we have

$$||Tf||_2^2 \le \frac{1}{2}||f||_2^2.$$

In order to compute the adjoint we have to find a T^* such that

$$\langle Tf, g \rangle = \int_0^1 Tf(x)g(x) \, \mathrm{d}x = \int_0^1 f(x)T^*g(x) \, \mathrm{d}x = \langle f, T^*g \rangle.$$

The two terms have to be equal, i.e.

$$\int_0^1 f(x)T^*g(x) = \int_0^1 Tf(x)g(x) dx$$

$$= \int_0^1 \int_0^x f(t) dtg(x) dx \qquad (def.)$$

$$= \int_0^1 \int_t^1 f(t)g(x) dx dt \qquad (0 \le t \le x \le 1)$$

$$= \int_0^1 f(t) \left(\int_t^1 g(x) dx \right) dt$$

Therefore, we conclude that

$$T^*g(t) = \int_t^1 g(x) \, \mathrm{d}x.$$

Def. (Self adjoint – symmetric – operator)

If $T^* = T$ we say that T is a **self-adjoint** (or **symmetric**) operator.

Def. (Eigenfunction)

Given an operator $T: H \to H$, we say that $\lambda \in \mathbb{R}$ is an *eigenvalue* of T if there exists an $h \in H$ such that

$$Th = \lambda h$$
,

and h is called *eigenfunction* (or *eigenvector*).

Eigenvalues of T are contained in the so-called **spectrum** of T, that is

$$\{\lambda \in \mathbb{R} \text{ such that } (T - \lambda I) \text{ is not invertible}\},$$

where the operator $T - \lambda I$ is defined as

$$T - \lambda I : H \longrightarrow H$$

 $h \longmapsto Th - \lambda h.$

If λ is an eigenvalue of T then there exists $h \in H$ such that $(T - \lambda I)(h) = 0$. In the finite-dimensional case, the set of eigenvalues coincides with the spectrum, whereas in the infinite-dimensional case they form the **point spectrum** of T.

Def. (Compact linear operator)

Operators for which the spectrum is given by the set of eigenvalues are called *compact linear operators*.

Prop. 6 (Compact operators compress the space)

An operator $T: H \longrightarrow H$ is compact if and only if given any bounded sequence $h_n \in H$ not necessarily converging, then $Th_n \xrightarrow{n \to \infty} Th$.

Transition kernel A compact operator in L^2 is typically of the form of an integral w.r. to a transition kernel K,

$$f \longrightarrow Tf(x) = \int_A K(x, y) f(y) \, \mathrm{d}y,$$

and in the example above this was exactly

$$f \longrightarrow Tf(x) = \int_0^1 \mathbb{1}_{[0,x]}(y)f(y) \, \mathrm{d}y.$$

Symmetric compact operator are diagonalizable just like symmetric matrices in the finite-dimensional case.

Prop. 7 (Existence of an orthonormal basis)

If $T: H \longrightarrow H$ a is compact and symmetric operator on a separable Hilbert space H, then there exists an orthonormal basis $\{e_i, i \in I\}$ of H made of eigenfunctions of T, i.e.

$$Te_i = \lambda e_i$$
, for all $i \in I$.

Vanishing eigenvalues Moreover, the eigenvalues of T are such that $\{\lambda_i, i \in I\}$ can be either a finite set (i.e. some of them are equal) or an infinite set and

$$\lim_{i \to \infty} \lambda_i = 0.$$

Def. (Hilbert-Schmidt operator)

A compact and symmetric operator T is said to be a Hilbert-Schmidt (or trace-class) operator if T also satisfies the following condition,

$$\sum_{i=1}^{\infty} |\lambda_i|^2 < \infty.$$

Example (Continued)

Let $Tf(x) = \int_0^x f(t) dt$, then T is linear bounded from before. Moreover, it can be proven that T is compact but not symmetric, therefore we want to compute its eigenvalues if they exist: we want to find $\lambda \in \mathbb{R}$ such that

$$Tf(x) = \lambda f(x)$$
 for all x ,

therefore

$$\int_0^x f(t) dt = \lambda f(x) \iff f \text{ continuous and differentiable},$$

therefore if we differentiate both sides we find

$$\frac{\partial}{\partial x} \int_0^x f(t) dt = \frac{\partial}{\partial x} \lambda f(x) \iff \begin{cases} f(x) = \lambda f'(x) \\ 0 = \lambda f(0) \end{cases}$$

therefore $f(x) = e^{x/\lambda} \implies f'(x) = \frac{1}{\lambda} e^{x/\lambda}$. However, there $\lambda e^0 = \lambda$ but since $\lambda = 0$ we have f = 0 is not an eigenvalue. Therefore, T has no eigenvalues.

Theorem 15 (Riesz's representation)

Let $T: H \to H$ be a linear continuous operator on a Hilbert space H. Then there exists a unique $v \in H$ such that for all $h \in H$.

$$Th = \langle h, v \rangle,$$

and moreover ||T|| = ||v||.