Functional Analysis II<br>held by Prof. Dr. Moritz Weber and Dr. Tobias Mai Summer '18

## General information on organisation

Tutorials and admission for the final exam To take part in the final exam of this course, $50 \%$ of the total points on all exercise sheets have to be achieved. Furthermore there is compulsory attendence for the tutorials associated to this course, one shouldn't miss more than two tutorials without good reason.

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## 0 A brief reminder on Hilbert spaces and operators on Hilbert spaces

Throughout these lecture notes, the analytic closure in a topological $(X, \mathfrak{T})$ for $A \subseteq X$ will be denoted $\operatorname{cl}_{\mathfrak{T}}(A)$ or just $\operatorname{cl}(A)$, if no confusion is to be feared.
Definition 0.1: A pre-Hilbert space $(H,\langle\cdot, \cdot\rangle)$ is a $\mathbb{K}$-vector space $H$ (where $\mathbb{K}=\mathbb{R}$ or $\mathbb{K}=\mathbb{C}$ ) that is endowed with an inner product $\langle\cdot, \cdot\rangle$, i. e. a map

$$
\langle\cdot, \cdot\rangle: H \times H \longrightarrow \mathbb{K}
$$

that satisfies
(i) $\langle\lambda x+\mu y, z\rangle=\lambda\langle x, z\rangle+\mu\langle y, z\rangle$ for all $x, y, z \in H, \lambda, \mu \in \mathbb{K}$,
(ii) $\overline{\langle x, y\rangle}=\langle y, x\rangle$ for all $x, y \in H$,
(iii) $\langle x, x\rangle \geq 0$ for all $x \in H$,
(iv) If $\langle x, x\rangle=0$, then $x=0$.

A Hilbert space $(H,\langle\cdot, \cdot\rangle)$ is a pre-Hilbert space $(H,\langle\cdot, \cdot\rangle)$ that is complete (i. e. a Banach space) with respect to the norm

$$
\begin{aligned}
\|\cdot\|: H & \longrightarrow[0, \infty) \\
x & \longmapsto\langle x, x\rangle^{\frac{1}{2}}
\end{aligned}
$$

that is induced by the inner product.
Remark 0.2 (Properties of pre-Hilbert spaces): Let $(H,\langle\cdot, \cdot\rangle)$ be a pre-Hilbert space over $\mathbb{K}$.
(i) For all $x, y, z \in H, \lambda, \mu \in \mathbb{K}$, we have that

$$
\langle z, \lambda x+\mu y\rangle=\bar{\lambda}\langle z, x\rangle+\bar{\mu}\langle z, y\rangle .
$$

(ii) $\|\cdot\|$ is indeed a norm and the Cauchy-Schwarz inequality holds, i. e. for all $x, y \in H$ :

$$
|\langle x, y\rangle| \leq\|x\|\|y\| .
$$

(iii) $\|\cdot\|$ satisfies the parallelogram identity, i. e. for all $x, y \in H$ :

$$
\|x+y\|^{2}+\|x-y\|^{2}=2\left(\|x\|^{2}+\|y\|^{2}\right)
$$

In fact, a normed space is a pre-Hilbert space if and only if the norm satsifies the parallelogram identity. Indeed, the inner product can be recovered by the polarisation identities

$$
\langle x, y\rangle=\frac{1}{4} \sum_{k=0}^{3} \mathrm{i}^{k}\left\|x+\mathrm{i}^{k} y\right\|^{2} \quad(\text { if } \mathbb{K}=\mathbb{C})
$$

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$$
\left.\langle x, y\rangle=\frac{1}{4}\left(\|x+y\|^{2}-\|x-y\|^{2}\right) \quad \quad \text { (if } \mathbb{K}=\mathbb{R}\right)
$$

(iv) If $x, y \in H$ are orthogonal $(x \perp y)$, i.e. $\langle x, y\rangle=0$, then the Pythagorean identiy holds, i.e.

$$
\|x+y\|^{2}=\|x\|^{2}+\|y\|^{2} .
$$

(v) The completion of a pre-Hilbert space $(H,\langle\cdot, \cdot\rangle)$ is a Hilbert space.

Remark 0.3 (Properties of Hilbert spaces): Let $(H,\langle\cdot, \cdot\rangle)$ be a Hilbert space.
(i) If $C \subseteq H$ is closed and convex, then for each $x \in H \backslash C$ there is a unique point $y_{0} \in C$, called the best approximation to $x$ such that

$$
\left\|x-y_{0}\right\|=\operatorname{dist}(x, C):=\inf _{y \in C}\|x-y\|
$$

(ii) If $K \subseteq H$ is a closed linear subspace, then $y_{0} \in K$ is the best approximation to $x \in H \backslash K$ if and only if $\left\langle x-y_{0}, y\right\rangle=0$ for all $y \in K$.
(iii) Let $M \subseteq H$ be any subset. We call

$$
M^{\perp}:=\{x \in H \mid \forall y \in M:\langle x, y\rangle=0\}
$$

the orthogonal complement of $M$ in $H$. Note that $M^{\perp}$ is a closed linear subspace of $H$.
(iv) Let $K_{1}, K_{2} \subseteq H$ be closed linear subspaces with $K_{1} \perp K_{2}$, i. e. $\left\langle y_{1}, y_{2}\right\rangle=0$ for all $y_{1} \in K_{1}, y_{2} \in K_{2}$. Then the (orthogonal) direct sum

$$
K_{1} \oplus K_{2}:=\left\{y_{1}+y_{2} \mid y_{1} \in K_{1}, y_{2} \in K_{2}\right\}
$$

is again a closed linear subspace of $H$. Furthermore, $K_{1} \cap K_{2}=\{0\}$ and each $y \in K_{1} \oplus K_{2}$ has a unique decomposition $y=y_{1}+y_{2}$ with $y_{1} \in K_{1}$ and $y_{2} \in K_{2}$.
(v) The projection theorem (Theorem 5.18, Functional Analysis I) says that $H$ decomposes as $H=K \oplus K^{\perp}$ for any closed linear subspace $K \subseteq H$; in fact each $x \in H$ can be written as $x=x_{1}+x_{2}$, where $x_{1} \in K$ and $x_{2} \in K^{\perp}$ are the best approximations in $K$ and $K^{\perp}$ respectively to $x$.
(vi) The Riesz representation theorem (Theorem 5.20, Functional Analysis I) says that the dual space $H^{\prime}:=B(H, \mathbb{K})=\{f: H \rightarrow \mathbb{K}$ linear and bounded $\}$ can be identified with $H$ via the anti-linear isometric isomorphism

$$
\begin{aligned}
j: H & \longrightarrow H^{\prime} \\
y & \longmapsto f_{y}:=\langle\cdot, y\rangle .
\end{aligned}
$$

In particular $H \cong H^{\prime \prime}$, i. e. $H$ is reflexive.
(vii) A family $\left(e_{i}\right)_{i \in I}$ is called orthonormal system, if $\left\langle e_{i}, e_{j}\right\rangle=\delta_{i, j}$ for all $i, j \in I$. For such $\left(e_{i}\right)_{i \in I}$, Bessels inequality holds, i. e. for all $x \in H$ it holds

$$
\sum_{i \in I}\left|\left\langle x, e_{i}\right\rangle\right|^{2} \leq\|x\|^{2}
$$

with equality if and only if $x=\sum_{i \in I}\left\langle x, e_{i}\right\rangle e_{i}$ (Theorem 5.27, Functional Analysis I). We call $\left(e_{i}\right)_{i \in I}$ an orthonormal basis, if one of the following equivalent conditions is satisfied: ${ }^{1}$
(1) $\left(e_{i}\right)_{i \in I}$ is a maximal orthonormal system,
(2) If $x \perp e_{i}$ for all $i \in I$, then $x=0$,
(3) For all $x \in H$ it holds $x=\sum_{i \in I}\left\langle x, e_{i}\right\rangle e_{i}$,
(4) For all $x \in H$ it holds $\|x\|^{2}=\sum_{i \in I}\left|\left\langle x, e_{i}\right\rangle\right|^{2}$,
(5) The set $\left\{\sum_{i \in F} \alpha_{i} e_{i} \mid F \subseteq_{\text {fin }} I, \alpha_{i} \in \mathbb{K}\right\}$ is dense in $H$.

If $\left(f_{j}\right)_{j \in J}$ is another orthonormal basis of $H$, then $|I|=|J|$; the unique cardinality of an orthonormal basis of $H$ is called the (Hilbert space) dimension of $H$, denoted by $\operatorname{dim} H$. Every Hilbert space admits an orthonormal basis. Two Hilbert spaces $H$ and $K$ are isomorphic $(H \cong K)$, i. e. there is a sujective linear map $U: H \rightarrow K$ that satisfies

$$
\langle U x, U y\rangle_{K}=\langle x, y\rangle_{H} \quad \text { for all } x, y \in H
$$

if and only if $\operatorname{dim} H=\operatorname{dim} K$.
If $\operatorname{dim} H$ is countable, we call $H$ separable; this is equivalent to $H$ being separable as a Banach space (i.e. there is a countable dense subset).

Example 0.4: (i) $\mathbb{K}^{n}$ is a Hilbert space with the inner product

$$
\langle x, y\rangle=\sum_{i=1}^{n} x_{j} \bar{y}_{j} \quad \text { for } x=\left(x_{1}, \ldots, x_{n}\right)^{\top}, y \in\left(y_{1}, \ldots, y_{n}\right)^{\top} \in \mathbb{K}^{n} .
$$

(ii) If $(\Omega, \mathfrak{F}, \mu)$ is a measure space, then

$$
L^{2}(\Omega, \mu):=\left\{f: \Omega \rightarrow \mathbb{C} \text { measurable }: \int_{\Omega}|f(\omega)|^{2} d \mu(\omega)<\infty\right\} / \mathcal{N}
$$

where $\mathcal{N}:=\{f: \Omega \rightarrow \mathbb{C}$ measurable $\mid \mu(\{\omega \in \Omega: f(\omega) \neq 0\})=0\}$, is a Hilbert space with inner product

$$
\langle f, g\rangle=\int_{\Omega} f(\omega) \overline{g(\omega)} d \mu \quad \text { for all } f, g \in L^{2}(\Omega, \mu)
$$

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In particular, if a set $I$ is endowed with the counting measure $\sigma_{I}$, we obtain the Hilbert space $\ell^{2}(I):=L^{2}\left(I, \sigma_{I}\right)$ with $\operatorname{dim} \ell^{2}(I)=|I|$. Each separable Hilbert space $H$ that is not finite dimensional satisfies $H \cong \ell^{2}(\mathbb{N})$.

An orthonormal basis of $L^{2}\left([0,2 \pi), \lambda^{1}\right)$ with the Lebesgue measure $\lambda^{1}$ on $[0,2 \pi)$ is given by $\left(e_{n}\right)_{n \in \mathbb{Z}}$ where

$$
e_{n}(t)=\frac{1}{\sqrt{2 \pi}} e^{\mathrm{i} n t} \quad \text { for } t \in[0,2 \pi)
$$

For $f \in L^{2}\left([0,2 \pi), \lambda^{1}\right)$, the representation $f=\sum_{n \in \mathbb{Z}}\left\langle f, e_{n}\right\rangle e_{n}$ (in the $L^{2}$-sense) is called the Fourier series of $f$ and

$$
\hat{f}_{n}:=\frac{1}{\sqrt{2 \pi}}\left\langle f, e_{n}\right\rangle=\frac{1}{2 \pi} \int_{0}^{2 \pi} f(t) e^{-\mathrm{i} n t} d \lambda^{1}(t)
$$

are called the Fourier coefficients; in particular $\|f\|^{2}=2 \pi \sum_{n \in \mathbb{Z}}\left|\hat{f}_{n}\right|^{2}$.
(iii) Hardy spaces and Bergman spaces are Hilbert spaces of holomorphic functions, especially so-called reproducing Hilbert spaces.

Remark 0.5: Let $X$ be a vector space over $\mathbb{K}$.
(i) A map $p: X \rightarrow[0, \infty)$ is called a seminorm, if it satisfies the following statements:

- For all $x \in X$ and $\lambda \in \mathbb{K}$ it holds $p(\lambda x)=|\lambda| p(x)$,
- For all $x, y \in X$ it holds $p(x+y) \leq p(x)+p(y)$

For $x \in X$ and $r>0$, we put $B_{p}(x, r):=\{y \in X \mid p(y-x)<r\}$.
(ii) We call $X$

- topological vector space, if $X$ is endowed with a topology $\mathfrak{T}$ with respect to which

$$
\begin{aligned}
+: X \times X & \longrightarrow X & \cdot: \mathbb{K} \times X & \longrightarrow X \\
(x, y) & \longmapsto x+y & (\alpha, x) & \longmapsto \alpha x
\end{aligned}
$$

are both continuous.

- locally convex vector space, if $X$ is a topological vector space whose topology $\mathfrak{T}$ is generated by some family $P$ of seminorms, i. e. $U \subseteq X$ is open in $(X, \mathfrak{T})$ if and only if

$$
\forall x \in U \exists n \in \mathbb{N} \exists p_{1}, \ldots, p_{n} \in P \exists \varepsilon_{1}, \ldots, \varepsilon_{n}>0: \bigcap_{i=1}^{n} B_{p_{i}}\left(x, \varepsilon_{i}\right) \subseteq U .
$$

That topology is Hausdorff if and only if for all $0 \neq x \in X$, there is $p \in P$ such that $p(x)>0$.

Definition 0.6: Let $(H,\langle\cdot, \cdot\rangle)$ be a Hilbert space over $\mathbb{K}$.
(i) The topology induced by $\|\cdot\|$ is called the strong topology on $H$.
(ii) The locally convex (Hausdorff) topology induced by the family $P=\left\{p_{x} \mid\right.$ $x \in H\}$ of seminorms

$$
\begin{aligned}
p_{x}: H & \longrightarrow[0, \infty) \\
y & \longmapsto|\langle y, x\rangle|,
\end{aligned}
$$

is called the weak topology on $H$.
Theorem 0.7: Let $(H,\langle\cdot, \cdot\rangle)$ be a Hilbert space over $\mathbb{K}$.
(i) If a sequence $\left(x_{n}\right)_{n \in \mathbb{N}}$ in $H$ converges weakly to some point $x \in H$ and $\left\|x_{n}\right\| \rightarrow\|x\|$ as $n \rightarrow \infty$, then $\left(x_{n}\right)_{n \in \mathbb{N}}$ converges strongly to $x$.
(ii) Every bounded sequence in $H$ has a weakly convergent subsequence.
(iii) Every bounded sequence $\left(x_{n}\right)_{n \in \mathbb{N}}$ has a weakly convergent subsequence $\left(x_{n_{k}}\right)_{k \in \mathbb{N}}$ such that

$$
\left(\frac{1}{K} \sum_{k=1}^{K} x_{n_{k}}\right)_{K \in \mathbb{N}}
$$

converges strongly (Theorem of Banach-Saks).
Proof: Part (ii) and (iii) are exercises, we want to show part (i). We have

$$
\left\|x-x_{n}\right\|^{2}=\|x\|^{2}-2 \operatorname{Re}\left(\left\langle x_{n}, x\right\rangle\right)+\left\|x_{n}\right\|^{2}
$$

and by weak convergence of $\left(x_{n}\right)_{n \in \mathbb{N}}$ we also have

$$
\left|\left\langle x_{n}, x\right\rangle-\|x\|^{2}\right|=\left|\left\langle x_{n}-x, x\right\rangle\right|=p_{x}\left(x_{n}-x\right) \longrightarrow 0 \text { as } n \rightarrow \infty,
$$

i. e. $\left\langle x_{n}, x\right\rangle \rightarrow\|x\|^{2}$ as $n \rightarrow \infty$, hence $\operatorname{Re}\left(\left\langle x_{n}, x\right\rangle\right) \rightarrow\|x\|^{2}$ as $n \rightarrow \infty$. Thus it holds $\left\|x-x_{n}\right\| \rightarrow 0$ as $n \rightarrow \infty$.

Remark 0.8: (i) One can show that the closed unit ball $\{x \in H \mid\|x\| \leq 1\}$ is compact with respect to the weak topology.
(ii) If the Hilbert space is separable, then the weak topology on $\{x \in H \mid\|x\| \leq 1\}$ is metrisable.

Remark 0.9: (i) Let $\left(H,\langle\cdot, \cdot\rangle_{H}\right)$ and $\left(K,\langle\cdot, \cdot\rangle_{K}\right)$ be Hilbert spaces over $\mathbb{K}$. For each $A \in B(H, K)$, there is a unique operator $A^{*} \in B(K, H)$ with

$$
\langle A x, y\rangle_{K}=\left\langle x, A^{*} y\right\rangle_{H} \quad \text { for all } x \in H, y \in K
$$

We call $A^{*}$ the adjoint operator to $A$; with $A^{\prime} \in B\left(K^{\prime}, H^{\prime}\right)$ defined via $A^{\prime} f:=f \circ A$, it is given as $A^{*}=j_{H}^{-1} \circ A^{\prime} \circ j_{K}$ and we have that $\left\|A^{*}\right\|=\|A\| .^{2}$

[^1](ii) Let $(H,\langle\cdot, \cdot\rangle)$ be a Hilbert space over $\mathbb{C}$. Then $B(H):=B(H, H)$ is a Banach algebra with respect to the operator norm $\|\cdot\|$, in particular $\|S T\| \leq\|S\|\|T\|$ for all $S, T \in B(H)$. The map * : $B(H) \rightarrow B(H), A \mapsto A^{*}$ satisfies

- $(\lambda A+\mu B)^{*}=\bar{\lambda} A^{*}+\bar{\mu} B^{*}$,
- $(A B)^{*}=B^{*} A^{*}$,
- $\left(A^{*}\right)^{*}=A$,
- $\left\|A^{*}\right\|=\|A\|$ and $\left\|A^{*} A\right\|=\|A\|^{2}$.
(iii) Note that $\operatorname{ker}(A)=\operatorname{im}\left(A^{*}\right)^{\perp}$ for all $A \in B(H)$.

Definition 0.10: Let $\left(H,\langle\cdot, \cdot\rangle_{H}\right)$ and $\left(K,\langle\cdot, \cdot\rangle_{K}\right)$ be Hilbert spaces.
(i) $A \in B(H)$ is called selfadjoint (or hermitian), if $A=A^{*}$.
(ii) $A \in B(H)$ is called normal, if $A A^{*}=A^{*} A$.
(iii) $A \in B(H)$ is called positive, if $A=A^{*}$ and $\langle A x, x\rangle \geq 0$ for all $x \in H$.
(iv) $V \in B(H, K)$ is called isometry if one of the following equivalent statements hold:

- $V^{*} V=\operatorname{id}_{H}$,
- $\|V x\|_{K}=\|x\|_{H}$ for all $x \in H$,
- $\langle V x, V y\rangle_{K}=\langle x, y\rangle_{H}$ for all $x, y \in H$.
(v) $U \in B(H, K)$ is called unitary if $U^{*} U=\mathrm{id}_{H}$ and $U U^{*}=\mathrm{id}_{K}$ or equivalently if $U$ is a surjective isometry.
(vi) $P \in B(H)$ is called (orthogonal) projection if $P^{2}=P=P^{*}$. Then im $(P)$ is a closed linear subspace of $H$ (in fact $\operatorname{im}(P)=\operatorname{ker}(1-P)$ where $\left.1=\mathrm{id}_{H}\right)$ and $H=\operatorname{im}(P) \oplus \operatorname{ker}(P)$.

Conversely: If $K \subseteq H$ is a closed linear subspace, then

$$
\begin{aligned}
P: H=K \oplus K^{\perp} & \longrightarrow H \\
x=x_{1}+x_{2} & \longmapsto x_{1}
\end{aligned}
$$

is an orthogonal projection; note that $P x$ is the best approximation to $x$ in $K$.
(vii) $V \in B(H)$ is called partial isometry if one of the following equivalent conditions holds:

- $V V^{*} V=V$,
- $V^{*} V$ is a projection (initial projection),
- $V V^{*}$ is a projection (final projection),
- There is a closed linear subspace $K \subseteq H$ such that $\left.V\right|_{K}: K \rightarrow H$ is an isometry and $\left.V\right|_{K^{\perp}} \equiv 0 .^{3}$

[^2]Theorem 0.11 (Polar decomposition): For each $T \in B(H)$, there is a unique partial isometry $V \in B(H)$ such that $T=V|T|$ and $\operatorname{ker}(V)=\operatorname{ker}(T)$ where we define $|T|:=\left(T^{*} T\right)^{\frac{1}{2}}$.

In fact, $V^{*} V$ is the projection onto $\operatorname{ker}(T)^{\perp}$ and $V V^{*}$ is the projection onto $\operatorname{cl}(\operatorname{im}(T))^{4}$

[^3]
## 1 Locally convex topologies on the space of operators

Let $(H,\langle\cdot, \cdot\rangle)$ be a Hilbert space over $\mathbb{K}$.
Definition 1.1: (i) The weak operator topology (WOT) on $B(H)$ is the locally convex topology that is defined by the family of seminorms

$$
\begin{aligned}
p_{x, y}: B(H) & \longrightarrow[0, \infty) \\
T & \longmapsto|\langle T x, y\rangle| \quad \text { for all } x, y \in H .
\end{aligned}
$$

(ii) The strong operator topology (SOT) on $B(H)$ is the locally convex topology induced by the family of seminorms

$$
\begin{aligned}
p_{x}: B(H) & \longrightarrow[0, \infty) \\
T & \longmapsto\|T x\|
\end{aligned} \quad \text { for all } x, y \in H .
$$

(iii) The operator norm topology (ONT) on $B(H)$ is the topology induced by the operator norm $\|\cdot\|$.

Remark 1.2: We have $\mathfrak{T}_{\text {WOT }} \subseteq \mathfrak{T}_{\text {SOT }} \subseteq \mathfrak{T}_{\text {ONT }}$; note that

$$
|\langle T x, y\rangle| \leq\|T x\|\|y\| \leq\|T\|\|x\|\|y\| \quad \forall x, y \in H, T \in B(H)
$$

Remark 1.3: (i) $T \mapsto T^{*}$ is continuous with respect to ONT and WOT, but is not continuous with respect to SOT.
(ii) Multiplication : : $B(H) \times B(H) \rightarrow B(H),(S, T) \mapsto S T$ is continuous with respect to ONT but not continuous with respect to SOT or WOT.
(iii) However, for fixed $S \in B(H)$, both mappings $T \mapsto S T$ and $T \mapsto T S$ are continuous with respect to SOT and WOT.
(iv) Furthermore: If $\left(S_{n}\right)_{n \in \mathbb{N}},\left(T_{n}\right)_{n \in \mathbb{N}}$ are sequences in $B(H)$ that are strongly convergent to $S$ and $T$ respectively, then $\left(S_{n} T_{n}\right)_{n \in \mathbb{N}}$ converges strongly to $S T$. Indeed, $\left(\left\|S_{n}\right\|\right)_{n \in \mathbb{N}},\left(\left\|T_{n}\right\|\right)_{n \in \mathbb{N}}$ are bounded by the uniform boundedness principle, so that

$$
\left\|S_{n} T_{n} x-S T x\right\| \leq\left\|S_{n}\right\|\left\|\left(T_{n}-T\right) x\right\|+\left\|\left(S_{n}-S\right) x\right\| \longrightarrow 0 \text { as } n \rightarrow \infty
$$

If now the strong - is replaced by the weak topology, the statement is not valid anymore.

Operator algebras: ONT $\rightsquigarrow C^{*}$-algebras, WOT, SOT $\rightsquigarrow$ von Neumann algebras.

## 2 Unitisation of $C^{*}$-algebras

Definition 2.1: (i) A $C^{*}$-algebra is a Banach algebra $A$ (i.e. $A$ is a $\mathbb{C}$-algebra that is complete and normed such that $\|x y\| \leq\|x\|\|y\|$ for all $x, y \in A$ ) with an involution, i.e. an anti-linear map

$$
\begin{aligned}
*: A & \longrightarrow A \\
x & \longmapsto x^{*}
\end{aligned}
$$

such that $x^{* *}=x,(x y)^{*}=y^{*} x^{*}$ and $\left\|x^{*} x\right\|=\|x\|^{2}$.
(ii) $A$ is unital, if $A$ is unital as an algebra $(1 \in A)$.

Motivation 2.2: (i) So far, we only considered unital $C^{*}$-algebras such as

$$
B(H):=\{T: H \longrightarrow H \text { linear, bounded }\}
$$

for instance $M_{n}(\mathbb{C})=B\left(\mathbb{C}^{n}\right)$ or $C(X):=\{f: X \rightarrow \mathbb{C}$ continuous $\}$ for a compact metric space $X$. These are natural as well as unital examples and we have strong theorems for them:

Theorem (First funamental theorem for $C^{*}$-algebras): Let $A$ be a commutative and unital $C^{*}$-algebra. Then $A \cong C(X)$ for some compact metric space $X$.

In the proof of the Stone-Weierstraß Theorem, we also needed the unit:

$$
\sqrt{f}=\sqrt{1-g}=1-\sum_{n=1}^{\infty} a_{n} g^{n} \leftarrow 1-\sum_{n=1}^{N} a_{n} g^{n} \in A
$$

for $0 \leq f \leq 1, f \in A$ and $g:=1-f$.
(ii) The compact operators $K(H):=\{T: H \rightarrow H$ compact $\} \subseteq B(H)$ satisfy Definition 2.1 (i), but not (ii). One checks, that $K(H)$ is infact a $C^{*}$-algebra.

More general: If $A$ is a $C^{*}$-algebra and $I \triangleleft A$ an ideal (two-sided and closed, $I^{*} \subseteq I$ ), then $I$ is a (non-unital) $C^{*}$-algebra. Hence, there are important examples of non-unital $C^{*}$-algebras.
(iii) Another important example (for a unital $C^{*}$-algebra) was $C(X)$ for a compact metric space $X$. Why "compact"?

$$
\|f\|_{\infty}:=\sup \{|f(t)| \mid t \in X\}<\infty \text { if } X \text { is compact. }
$$

Let now $X$ be locally compact (i.e. for all $x \in X$ and $U \subseteq X$ open with $x \in U$, there is a compact neighbourhood $K \subseteq X$ such that $x \in K \subseteq U$. Then $\|f\|_{\infty}=\infty$ for $f \in C(X)$ is possible.

Consider

$$
\begin{aligned}
& C_{0}(X) \\
& :=\{f: X \rightarrow \mathbb{C} \text { continuous, vanishing at } \infty\} \\
& =\{f: X \rightarrow \mathbb{C} \text { continuous, } \forall \varepsilon>0 \exists K \subseteq X \text { compact }:|f(t)|<\varepsilon \forall t \in X \backslash K\},
\end{aligned}
$$

then $\|f\|_{\infty}<\infty$ for all $f \in C_{0}(X)$. One can check that $\left(C_{0}(X),\|\cdot\|_{\infty}\right)$ is a commutative (non-unital) $C^{*}$-algebra.
(iv) We shall find ways of dealing with non-unital $C^{*}$-algebras, for instance by "adding a unit": a minimal unit, a maximal unit or an approximitive unit. Amongst other, we need to study how to find $C^{*}$-norms.

Proposition 2.3: Let $B$ be a unital $C^{*}$-algebra.
(i) Let $x \in B$ with $x=x^{*}$. Then $r(x)=\|x\| .{ }^{1} \quad$ ("the norm is algebraic")
(ii) The $C^{*}$-norm on $B$ is unique.
(iii) If $A$ is a unital ${ }^{*}$-Banach algebra such that $\left\|x^{*}\right\|=\|x\|$ and $\|1\|=1$ and $\varphi: A \rightarrow B$ is a unital ${ }^{*}$-homomorphism (i.e. $\varphi: A \rightarrow B$ is an algebrahomomorphism that satisfies $\varphi(x)^{*}=\varphi\left(x^{*}\right)$ ), then $\|\varphi(x)\| \leq\|x\|$.
(iv) On a unital ${ }^{*}$-Banach algebra, the $C^{*}$-norm (if it exists) is the minimal norm.

Proof: (i) From Theorem 8.13, Functional Analysis I, we know that

$$
r(x)=\lim _{n \rightarrow \infty} \sqrt[n]{\left\|x^{n}\right\|}=\lim _{n \rightarrow \infty} \sqrt[2^{n}]{\left\|x^{2^{n}}\right\|}=\lim _{n \rightarrow \infty} \sqrt[2^{n}]{\|x\|^{2^{n}}}=\|x\|
$$

as $\left\|x^{2^{n+1}}\right\|=\left\|\left(x^{2^{n}}\right)^{*}\left(x^{2^{n}}\right)\right\|=\left\|x^{2^{n}}\right\|^{2}=\|x\|^{2^{n+1}}$ via induction.
(ii) Let $\|\cdot\|_{1},\|\cdot\|_{2}$ be $C^{*}$-norms on $B$. Then

$$
\|x\|_{1}^{2}=\left\|x^{*} x\right\|_{1} \stackrel{(\mathrm{i})}{=} r\left(x^{*} x\right)=\left\|x^{*} x\right\|_{2}=\|x\|_{2}^{2}
$$

(iii) We can estimate

$$
\|\varphi(x)\|_{B}^{2}=\left\|\varphi\left(x^{*} x\right)\right\|_{B} \stackrel{(\mathrm{i})}{=} r\left(\varphi\left(x^{*} x\right)\right) \leq r\left(x^{*} x\right) \leq\left\|x^{*} x\right\|_{A} \leq\left\|x^{*}\right\|_{A}\|x\|_{A}=\|x\|_{A}^{2}
$$

because if $\lambda-\varphi(y)=\varphi(\lambda 1-y)$ is not invertible, then $\lambda 1-y$ is not invertible: If $z$ is invertible, then there is $z^{-1}$ such that $z z^{-1}=1$, thus $\varphi(z) \varphi\left(z^{-1}\right)=1$.
(iv) The map

$$
\begin{aligned}
\varphi:\left(A,\|\cdot\|_{\text {Banach norm }}\right) & \longrightarrow\left(A,\|\cdot\|_{C^{*}-\text { algebra }}\right) \\
x & \longmapsto x
\end{aligned}
$$

is a ${ }^{*}$-homomorphism, hence $\|x\|_{C^{*} \text {-algebra }}=\|\varphi(x)\| \leq\|x\|_{\text {Banach norm }}$.

[^4]Reminder 2.4 (Theorem of Gelfand-Naimark ${ }^{2}$ ): Let $A$ be a commutative unital $C^{*}$-algebra. Then the Gelfand transformation

$$
\begin{array}{rlrl}
\chi: A & \longrightarrow C(\operatorname{Spec}(A)) & & (\operatorname{Spec}(A):=\{\varphi: A \rightarrow \mathbb{C} \text { unital algebra homomorphism }\}) \\
x & \longmapsto \hat{x} & (\hat{x}(\varphi):=\varphi(x))
\end{array}
$$

is an isometric *-isomorphism.
$\operatorname{Spec}(A)$ is compact, refer to Proposition 10.13 of Functional Analysis I.
Proof: $\chi$ is a *-algebra homomorphism because $\varphi$ is an algebra homomorphism; Lemma 10.10 shows that $\chi$ respects *. $\chi$ is isometric (and thus injective), because

$$
\|\hat{x}\|_{\infty}^{2}=\left\|\chi\left(x^{*} x\right)\right\|_{\infty}=r\left(x^{*} x\right)=\left\|x^{*} x\right\|=\|x\|^{2}
$$

where you use Corollary 10.12 from Functional Analysis I for the second equality. $\chi$ is surjective because of the Stone-Weierstraß Theorem.

Remark 2.5: (i) Gelfand-Naimarks's Theorem (from the 1940's) is a milestone in the theory of $C^{*}$-algebras since:
(1) It justifies the view point of "non-commutative topology".
(2) It yields the very useful tool of "functional calculus".

Can we get (1) and (2) in the non-unital situation, too?
(ii) About "non-commutative topology" as a part of the "non-commutative world":


Figure 2.1: The noncommutiser.

Topology $\longrightarrow C^{*}$-algebras
Measure theory $\longrightarrow$ von Neumann algebras
Probability theory / Independce $\longrightarrow$ Free Probability / Free Independence

[^5]
## Differential Geometry $\longrightarrow$ Noncommutative Geometry

Symmetry / Groups $\longrightarrow$ Quantum symmetry / Quantum groups

Information theory $\longrightarrow$ Quantum information
Complex analysis $\longrightarrow$ Free Analysis
(iii) A dictionary of non-commutative topolgy:
topology ( $X$ locally compact) $\quad C^{*}$-algebra $\left(C_{0}(X)\right)$
compact unital
open subset closed ideal
closed subset quotient
metrisable separable
connected projectionless (no trivial projections)
See Gracia-Bondia, Vanilly, Figuenoa : Elements of non-commutative geometry, Introduction to chapter 1 or end of chapter 3; see also Wegge-Olsen, Chapter 1.11.
(iv) History / Importance of $C^{*}$-algebras:

- The name " $C^{*}$-algebras" was introduced by Segal in 1947 where " $C$ " stands for "continuous" and "*" stands for "involution".
- The study of $C^{*}$-algebras may be seen as the study of operators on a Hilbert space by algebraic means (refer to the Introduction of " $C^{*}$-algebras and their automorphism groups" by Pedersen). Reminder: Suppose that $A \subseteq B(H)$ is a subalgebra such that " $x \in A \Rightarrow x^{*} \in A$ " holds. Then $\mathrm{cl}_{\|\cdot\|}(A) \subseteq B(H)$ is a $C^{*}$-algebra and $\operatorname{cl}_{\text {SOT }}(A)$ and $\mathrm{cl}_{\text {WOT }}(A)$ are von Neumann algebras. One can show: There always is a faithful representation

$$
\pi: A \hookrightarrow B(H)
$$

Lemma 2.6: (i) Let $A$ and $B$ be $C^{*}$-algebras. Then $A \oplus B$ is again a $C^{*}$-algebra via

$$
A \oplus B:=\{(a, b) \mid a \in A, b \in B\}
$$

with pointwise operations and $\|(x, y)\|:=\max \{\|x\|,\|y\|\}$.
(ii) Let A be $a^{*}$-Banach algebra. Then

$$
\tilde{A}:=\{(x, \lambda) \mid x \in A, \lambda \in \mathbb{C}\}
$$

is $a^{*}$-Banach algebra with unit $(0,1)$ with the operations

$$
\begin{aligned}
(x, \lambda)+(y, \mu) & :=(x+y, \lambda+\mu) \\
(x, \lambda)^{*}: & =\left(x^{*}, \bar{\lambda}\right)
\end{aligned}
$$

$$
(x, \lambda)(y, \mu):=(x y+\lambda y+\mu x, \lambda \mu)
$$

and the norm

$$
\|(x, \lambda)\|_{B A}:=\|x\|+|\lambda| .
$$

Proof: For the Banach property of the new algebra defined in (ii):

$$
\|(x, \lambda)(y, \mu)\| \leq\|x\|\|y\|+|\lambda|\|y\|+|\mu|\|x\|+|\lambda||\mu|=\|(x, \lambda)\|\|(y, \mu)\| .
$$

Note that $\left(\tilde{A},\|\cdot\|_{B A}\right)$ is a unital *-Banach algebra, but not a $C^{*}$-algebra, since $\left\|(x, \lambda)^{*}(x, \lambda)\right\| \neq\|(x, \lambda)\|^{2}$. We write $\lambda 1+x:=(x, \lambda)$ for the elements in $\tilde{A}$ making the multiplication intuitive. The embedding

$$
\begin{aligned}
& A \hookrightarrow \tilde{A} \\
& x \longmapsto(x, 0)
\end{aligned}
$$

is an injective *-homomorphism.
Theorem 2.7: Let $A$ be a $C^{*}$-algebra. On $\tilde{A}$ there is a unique norm turning it into a $C^{*}$-algebra. Then

$$
\begin{aligned}
& A \longleftrightarrow \tilde{A} \\
& x \longmapsto(x, 0)
\end{aligned}
$$

is isometric.
Proof: The uniqueness is granted by Proposition 2.3, we thus only have to show exsistence.
(1) Let $A$ be unital with unit $e$. Then

$$
\begin{aligned}
\tilde{A} & \longrightarrow A \oplus \mathbb{C} \\
(x, \lambda) & \longmapsto(\lambda e+x, \lambda)
\end{aligned}
$$

is a bijective unital *-algebra homomorphism (so $\tilde{A} \cong A \oplus \mathbb{C}$ as ${ }^{*}$-algebras). Hence $\tilde{A}$ has a $C^{*}$-norm, namely the one of $A \oplus \mathbb{C}$.
(2) Let $A$ be non-unital. For $x \in \tilde{A}$, write $x=\lambda 1+a$. Consider

$$
L: \tilde{A} \longrightarrow B(A):=\{T: A \rightarrow A \text { bounded, linear }\}
$$

where $L_{x}(b):=x b=\lambda b+a b$ for $b \in A$. The linearity of $L$ is clear. For the boundedness we notice that

$$
\left\|L_{x}(b)\right\|=\|x b\|_{A}=\|\lambda b+a b\|_{A} \leq\left(|\lambda|+\|a\|_{A}\right)\|b\|_{A}
$$

and thus $\left\|L_{x}\right\| \leq|\lambda|+\|a\|=\|x\|_{B A}$. Now put

$$
\left.\|x\|_{\tilde{A}}:=\left\|L_{x}\right\|=\sup \{\| \lambda 1+a) b\left\|_{A} \mid b \in A,\right\| b \| \leq 1\right\}
$$

with $x=\lambda 1+a$. We have $\|a\|_{\tilde{A}}=\|a\|_{A}$ for all $a \in A$, as

$$
\|a\|\left\|a^{*}\right\|=\|a\|^{2}=\left\|a a^{*}\right\|=\left\|L_{a}\left(a^{*}\right)\right\| \leq\left\|L_{a}\right\|\left\|a^{*}\right\|
$$

and thus $\|a\|_{A} \leq\|a\|_{\tilde{A}}$. On the other hand we know that $\|a\|_{\tilde{A}}=\left\|L_{a}\right\| \leq\|a\|_{A}$ from the proof that $L_{a}$ is bounded.

To show that $\|\cdot\|_{\tilde{A}}$ is a norm, we only need to show that " $\|x\|_{A}=0 \Rightarrow x=0$ " holds. Let $x \in \tilde{A}$ such that $x=\lambda 1+a, \lambda \neq 0$. Assume that $\left\|L_{x}\right\|=0$, hence $x b=0$ for all $b \in A$ and thus $\lambda b+a b=0$ for all $b \in A$. Therefore, $-\lambda^{-1} a$ is a unit in $A$, since $b=-\lambda^{-1} a b$. This is a contradiction as $A$ was assumed to be non-unital.

Now we need to show that $\|\cdot\|_{\tilde{A}}$ is submultiplicative: It holds that

$$
\left\|L_{x y}\right\|=\left\|L_{x} L_{y}\right\| \leq\left\|L_{x}\right\|\left\|L_{y}\right\| .
$$

$\left(\tilde{A},\|\cdot\|_{\tilde{A}}\right)$ is in fact a *-Banach algebra, but we will omit the proof of completeness.
Finally, $\|\cdot\|_{\tilde{A}}$ satisfies the $C^{*}$-condition: Let $x \in \tilde{A}$ and $\varepsilon>0$. There is a $b \in A$ with $\|b\| \leq 1$ such that $\|x b\|_{A} \geq\left\|L_{x}\right\|-\varepsilon$, thus

$$
\left(\left\|L_{x}\right\|-\varepsilon\right)^{2} \leq\|x b\|_{A}^{2}=\left\|b^{*} x^{*} x b\right\|_{A} \leq\left\|b^{*}\right\|_{A}\left\|L_{x^{*} x}(b)\right\|_{A} \leq\left\|L_{x^{*} x}\right\|
$$

Therefore we get $\|x\|_{\tilde{A}}^{2}=\left\|L_{x}\right\|^{2} \leq\left\|L_{x^{*} x}\right\|=\left\|x^{*} x\right\|_{\tilde{A}}$. In general it holds: $\|x\|^{2} \leq$ $\left\|x^{*} x\right\| \Rightarrow\left\|x^{*} x\right\|$, because

$$
\begin{aligned}
\|x\|^{2} \leq\left\|x^{*} x\right\| \leq\left\|x^{*}\right\|\|x\| & \Rightarrow\|x\| \leq\left\|x^{*}\right\| \leq\left\|x^{* *}\right\|=\|x\| \\
& \Rightarrow\|x\|^{2} \leq\left\|x^{*} x\right\| \leq\|x\|^{2} .
\end{aligned}
$$

Remark 2.8: This is not the only unitisation of a $C^{*}$-algebra:
(i) Another possibiliy comes from the multiplier algebra

$$
\mathcal{M}(A):=\{(L, R) \text { double centralisers }\}
$$

with $L, R \in B(A)$ such that $L(a b)=L(a) b, R(a b)=a R(b)$ and $R(a) b=$ $a L(b)$. For instance let $L_{x}(a):=x a$ and $R_{x}(a):=a x$, then $\left(L_{x}, R_{x}\right) \in \mathcal{M}(A)$. We have $\|L\|=\|R\| . \mathcal{M}(A)$ is a $C^{*}$-algebra via

$$
\begin{aligned}
\left(L_{1}, R_{1}\right)+\left(L_{2}, R_{2}\right) & :=\left(L_{1}+L_{2}, R_{1}+R_{2}\right) \\
\lambda(L, R) & :=(\lambda L, \lambda R) \\
\left(L_{1}, R_{1}\right)\left(L_{2}, R_{2}\right) & :=\left(L_{1} L_{2}, R_{2} R_{1}\right) \\
(L, R)^{*} & :=\left(R^{*}, L^{*}\right) \\
L^{*}(a) & :=L\left(a^{*}\right)^{*}
\end{aligned}
$$

$$
\begin{aligned}
R^{*}(a) & :=R\left(a^{*}\right)^{*} \\
\|(L, R)\| & :=\|L\|=\|R\| \\
1 & :=(\mathrm{id}, \mathrm{id}) .
\end{aligned}
$$

Then $\mathcal{M}(A)$ is a unital $C^{*}$-algebra and

$$
\begin{aligned}
A & \longleftrightarrow \mathcal{M}(A) \\
x & \longmapsto\left(L_{x}, R_{x}\right)
\end{aligned}
$$

is an isometric embedding.
(ii) In fact $A \subseteq \tilde{A}$ and $A \subseteq \mathcal{M}(A)$ are ideals.
(iii) Let $A, B$ be $C^{*}$-algebras, $B$ unital, $A \subseteq B$ as an ideal. Then the diagrams

commute. In this sense, $\tilde{A}$ is the minimal and $\mathcal{M}(A)$ is the maximal unitisation.
(iv) Let $X$ be locally compact. Then

$$
\mathcal{M}\left(C_{0}(X)\right)=C_{b}(X):=\{f: X \rightarrow \mathbb{C} \text { continuous, bounded }\}
$$

and

$$
\widetilde{C_{0}(X)}=C_{0}(X) \oplus \mathbb{C} 1=C(\hat{X})
$$

where $\hat{X}$ is the one point compactification.
(v) Let $H$ be a Hilbert space with $\operatorname{dim} H=\infty$. Then $\mathcal{M}(K(H))=B(H)$ and

$$
\widetilde{K(H)}=K(H) \oplus \mathbb{C} 1 \subsetneq B(H)
$$

where $1 \in B(H)$.
Remark 2.9: Let $A, B$ be algebras, $\varphi: A \rightarrow B$ be an algebra homomorphism. Then

$$
\begin{aligned}
& \tilde{\varphi}: \tilde{A} \longrightarrow \tilde{B} \\
& \lambda 1+a \longmapsto \lambda 1+\varphi(a)
\end{aligned}
$$

is a unital algebra homomorphism with $\left.\tilde{\varphi}\right|_{A}=\varphi$. Hence Proposition 2.3 is also true for non-unital $C^{*}$-algebras: Let $\varphi: A \rightarrow B$ as in Proposition 2.3 (iii), extend it to $\tilde{\varphi}: \tilde{A} \rightarrow \tilde{B}$. Hence $\|\tilde{\varphi}(x)\| \leq\|x\|$ for all $x \in \tilde{A}$, thus $\|\varphi(x)\| \leq x \|$ for all $x \in A$.

Definition 2.10: Let $A$ be a (not necessarily unital) $C^{*}$-algebra, $x \in A$. Then we define the spectrum of $x$ as

$$
\operatorname{Sp}(x):= \begin{cases}\operatorname{Sp}_{A} x & \text { if } A \text { is unital } \\ \operatorname{Sp}_{\tilde{A}}(x) & \text { if } A \text { is non-unital. }\end{cases}
$$

Remark 2.11: (i) If $A$ is non-unital, then $0 \in \operatorname{Sp}(x)$ for all $x \in A$.
(ii) If $A$ is unital and $u \in A$ is unitary (i. e. $u^{*} u=u u^{*}=1$ ), then we have $\operatorname{Sp}(u) \subseteq \mathbb{S}^{1}:=\{\lambda \in \mathbb{C}| | \lambda \mid=1\} \subseteq \mathbb{C}$.
(iii) If $x \in A$ is selfadjoint (i. e. $x=x^{*}$ ), then $\operatorname{Sp}(x) \subseteq \mathbb{R}$ (even for $1 \notin A$ ).
(iv) If $B \subseteq A$ is a $C^{*}$-subalgebra and $x \in B$, then $\operatorname{Sp}_{B}(x)=\operatorname{Sp}_{A}(x)$.

Definition 2.12: Let $A$ be a commutative $C^{*}$-algebra (possibly $1 \notin A$ ). Then we define the spectrum of $A$ as $\operatorname{Spec}(A):=\operatorname{Spec}(\tilde{A}) \backslash\{\tilde{0}\}$ where

$$
\begin{gathered}
\tilde{0}: \tilde{A} \longrightarrow \mathbb{C} \\
\lambda+x \longmapsto \lambda
\end{gathered}
$$

(notice that $\left.\tilde{0}\right|_{A}=0$ ).
Theorem 2.13: Let $A$ be a commutative $C^{*}$-algebra. Then $\operatorname{Spec}(A)$ is locally compact and we have


Definition 2.14: Let $A$ be a $C^{*}$-algebra, $M \subseteq A$ a subset. Then

$$
C^{*}(M):=\text { smallest } C^{*} \text {-algebra of } A \text { containing } M
$$

is called the enveloping $C^{*}$-algebra of $M$.
Remark 2.15: (i) Let $A, B$ be $C^{*}$-algebras and $M \subseteq A$. Furthermore let $\varphi, \psi: C^{*}(M) \rightarrow B$ be *-homomorphisms. Then it holds: If $\left.\varphi\right|_{M}=\left.\psi\right|_{M}$, then $\varphi=\psi$ (because $\left\{x \in C^{*}(M) \mid \varphi(x)=\psi(x)\right\} \subseteq A$ is a $C^{*}$-algebra containing $M$ ).
(ii) $x \in A$ is normal if and only if $x^{*} x=x x^{*}$ : We know that $C^{*}(x)$ is commutative if and only if $x$ is normal, as

$$
C^{*}(x)=\operatorname{cl}\left(\left\{\text { non-commutative polynomials in } x, x^{*}\right\}\right)
$$

Proposition 2.16: Let $A$ be a $C^{*}$-algebra and $x \in A$ be normal.
(i) If $A$ is unital, then

$$
\begin{aligned}
& \operatorname{Spec}\left(C^{*}(x, 1)\right) \cong \\
& \varphi \longmapsto \operatorname{Sp}(x) \\
& \varphi(x)
\end{aligned}
$$

is a homeomorphism.
(ii) If $A$ is possibly non-unital, $C^{*}(x)$ non-unital, then

$$
\begin{aligned}
\operatorname{Spec}\left(C^{*}(x)\right) & \cong \\
\varphi & \longmapsto p(x) \backslash\{0\} \\
& \longmapsto(x)
\end{aligned}
$$

is a homeomorphism.
Proof: (i) As in Functional Analysis I: $\varphi \mapsto \varphi(x)$ is injective by Remark 2.15 and surjective, since $\hat{x}\left(\operatorname{Spec}\left(C^{*}(x, 1)\right)\right)=\operatorname{Sp}(x)$; it is continuous with respect to the pointwise topolgy.
(ii) We have


Corollary 2.17: We have the following functional calculus
(i) Let $A$ be a unital $C^{*}$-algebra and let $x \in A$ be normal. Then

$$
\begin{aligned}
\Psi: C(\operatorname{Sp}(x)) & \cong \\
f & \longmapsto C^{*}(x, 1) \\
& f(x)
\end{aligned}
$$

is an isometric *-isomorphism mapping $\operatorname{id}_{\operatorname{Sp}(x)} \mapsto x$. Note that $C^{*}(x, 1) \cong$ $C\left(\operatorname{Spec}\left(C^{*}(x, 1)\right)\right) \cong C(\operatorname{Sp}(x))$ by The theorem of Gelfand-Naimark (Reminder 2.4) and Remark 2.15. We have the diagramm

$$
\begin{aligned}
& \Phi^{-1}: C^{*} \longrightarrow C\left(\operatorname{Spec}\left(C^{*}(x, 1)\right)\right) \longrightarrow C(\operatorname{Sp}(x)) \\
& x \longmapsto \hat{x}, \hat{x}(\varphi)=\varphi(x) f \\
& \longmapsto f \circ \alpha
\end{aligned}
$$

where

$$
\begin{aligned}
\alpha: \operatorname{Sp}(x) & \longrightarrow \operatorname{Spec}\left(C^{*}(x, 1)\right) \\
\lambda & \longmapsto \varphi, \quad \varphi(x)=\lambda .
\end{aligned}
$$

2 Unitisation of $C^{*}$-algebras
(ii) Let $A$ be non-unital and $x \in A$ be normal. Then we can extend the functional calculus

$$
\Phi: C_{0}(\operatorname{Sp}(x))=\{f \in C(\operatorname{Sp}(x)) \mid f(0)=0\} \longrightarrow C^{*}(x) .
$$

Some properties:

- *-homomorphism, so $(f+g)(x)=f(x)+g(x),(f g)(x)=f(x) g(x), \bar{f}(x)=$ $f(x)^{*}$,
- $\operatorname{Sp}(f(x))=f(\operatorname{Sp}(x)), g$ continuous on $f(\operatorname{Sp}(x))$ thus $(g \circ f)(x)=g(f(x))$. (see Sheet 3, Exercise 1)

If $f$ is a polynomial, it is clear what $f(x)$ is. If $f$ is a continuous function, that is not a polynomial, think of the Stone-Weierstraß theorem.

Example: (i) Let $x \in A$ be selfadjoint (i. e. $x=x^{*}$ ). Then $\operatorname{Sp}(x) \subseteq \mathbb{R}$. Consider $f_{+}, f_{-}: \mathbb{R} \rightarrow[0, \infty)$. Then for $x_{+}:=f_{+}(x)$ and $x_{-}:=f_{-}(x)$ we have that $x=$ $x_{+}-x_{-}$and $x_{+}, x_{-}$are selfadjoint, too. Furthermore, we know that $f_{+}(\operatorname{Sp}(x))=$ $\operatorname{Sp}\left(x_{+}\right), f_{-}(\operatorname{Sp}(x))=\operatorname{Sp}\left(x_{-}\right) \subseteq[0, \infty)$ and $x_{+} x_{-}=x_{-} x_{+}=0$.
(ii) Let $x$ be again selfadjoint with $\operatorname{Sp}(x) \subseteq[0, \infty)$. Then there is $\sqrt{x} \in A$ that is selfadjoint with $\operatorname{Sp}(\sqrt{x}) \subseteq[0, \infty)$ and $\sqrt{x}^{2}=x$.

Proposition 2.18: Let $A, B$ be $C^{*}$-algebras and let $\varphi: A \rightarrow B$ be $a^{*}$-homomorphism, $x \in A$ be normal and $f \in C(\operatorname{Sp}(x))$ (or $f \in C_{0}(\operatorname{Sp}(x))$ ). Then

$$
\varphi(f(x))=f(\varphi(x))
$$

Proof: Since $\varphi$ is an algebra homomorphism, the statement is clear for polynomials. Then use Stone-Weierstraß approximation.

Remark 2.19: This is the continous functional calculus ( $f$ continuous). There is also a measurable functional calculus ( $f$ measurable, von Neumann algebras). Furthermore there is an analytic - or holomorphic functional calculus, which doesn't require $x$ to be normal.

## 3 Positive elements in $C^{*}$-algebras

Definition 3.1: Let $A$ be a $C^{*}$-algebra and $a \in A$. We say that $a$ is positive (we write $a \geq 0$ ), if $a=a^{*}$ and $\operatorname{Sp}(a) \subseteq[0, \infty)$. We write $a \leq b$, if $b-a \geq 0$.

Remark 3.2: Every positive element in a $C^{*}$-algebra possesses a unique positive square root $\sqrt{a} \in A$.

Lemma 3.3: Let $A$ be a unital $C^{*}$-algebra, $a \in A$ be selfadjoint and let $\lambda \geq\|a\|$. Then $a \geq 0$ if and only if $\|\lambda 1-a\| \leq \lambda$.

Proof: This statement is clear for functions (sketch missing).
For the fomal proof: $a \geq 0$ holds if and only if $\operatorname{Sp}(a) \subseteq[0, \infty)$ which holds if and only if $\mathrm{id}_{\operatorname{Sp}(a)} \geq 0$. $\mathrm{id}_{\operatorname{Sp}(a)} \geq 0$ holds if and only if

$$
\|\lambda 1-a\|=\left\|\Phi^{-1}(\lambda 1-a)\right\|_{\infty}=\left\|\lambda 1-\operatorname{id}_{\operatorname{Sp}(a)}\right\|_{\infty} \leq \lambda
$$

Proposition 3.4: If $a, b \geq 0$, then also $a+b \geq 0$.
Proof: Put $\lambda:=\|a\|+\|b\| \geq\|a+b\|$. Then

$$
\|\lambda-(a+b)\| \leq\| \| a\|1+a\|+\| \| b\|1-b\| \leq\|a\|+\|b\| \leq \lambda
$$

due to $a \geq 0$ and $b \geq 0$ and Lemma 3.3.
Lemma 3.5: Let $A$ be a $C^{*}$-algebra, $x \in A$ and $-x^{*} x \geq 0$. Then $x=0$.
Proof: Write $x=x_{1}+\mathrm{i} x_{2}$ where $x_{1}=\frac{1}{2}\left(x+x^{*}\right)$ and $x_{2}=\frac{1}{2 \mathrm{i}}\left(x-x^{*}\right)$, hence $x_{1}$ and $x_{2}$ are selfadjoint ("decomposition in real- and imaginary part"). We then have
$x^{*} x+x x^{*}=\left(x_{1}^{2}+\mathrm{i} x_{1} x_{2}-\mathrm{i} x_{2} x_{1}+x_{2}^{2}\right)+\left(x_{1}^{2}+\mathrm{i} x_{2} x_{1}-\mathrm{i} x_{1} x_{2}+x_{2}^{2}\right)=2 x_{1}^{2}+2 x_{2}^{2}$,
thus $x x^{*}=2 x_{1}^{2}+2 x_{2}^{2}-x^{*} x$. Now $2 x_{1}^{2}$ is positive, because $x_{1}=x_{1}^{*}$ and for $f(z):=z^{2}$ we have $\operatorname{Sp}\left(f\left(x_{1}\right)\right)=f\left(\operatorname{Sp}\left(x_{1}\right)\right) \subseteq[0, \infty)$; the same holds for $2 x_{2}^{2}$ and $-x x^{*}$ is positive by assumption. Via Proposition 3.4, $x x^{*}$ is positive. Now we have

$$
\operatorname{Sp}\left(x^{*} x\right) \cup\{0\}=\operatorname{Sp}\left(x x^{*}\right) \cup\{0\} \subseteq[0, \infty), \quad(\text { refer to Sheet } 3, \text { Exercise } 2),
$$

on the other hand $\operatorname{Sp}\left(x^{*} x\right) \subseteq(-\infty, 0]$, because $-x x^{*}$ is positive, thus $\operatorname{Sp}\left(x^{*} x\right)=\{0\}$ and via $\left\|x^{*} x\right\|=r\left(x^{*} x\right)=0$ we infer that $x=0$.

Proposition 3.6: Let $A$ be a $C^{*}$-algebra and $a \in A$. Then the following are equivalent:
(i) $a$ is positive,
(ii) There is $h \in A$ with $h=h^{*}$ such that $a=h^{2}$,
(iii) There is $x \in A$ such that $a=x^{*} x . \quad$ (algebraic way to express positivity)

Proof: "(i) $\Rightarrow$ (ii)": We define $h:=\sqrt{a}$ by Remark 3.2.
"(ii) $\Rightarrow$ (iii)": Put $x:=h$.
"(iii) $\Rightarrow$ (i)": Write $x^{*} x=u-v$, where $u:=\left(x^{*} x\right)_{+} \geq 0, v:=\left(x^{*} x\right)_{-} \geq 0$. For $u, v$ it holds $u v=v u=0$. Put $y:=x v$. Then

$$
-y^{*} y=-v x^{*} x v=-v(u-v) v=v^{3} \geq 0
$$

by $f(\operatorname{Sp}(v))=\operatorname{Sp}(f(v))$. Now from Lemma 3.5 we deduce that $y=0$, thus $v^{3}=0$ and consequently $v=0$ (If $\left\|v^{2}\right\|^{2}=\left\|\left(v^{2}\right)^{*}\left(v^{2}\right)\right\|=\left\|v^{4}\right\|=0$, then $v^{2}=0$ and thus $\|v\|^{2}=\left\|v^{*} v\right\|=\left\|v^{2}\right\|=0$ ). Now it holds $x^{*} x=u-v=u \geq 0$.

Corollary 3.7: In a $C^{*}$-algebra $A$, we put

$$
\begin{aligned}
& A_{+}:=\left\{h \in A \mid h=h^{*}, h \geq 0\right\}=\left\{x^{*} x \mid x \in A\right\} \subseteq A \quad \text { (Proposition 3.6) } \\
& A_{\mathrm{sa}}:=\left\{h \in A \mid h=h^{*}\right\} \subseteq A_{+}
\end{aligned}
$$

Then $A_{+}$is a cone, i.e.
(i) If $h+A_{+}$and $\lambda \geq 0$, then $\lambda h \in A_{+}$,
(ii) If $h_{1}, h_{2} \in A_{+}$, then $h_{1}+h_{2} \in A_{+}$.

Moreover $A_{+} \cap\left(-A_{+}\right)=\{0\}$ by Lemma 3.5 and Proposition 3.6, $A_{\mathrm{sa}}=A_{+}-A_{-}$ by (Chapter 2) and $A_{+}$is closed by Lemma 3.3.

We observe that due to $\left\|x^{*} x\right\|=\|x\|^{2}$, the positive elements play a special role in the theory of $C^{*}$-algebras.

Proposition 3.8: In a $C^{*}$-algebra, the following holds:
(i) If $a, b \in A_{\mathrm{sa}}, c \in A$ and $a \leq b$, then $c^{*} a c \leq c^{*} b c$.
(ii) If $0 \leq a \leq b$, then $\|a\| \leq\|b\|$.
(iii) If $A$ is unital and $0 \leq a \leq b$ are invertible elements of $A$, then $0 \leq b^{-1} \leq a^{-1}$,
(iv) For $\beta \in \mathbb{R}$ with $0 \leq b \leq 1$ and $0 \leq a \leq b$, it holds that $0 \leq a^{\beta} \leq b^{\beta}$. In particular we have $0 \leq \sqrt{a} \leq \sqrt{b}$.

Proof: (i) It holds $c^{*} b c-c^{*} a c=c^{*}(b-a) c=c^{*} x^{*} x c \geq 0$ for some $x \in A$.
(ii) Without loss of generality let $1 \in A$. We view $a$ and $b$ as functions $\operatorname{id}_{\operatorname{Sp}(a)}$ and $\operatorname{id}_{\operatorname{Sp}(b)}$, thus

$$
\|a\|=\inf \{\lambda \geq 0 \mid \lambda 1 \geq a\}, \quad\|b\|=\inf \{\lambda \geq 0 \mid \lambda 1 \geq b\}
$$

Now we infer $\|b\| 1 \geq b \geq a$, hence $\|b\| \geq\|a\|$.
(iii) Because $a \geq 0$, it also holds $a^{-1} \geq 0$ by the functional calculus. Futhermore the functional calculus gives that if $a \geq 1$, then $a^{-1} \leq 1$. Hence

$$
1=\sqrt{a^{-1}} a \sqrt{a^{-1}} \stackrel{\text { (i) }}{\leq} \sqrt{a^{-1}} b \sqrt{a^{-1}}
$$

and thus $d:=\sqrt{a} b^{-1} \sqrt{a}=\left(\sqrt{a^{-1}} b \sqrt{a^{-1}}\right)^{-1} \leq 1$. Again using (i), we get the estimate $b^{-1}=\sqrt{a^{-1}} d \sqrt{a^{-1}} \leq \sqrt{a^{-1}} 1 \sqrt{a^{-1}}=a^{-1}$.
(iv) This proof is complicated and therefore left out.

Remark 3.9: We do not have Proposition 3.8 (iv) for $\beta>1$, in particular $0 \leq a \leq b$ but $a^{2} \not \leq b^{2}$ may occur. We may even prove

$$
\left(\exists \beta>1 \forall x, y \in A:\left(0 \leq x \leq y \Rightarrow x^{\beta} \leq y^{\beta}\right)\right) \Rightarrow(a b=b a \forall a, b \in A)
$$

See also Sheet 3, Exercise 4 for other strange things.

## 4 Approximate units, ideals, quotients

Definition 4.1: Let $X$ be a topological space. A family $\left(x_{\lambda}\right)_{\lambda \in \Lambda} \subseteq X$ is a net, if $\Lambda$ is a filtration (i.e. we have an order " $\leq$ " on $\Lambda$, such that $\lambda \leq \lambda$ for all $\lambda \in \Lambda$, if $\lambda \leq \mu \leq \lambda$, then $\lambda=\mu$, if $\lambda_{1} \leq \lambda_{2} \leq \lambda_{3}$, then $\lambda_{1} \leq \lambda_{3}$ and for all $\lambda, \mu \in \Lambda$, there is $\nu \in \Lambda$ such that $\lambda \leq \nu$ and $\mu \leq \nu)$.
$\left(x_{\lambda}\right)_{\lambda \in \Lambda}$ converges to $x \in X$, if for any neighbourhood $U$ of $x$, there is a $\lambda_{0} \in \Lambda$ such that $x_{\lambda} \in U$ for all $\lambda \geq \lambda_{0}$.

Definition 4.2: Let $A$ be a $C^{*}$-algebra. An approximate unit $\left(u_{\lambda}\right)_{\lambda \in \Lambda}$ in $A$ is a net such that $\lambda \leq \mu \Rightarrow u_{\lambda} \leq u_{\mu}, 0 \leq u_{\lambda} \leq 1$ (i.e. $\left\|u_{\lambda}\right\| \leq 1$ ), and $u_{\lambda} x \rightarrow x$ and $x u_{\lambda} \rightarrow x$ for all $x \in A$.

Example 4.3: (i) Let $A=C_{0}(\mathbb{R})$. Then the net $\left(u_{N}\right)_{N \in \mathbb{N}}$, where $u_{N}$ is the function

is an approximate unit.
(ii) Let $A=K(H)$ for a separable Hilbert space. Then the net of projections $\left(u_{n}\right)_{n \in \mathbb{N}}$, where $u_{n}\left(\sum_{i \in \mathbb{N}} \alpha_{i} e_{i}\right):=\sum_{i=1}^{n} \alpha_{i} e_{i}$ and $\left(e_{i}\right)_{i \in \mathbb{N}}$ is an orthonormal basis, is an approximate unit.

Theorem 4.4: Let $A$ be a $C^{*}$-algebra (or an ideal in a $C^{*}$-algebra). Then $A$ has an approximate unit.

Proof: (i) The set $\Lambda:=\{h \in A \mid h \geq 0,\|h\|<1\}$ is a filtration.
Proof: The order shall be the order of $A_{+}, h \leq h$ for all $h \in \Lambda$ is clear. "If $h \leq g \leq h$, then $g=h$ " follows from $A_{+} \cap-A_{+}=\{0\}$. For $f, g, h \in \Lambda$ with $h \leq g \leq f, h \leq f$ follows from Proposition 3.4.

Let now $a, b \in \Lambda$. We need to find an element $c \in \Lambda$ such that $a \leq c$ and $b \leq c$. Put $a^{\prime}:=a(1-a)^{-1} \geq 0$ (note that $1 \notin \operatorname{Sp}(a)$, since $\|a\|<1$ ). $a^{\prime}$ is positive because of $t /(1-t) \geq 0$ on $[0,1)$ and the functional calculus.

Using Proposition 3.8 (iii) one can check that if $(1+y)^{-1} \leq(1+x)^{-1}$, then $(x(1+x))^{-1}=1-(1+x)^{-1} \leq 1-(1+y)^{-1}=y(1+y)^{-1}$. Hence

$$
a=a^{\prime}\left(1+a^{\prime}\right)^{-1} \leq\left(a^{\prime}+b^{\prime}\right)\left(1+a^{\prime}+b^{\prime}\right)^{-1}=: c .
$$

With what we have shown above we know that $a \leq c$ and $b \leq c$. Now $c \geq 0$, since $a^{\prime}+b^{\prime} \geq 0,1 /(1+t) \geq 0$ and $\|c\|<1$ since $1 /(1+t)<1$, thus $c \in \Lambda$.
(ii) Let $h \geq 0, h \in A$ and $n \in \mathbb{N}$. Then $h\left(\frac{1}{n}+h\right)^{-1} \in A, h\left(1-h\left(\frac{1}{n}-h\right)^{-1}\right) \leq \frac{1}{n}$.

Proof: Because of $0 \leq t /\left(t+\frac{1}{n}\right) \leq 1$ for all $t \geq 0$ and $t\left(1-t /\left(\frac{1}{n}-t\right)\right) \leq \frac{1}{n}$ for all $t \geq 0$, the assertion follows from the functional calculus.
(iii) Let $h \geq 0, g \in \Lambda, h\left(\frac{1}{n}+h\right)^{-1} \leq g$. Then $\|h-g h\|^{2} \leq \frac{1}{n}\|h\|$ as well as $\|h-h g\|^{2} \leq \frac{1}{n}\|h\|$.

Proof: We have

$$
\|h-g h\|^{2}=\left\|h(1-g)^{2} h\right\| \leq\|h(1-g) h\| \leq\left\|h\left(1-h\left(\frac{1}{n}-h\right)^{-1}\right) h\right\| \leq \frac{1}{n}\|h\|
$$

for the second inequality we use

$$
(1-g)-(1-g)^{2}=g(1-g) \geq 0 \quad(\text { for } 0 \leq g \leq 1)
$$

hence $(1-g)^{2} \leq(1-g)$, then use Proposition 3.8 (i), (ii).
Finally, let $x \in A$ and $\varepsilon>0$. Put $h:=x x^{*} \geq 0$. By (ii), we know that $\lambda_{0}:=h\left(\frac{1}{n}-h\right)^{-1} \in \Lambda$ for $n \in \mathbb{N}$ with $\frac{1}{n}\|h\|<\varepsilon^{2}$. Hence for all $g \in \Lambda$ with $\lambda_{0} \leq g$ :

$$
\|x-g x\|^{2}=\|(1-g) h(1-g)\| \leq\|h-g h\|(1+\|g\|),
$$

if we put $u_{g}:=g$, we conclude $x-g x \rightarrow 0$ for $g \rightarrow \infty$. Likewise $x-x g \rightarrow 0$.
Remark 4.5: If $A$ is a separable $C^{*}$-algebra (i. e. it has a dense countable subset), then there is an approximate unit $\left(u_{n}\right)_{n \in \mathbb{N}}, u_{1} \leq u_{2} \leq \ldots$ (i.e. $A$ is $\sigma$-unital, i.e. $\Lambda$ is countable).

Definition 4.6: Let $A$ be $C^{*}$-algebra. $I$ is an ideal in $A$, if $I \subseteq A$ is a closed linear subspace such that $A I, I A \subseteq I$. We write $I \triangleleft A$.

Example 4.7: (i) Let $H$ be a Hilbert space. Then $K(H) \triangleleft B(H)$.
(ii) Let $A$ be a $C^{*}$-algebra. Then $A \triangleleft \tilde{A}$ using the notation from Chapter 2.

Lemma 4.8: Let $A$ be a $C^{*}$-algebra.
(i) If $I \triangleleft A$, then $I=I^{*}$ (hence $I$ is a $C^{*}$-algebra),
(ii) If $I \triangleleft J \triangleleft A$, then $I \triangleleft A$,
(iii) If $I \triangleleft A, I \subsetneq A$ and $\left(u_{\lambda}\right)_{\lambda \in \Lambda}$ is an approximate unit of $I$, then $\left(u_{\lambda}\right)_{\lambda \in \Lambda}$ is not an approximate unit for $A$

Proof: (i) Let $\left(u_{\lambda}\right)_{\lambda \in \Lambda}$ be an approximate unit for $I$ and $x \in I$. Then $u_{\lambda} x \rightarrow x$ and because ${ }^{*}$ is continuous and $I$ is closed, then $I \ni x^{*} u_{\lambda}=\left(u_{\lambda} x\right)^{*} \rightarrow x$.
(ii) Let $\left(u_{\lambda}\right)_{\lambda \in \Lambda}$ be an approximate unit for $I, x \in I$ and $a \in A$. Then $I \ni$ $x u_{\lambda} a \rightarrow x a$, thus $x a \in I, a x \in I$.
(iii) Assume $\left(u_{\lambda}\right)_{\lambda \in \Lambda}$ was also an approximate unit for $A$. Then for all $a \in A$ : $I \ni u_{\lambda} a \rightarrow a$, hence $I=A$ which is a contradiction.

Theorem 4.9: Let $A$ be a $C^{*}$-algebra and $I \triangleleft A$. Then $A / I$ is a $C^{*}$-algebra.
Proof: $A / I$ is a Banach algebra with the norm

$$
\|\dot{x}\|:=\inf \{\|x+z\| \mid z \in I\}
$$

for $\dot{x} \in A / I$. It has an involution $(\dot{x})^{*}:=\left(x^{*}\right)^{\bullet}$ (check that we have " $\dot{x}=\dot{y} \Rightarrow$ $\left.x-y \in I \Rightarrow(x-y)^{*} \in I \Rightarrow\left(x^{*}\right)^{\bullet}=\left(y^{*}\right)^{\bullet}\right)$. It remains to be shown that $\left\|\dot{x}^{*} \dot{x}\right\| \geq\|\dot{x}\|^{2}$ (due to Proposition 2.3, " $\leq$ " then follows).

Let $\left(u_{\lambda}\right)_{\lambda \in \Lambda}$ be an approximate unit for $I$ and let $x \in A$. Then $\|\dot{x}\|=\lim \| x-$ $u_{\lambda} x \|$.

Proof: Let $\varepsilon>0$. We find $z \in I$ such that $\|x+z\| \leq\|\dot{x}\|+\varepsilon$ and we find $\lambda_{0} \in \Lambda$ such that $\left\|z-u_{\lambda} z\right\|<\varepsilon$ for all $\lambda \geq \lambda_{0}$. Then we have

$$
\begin{aligned}
\|\dot{x}\| \leq\left\|x-u_{\lambda} x\right\| & \leq\left\|\left(1-u_{\lambda}\right)(x+z)\right\|+\left\|\left(1-u_{\lambda}\right) z\right\| \\
& \leq\left\|1-u_{\lambda}\right\|\|x+z\|+\left\|z-u_{\lambda} z\right\| \leq\|\dot{x}\|+2 \varepsilon
\end{aligned}
$$

as $\|x+z\| \leq\|\dot{x}\|+\varepsilon$ and $\left\|z-u_{\lambda} z\right\|<\varepsilon$ by assumption.
Now let $x \in A$. Then

$$
\begin{aligned}
\|\dot{x}\|^{2} & =\lim \left\|x-u_{\lambda} x\right\| \\
& =\lim \left\|\left(1-u_{\lambda}\right) x^{*} x\left(1-u_{\lambda}\right)\right\| \\
& =\lim \left\|\left(1-u_{\lambda}\right)\left(x^{*} x+z\right)\left(1-u_{\lambda}\right)\right\| \quad\left(\text { as } z \in I \text { and } u_{\lambda} z-z \rightarrow 0\right) \\
& \leq\left\|x^{*} x+z\right\| .
\end{aligned}
$$

By taking the Infimum over $z \in I$, we thus get $\mid \dot{x}\left\|^{2} \leq\right\| \dot{x}^{*} \dot{x} \|$.
Proposition 4.10: Let $A, B$ be $C^{*}$-algebras and let $\varphi: A \rightarrow B$ be $a^{*}$-homomorphism.
(i) If $\varphi$ is injective, then $\varphi$ is isometric.
(ii) $\varphi(A)$ is a $C^{*}$-algebra and $A / \operatorname{ker}(\varphi) \cong \varphi(A)$.

Proof: (i) We need to show that $\left\|\varphi\left(x^{*} x\right)\right\|=\left\|x^{*} x\right\|$, then $\|\varphi(x)\|^{2}=\left\|\varphi\left(x^{*} x\right)\right\|=$ $\left\|x^{*} x\right\|=\|x\|^{2}$. Let's assume that $\varphi$ is isometric, i. e. $\left\|\varphi\left(x^{*} x\right)\right\| \leq\left\|x^{*} x\right\|$ for some $x \in A$. Consider $f: \mathbb{R} \rightarrow \mathbb{R}$ with $0 \leq f \leq 1,\left.f\right|_{\left(\left\|x^{*} x\right\|, \infty\right)}=1$ and $\left.f\right|_{\left[0,\left\|\varphi\left(x^{*} x\right)\right\|\right]} \equiv 0$ (refer to Figure 4.1). Then

$$
\varphi\left(f\left(x^{*} x\right)\right)=f\left(\varphi\left(x^{*} x\right)\right)=0
$$

where $0 \neq\left\|f\left(x^{*} x\right)\right\|=\left\|\left.f\right|_{\left[0,\left\|x^{*} x\right\|\right]}\right\|_{\infty}=1$, since $r\left(x^{*} x\right)=\left\|x^{*} x\right\|$, hence $\left\|x^{*} x\right\| \in$ $\operatorname{Sp}\left(x^{*} x\right) \subseteq[0, \infty)$ (in principle: $\left\|x^{*} x\right\| \in \operatorname{Sp}\left(x^{*} x\right)$ or $-\left\|x^{*} x\right\| \in \operatorname{Sp}\left(x^{*} x\right)$ ) and $\left\|f\left(\varphi\left(x^{*} x\right)\right)\right\|=\left\|\left.f\right|_{\left[0,\left\|\varphi\left(x^{*} x\right)\right\|\right]}\right\|_{\infty}=0$, thus $\varphi$ is not injective.


Figure 4.1: Sketch of the function $f$ as mentioned in part (i) of the proof.
(ii) We have the diagramm

where $\dot{\varphi}(\dot{x}):=\varphi(x)$ is a well-defined ${ }^{*}$-homomorphism. Note that $\operatorname{ker}(\varphi) \triangleleft A$. By part (i), $\dot{\varphi}$ is isometric, therefore $\varphi(A)=\dot{\varphi}(A / \operatorname{ker}(\varphi))$ is complete and thus a $C^{*}$-algebra.

Remark 4.11 (Homological properties of $C^{*}$-algebras): The sequence

$$
0 \longrightarrow I \longrightarrow A \longrightarrow B \longrightarrow 0
$$

is exact if and only if $I \triangleleft A$ and $B \cong A / I$.

## 5 Positive linear functionals and the GNS construction

Definition 5.1: Let $A$ be a $C^{*}$-algebra and $\varphi: A \rightarrow \mathbb{C}$ be a linear functional. $\varphi$ is called positive (in signs $\varphi \geq 0$ ), if $\varphi(x) \geq 0$ for all $x \geq 0$.

A positive linear functional preserves the order, i. e. " $x \leq y \Rightarrow \varphi(x) \leq \varphi(y)$ " holds.

Example 5.2: (i) Let $A=C([0,1])$. Then $\varphi_{t}(f):=f(t)$ or $\varphi(f):=\int_{0}^{1} f(t) d t$ are positive functionals. More generally, we have the correspondence:

$$
\{\text { positive functionals on } A\} \xrightarrow{\sim}\{\text { Radon measures on }[0,1]\}
$$

$$
\varphi \longmapsto \mu \text { with } \varphi(f)=\int_{0}^{1} f d \mu
$$

(ii) Let $A=M_{n}(\mathbb{C})$. Then

$$
\begin{aligned}
\operatorname{tr}: A & \longrightarrow \mathbb{C} \\
\left(a_{i, j}\right)_{1 \leq i, j \leq n} & \longmapsto \sum_{i=1}^{n} a_{i, i}
\end{aligned}
$$

is a positive linear functional.
(iii) Let $A=B(H)$ and $\xi \in H$. Then $\varphi_{\xi}(x):=\langle x \xi, \xi\rangle$ is a positive functional.

Lemma 5.3: Let $A$ be a $C^{*}$-algebra and $\varphi: A \rightarrow \mathbb{C}$ be positive. Then $\langle x, y\rangle:=$ $\varphi\left(y^{*} x\right)$ is a positive sesquilinear form.

Proposition 5.4: Let $A$ be a $C^{*}$-algebra, $\varphi: A \rightarrow \mathbb{C}$ be positive.
(i) $\varphi$ is bounded (hence continuous),
(ii) $\varphi$ is involutive, i.e. $\varphi\left(x^{*}\right)=\overline{\varphi(x)}$ and $|\varphi(x)|^{2} \leq\|\varphi\| \varphi\left(x^{*} x\right)$ for all $x \in A$.

Proof: (i) (1) $\varphi$ is bounded on $S:=\{x \in A \mid x \geq 0,\|x\| \leq 1\}$.
Proof: Assume $\varphi$ was not bounded on $S$. Then there was a sequence $\left(a_{n}\right)_{n \in \mathbb{N}} \subseteq S$ such that $\varphi\left(a_{n}\right) \geq 2^{n}$. Then $a:=\sum_{n \in \mathbb{N}} \frac{1}{2^{n}} a_{n} \geq 0$ by (Corollary 3.7), but

$$
\varphi(a) \geq \varphi\left(\sum_{n=1}^{N} \frac{1}{2^{n}} a_{n}\right) \geq N
$$

which was a contradiction.
(2) For an arbitrary $z \in A$ write $z=(\operatorname{Re}(z))_{+}-\left(\operatorname{Re}(z)_{-}+\mathrm{i}\left((\operatorname{Im}(z))_{+}-(\operatorname{Im}(z))_{-}\right)\right.$ as a linear combination of four positive elements with norm smaller or equal to the norm of $z$. Thus $\|\varphi(z)\| \leq 4 K\|z\|$, where $K$ is the bound from (1).
(ii) By (Lemma 5.3), $\langle\cdot, \cdot\rangle$ is a positive sesquilinear form, hence the polarisation identity is satisfied (Theorem 5.7, Funtional Analaysis I). Hence $\langle x, y\rangle=\overline{\langle y, x\rangle}$. Thus

$$
\varphi\left(x^{*}\right) \leftarrow \varphi\left(x^{*} u_{\lambda}\right)=\left\langle u_{\lambda}, x\right\rangle=\overline{\left\langle x, u_{\lambda}\right\rangle}=\overline{\varphi\left(u_{\lambda}^{*} x\right)}=\overline{\varphi(u, x)} \rightarrow \overline{\varphi(x)},
$$

and moreover

$$
|\varphi(x)|^{2} \leftarrow\left|\varphi\left(u_{\lambda} x\right)\right|^{2}=\left|\left\langle x, u_{\lambda}\right\rangle\right|^{2} \leq\langle x, x\rangle\left\langle u_{\lambda}, u_{\lambda}\right\rangle=\varphi\left(x^{*} x\right) \varphi\left(u_{\lambda}^{2}\right)
$$

Proposition 5.5: Let $A$ be a $C^{*}$-algebra and $\varphi: A \rightarrow \mathbb{C}$ be continuous and linear. Then the following are equivalent:
(i) $\varphi$ is positive,
(ii) For all approximate units $\left(u_{\lambda}\right)_{\lambda \in \Lambda} \subseteq A$ it holds $\|\varphi\|=\lim \varphi\left(u_{\lambda}\right)$,
(iii) There is an approximate unit $\left(u_{\lambda}\right)_{\lambda \in \Lambda} \subseteq A$ such that $\|\varphi\|=\lim \varphi\left(u_{\lambda}\right)$.

Proof: "(i) $\Rightarrow$ (ii)": Without loss of generality, let $\|\varphi\|=1$. Then $\left(\varphi\left(u_{\lambda}\right)\right)_{\lambda \in \Lambda} \subseteq \mathbb{C}$ is a bounded, monontonically increasing net in $\mathbb{R}_{+}$, hence $\varphi\left(u_{\lambda}\right) \uparrow \alpha \leq 1$ for some $\alpha \in \mathbb{R}$. For $x \in A,\|x\| \leq 1$ we have

$$
|\varphi(x)|^{2} \leftarrow\left|\varphi\left(u_{\lambda} x\right)\right|^{2} \leq \varphi\left(u_{\lambda}^{2}\right) \varphi\left(x^{*} x\right) \leq \varphi\left(u_{\lambda}\right) \varphi\left(x^{*} x\right) \leq \alpha
$$

because of $0 \leq u_{\lambda} \leq u_{\lambda}^{2}$ and the functional calculus. Since $\|\varphi\|=1$ we find $x_{n} \in A$, $\left\|x_{n}\right\| \leq 1$ such that $\alpha \geq\left|\varphi\left(x_{n}\right)\right|^{2} \rightarrow 1$, thus $\alpha=1=\|\varphi\|$. "(ii) $\Rightarrow$ (iii)" is obviously true.
"(iii) $\Rightarrow$ (i)": Without loss of generality, let $\|\varphi\|=1$, thus $\varphi\left(u_{\lambda}\right) \rightarrow 1$.
(1) For $x \in A$ selfadjoint, $\|x\| \leq 1$, we have $\varphi(x) \in \mathbb{R}$.

Proof: Let $\varphi(x)=\alpha+\mathrm{i} \beta$, without loss of generality $\beta \leq 0$. Assume $\beta<0$. Then for all $n \in \mathbb{N}$ we have

$$
\begin{aligned}
\left\|x-\mathrm{i} n u_{\lambda}\right\|^{2}=\left\|\left(x+\mathrm{i} n u_{\lambda}\right)\left(x-\mathrm{i} n u_{\lambda}\right)\right\| & =\left\|x^{2}+n^{2} u_{\lambda}^{2}-\mathrm{i} n\left(x u_{\lambda} u_{\lambda} x\right)\right\| \\
& \leq 1+n^{2}+n \| x u_{\lambda}-u_{\lambda} x
\end{aligned}
$$

Therefore, $|\beta|^{2}+2|\beta| n+n^{2}=|\operatorname{Im}(\varphi(x)-\mathrm{i} n)|^{2} \leq|\varphi(x)-\mathrm{i} n|^{2} \leftarrow\left|\varphi\left(x-\mathrm{i} n u_{\lambda}\right)\right|$. Using the result of the above calulcation, we obtain

$$
|\beta|^{2}+2|\beta| n+n^{2} \leq 1+n^{2}+n\left\|x u_{\lambda}-u_{\lambda} x\right\|
$$

thus $2|\beta| n \leq 1-|\beta|^{2}$, hence $\beta=0$, which is a contradiction.
(2) Let now $x \geq 0,\|x\| \leq 1$. Then $-1 \leq u_{\lambda}-x \leq 1$ (where the last inequality holds because we have " $u_{\lambda} \leq 1 \Rightarrow 1-u_{\lambda} \geq 0 \Rightarrow 1-u_{\lambda}+x \geq 0$ "), hence $\left\|u_{\lambda}-x\right\| \leq 1$. By (i) we thus have

$$
1-\varphi(x) \leq|1-\varphi(x)| \leftarrow\left|\varphi\left(u_{\lambda}-1\right)\right| \leq 1
$$

which implies $\varphi(x) \geq 0$.

Corollary 5.6: Let $A$ be a unital $C^{*}$-algebra, $\varphi: A \rightarrow \mathbb{C}$ be continuous and linear. $\varphi$ is positive if and only if $\varphi(1)=\|\varphi\|$.

Corollary 5.7: Let $\varphi, \varphi^{\prime}$ be two positive functional. Then $\left\|\varphi+\varphi^{\prime}\right\|=\|\varphi\|+\left\|\varphi^{\prime}\right\|$.
Proof: We have

$$
\|\varphi\|+\left\|\varphi^{\prime}\right\| \leftarrow \varphi\left(u_{\lambda}\right)+\varphi^{\prime}\left(u_{\lambda}\right)=\left(\varphi+\varphi^{\prime}\right)\left(u_{\lambda}\right) \rightarrow\left\|\varphi+\varphi^{\prime}\right\| .
$$

Definition 5.8: A state on a $C^{*}$-algebra is a positive linear functional $\varphi$ with $\|\varphi\|=1$.

Remark 5.9: Let $A$ be a unital $C^{*}$-algebra. $\varphi$ is a state on $A$ if and only if $\varphi$ is positive and $\varphi(1)=1$.

Reminder 5.10 (Theorem of Hahn-Banach ${ }^{1}$ ): Let $E$ be a normed $\mathbb{C}$-vector space, $F \subseteq E$ be a linear subspace and $f: F \rightarrow \mathbb{C}$ continuous and linear. Then there is continuous and linear $\tilde{f}: E \rightarrow \mathbb{C}$ such that $\|\tilde{f}\|=\|f\|$ and $\left.\tilde{f}\right|_{F}=f$.

Theorem 5.11: Let $A$ be a $C^{*}$-algebra, $x \in A$ be normal. Then we find a state $\varphi: A \rightarrow \mathbb{C}$ such that $|\varphi(x)|=\|x\|$.

Proof: By the Gelfand isomorphism we find a character $\varphi_{0}: C^{*}(x, 1) \rightarrow \mathbb{C}$ such that $\varphi_{0}(1)=1$ with $\left|\varphi_{0}(x)\right|=\left|\hat{x}\left(\varphi_{0}\right)\right|=\|\hat{x}\|_{\infty}=\|x\|$. In particular, $\varphi_{0}$ is linear and continuous $\left(\left\|\varphi_{0}\right\|=1\right)$. By the Hahn-Banach theorem (Reminder 5.10), we find an extension $\tilde{\varphi}_{0}: \tilde{A} \rightarrow \mathbb{C}$ that is linear and continuous and fulfills that $\left\|\varphi_{0}\right\|=\left\|\tilde{\varphi}_{0}\right\|=1=\varphi_{0}(1)=\tilde{\varphi}_{0}(1)$.

By Corollary 5.6 we know that $\tilde{\varphi}_{0} \geq 0$ and thus $\varphi:=\left.\tilde{\varphi}_{0}\right|_{A}$ is positive and it holds that $|\varphi(x)|=\|x\|,\|\varphi\|=1$.

Definition 5.12: Let $H, H_{1}, H_{2},\left(H_{i}\right)_{i \in I}$ be Hilbert spaces and $A$ be a $C^{*}$-algebra.
(i) A representation of $A$ on $H$ is a *-homomorphism $\pi: A \rightarrow B(H)$.
(ii) Two representations $\left(\pi_{1}, H_{1}\right),\left(\pi_{2}, H_{2}\right)$ are equivalent, if there is a unitary map $U: H_{1} \rightarrow H_{2}$ such that $\pi_{2}(x)=U_{\pi_{1}}(x) U^{*}$ for all $x \in A$.
(iii) Let $\left(\pi_{i}, H_{i}\right)_{i \in I}$ be representations of $A$. Then $\left(\bigoplus_{i \in I} \pi_{i}, \bigoplus_{i \in I} H_{i}\right)$ is a representation on $\bigoplus_{i \in I} H_{i}$ given by

$$
\left(\bigoplus_{i \in I} \pi_{i}\right)(x) \xi_{j}:=\pi_{j}(x) \xi_{j} \quad j \in J
$$

(iv) A representation $\pi$ is non-degenerate if $\operatorname{cl}(\pi(A) H)=H$.
(v) A representation $\pi$ is cyclic if there is $\xi \in H$ (the so called cyclic vector) such that $\operatorname{cl}(\pi(A) \xi)=H$.

[^6]Every cyclic representation is in particular non-degenerate.
Remark 5.13: (i) Every representation is a direct sum of a non-degenerate and a zero representation.
(ii) Every non-degenerate representation is a direct sum of cyclic representations.

Lemma 5.14: Let $\left(\pi_{1}, H_{1}, \xi_{1}\right)$ and $\left(\pi_{2}, H_{2}, \xi_{2}\right)$ be two cyclic representations and $f_{1}, f_{2}: A \rightarrow \mathbb{C}$ be positive linear functionals with $f_{i}(x)=\left\langle\pi_{i}(x) \xi_{i}, \xi_{i}\right\rangle$ for $i=1,2$. If $f_{1}=f_{2}$ holds, then there is a unitary $U: H_{1} \rightarrow H_{2}$ with $U \xi_{1}=\xi_{2}$ such that $\pi_{2}(x)=U \pi_{1}(x) U^{*}$.

Proof: See Sheet 4, Exercise 3.
Theorem 5.15 (GNS-Construction): Let $f$ be a state on a $C^{*}$-algebra A. Then there is a unique (up to unitary equivalence) cyclic representation $\left(\pi_{f}, H_{f}, \xi_{f}\right)$ such that $f(x)=\left\langle\pi_{f}(x) \xi_{f}, \xi_{f}\right\rangle$.

Proof: (1) Firstly, given the data $A$ and $f$, we want to construct the Hilbert space for the representation of $A$ :

- $\langle x, y\rangle_{f}:=f\left(y^{*} x\right)$ is a positive sesquilinear form on $A \times A$
- Put $\mathcal{N}_{f}:=\left\{x \in A \mid\langle x, x\rangle_{f}=0\right\}$ and define $K_{f}:=A / \mathcal{N}_{f}$.

Then $K_{f}$ is a pre-Hilbert space: For the quotient map $\gamma: A \rightarrow A / \mathcal{N}_{f}$, the expression $\langle\gamma(x), \gamma(y)\rangle:=\langle x, y\rangle_{f}$ is well-defined and $\gamma$ continuous because

$$
\|\gamma(x)\|^{2}=\langle\gamma(x), \gamma(x)\rangle=f\left(x^{*} x\right) \leq\|x\|^{2}
$$

- Put $H_{f}:=\mathrm{cl}_{\|\cdot\|_{f}}\left(K_{f}\right)$. $H_{f}$ is then a Hilbert space with the inner product $\langle\gamma(x), \gamma(y)\rangle=f\left(y^{*} x\right)$.
(2) Secondly, we want to construct the representation of $A$ on $H_{f}$ :
- Define

$$
\begin{aligned}
\pi_{f}^{0}(x): K_{f} & \longrightarrow H_{f} \\
\gamma(y) & \longmapsto \gamma(x, y) .
\end{aligned}
$$

Then $\pi_{f}^{0}$ is continuous:

$$
\left\|\pi_{f}^{0}(x) \gamma(y)\right\|^{2}=\|\gamma(x y)\|^{2}=f\left(y^{*} x^{*} x y\right) \leq\left\|x^{*} x\right\| f\left(y^{*} y\right)=\|x\|^{2} \|\left.\gamma(y)\right|^{2}
$$

and thus $\left\|\pi_{f}^{0}(x)\right\| \leq\|x\|$. Also, this proves thast $\pi_{f}^{0}(x)$ is well-defined (If $\gamma(y)=\gamma(z)$, then $\left.\|\gamma(x y)-\gamma(x z)\|^{2}=\|\gamma(x(y-z))\|^{2} \leq\|x\|^{2}\|\gamma(y-z)\|^{2}=0\right)$.

- Extend $\pi_{f}^{0}$ to $\pi_{f}(x): H_{f} \rightarrow H_{f}$, then it holds $\left\|\pi_{f}(x)\right\| \leq\|x\|$ and due to $\pi_{f}^{0}(x) \pi_{f}^{0}(y)=\pi_{f}^{0}(x y)$, we also have $\pi_{f}(x) \pi_{f}(y)=\pi_{f}(x y)$.
- $\pi_{f}^{0}$ is a ${ }^{*}$-homomorphism, because

$$
\begin{aligned}
\pi_{f}^{0}\left(x^{*}\right)=\left\langle\pi_{f}(x) \gamma(y), \gamma(z)\right\rangle & =f\left(z^{*} x y\right)=f\left(\left(x^{*} z\right)^{*} y\right) \\
& =\left\langle\gamma(y), \gamma\left(x^{*} z\right)\right\rangle=\pi_{f}^{0}\left(x^{*}\right) \gamma(z)=\pi_{f}^{0}(x)^{*}
\end{aligned}
$$

This property transfers to the limit, thus $\pi_{f}$ is a ${ }^{*}$-homomorphism.
(3) Thirdly, we want to construct the cyclic vector $\xi_{f}$ :

- Let $\left(u_{\lambda}\right)_{\lambda \in \Lambda}$ be an approximate unit for $A$. Put $\xi_{f}:=\lim \gamma\left(u_{\lambda}\right)$. The limit $\lim \gamma\left(u_{\lambda}\right)$ exists: Let $\lambda \geq \mu$. Then

$$
\left.\left\|\gamma\left(u_{\lambda}\right)-\gamma\left(u_{\lambda}\right)\right\|^{2}=f\left(\left(u_{\lambda}-u_{\mu}\right)\right)^{2}\right) \leq f\left(u_{\lambda}-u_{\mu}\right)<\varepsilon
$$

for $\lambda, \mu$ large, as we have $u_{\lambda} \geq u_{\mu}, 1 \leq u_{\lambda} \geq u_{\lambda}-u_{\mu} \geq 0$ and via the functional calculus, we get $\left(u_{\lambda}-u_{\mu}\right)^{2} \leq\left(u_{\lambda}-u_{\mu}\right)^{2}$ and thus $\left(\gamma\left(u_{\lambda}\right)\right)_{\lambda \in \Lambda}$ is a Cauchy-net.

- $\xi_{f}$ is cyclic, as for $x \in A$ we have $\pi_{f}(x) \gamma\left(u_{\lambda}\right)=\gamma\left(x u_{\lambda}\right) \rightarrow \gamma(x)$, hence $\pi_{f}(x) \xi_{f}=\gamma(x)$ for all $\gamma(x) \in K_{f}$.
(4) Fourthly, we have

$$
\left\langle\pi_{f}(x) \xi_{f}, \xi_{f}\right\rangle=\lim \left\langle\pi_{f}(x) \gamma\left(u_{\lambda}\right), \gamma\left(u_{\lambda}\right)\right\rangle=\lim f\left(u_{\lambda} x u_{\lambda}\right)=f(x)
$$

(5) Finally, the uniqueness follows from Lemma 5.14.

Corollary 5.16 (Second Fundamental Theorem of $C^{*}$-algebras): Every $C^{*}$-algebra $A$ admits a faithful (i.e. injective) representation $\pi: A \hookrightarrow B(H)$. Hence, $A$ is isomorphic to a $C^{*}$-subalgebra of $B(H)$.

Of the above Corollary we can make "Every (abstractly defined) $C^{*}$-algebra has a concrete representation." or "The abstract (Chapter 2) and the concrete definition (as $\|\cdot\|$-closed ${ }^{*}$-subalgebra of $B(H)$ for a Hilbert space $H$ ) coincide."

Proof: Put $\pi:=\bigoplus_{f \text { state on } A} \pi_{f}$. We now need to show, that $\pi$ is faithful; let therefore $x \in A \backslash\{0\}$. By (Theorem 5.11) we find a state $f$ such that $f\left(x^{*} x\right)=\|x\|^{2}$. By Theorem 5.15, we get that $\left\|\pi_{f}(x) \xi_{f}\right\|^{2}=f\left(x^{*} x\right)=\|x\|^{2} \neq 0$, so $\pi_{f}(x) \neq 0$ and thus $\pi(x) \neq 0$ and $\pi$ is faithful.

Remark 5.17: If $A$ is separable, then $H$ in Corollary 5.16 may be chosen to be separable.

Proof: Let $\left\{x_{1}, x_{2}, \ldots,\right\} \subseteq A$ be dense and countable. We then only need to check that $f_{n}\left(x_{n}^{*} x_{n}\right)=\left\|x_{n}\right\|^{2}$. Then put $\pi:=\bigoplus_{n \in \mathbb{N}} \pi f_{n}$.

Remark 5.18: The GNS-Construction (Theorem 5.15) using pure states to irreducible representations.
(i) A representation $\pi: A \rightarrow H$ is called irreducible, if one of the following conditions is satisfied:
(1) $\pi=0$,
(2) $\operatorname{dim} H=1$,
(3) The only closed subspaces $K$ of $H$ such that $\pi(A) K \subseteq K$ (in this case it also holds $\pi(A) K^{\perp} \subseteq K^{\perp}$ ) are 0 and $H$.
(4) $\pi(A)^{\prime}=\mathbb{C} 1$, where $\pi(A)^{\prime}:=\{x \in B(H) \mid x \pi(y)=\pi(y) x \forall y \in A\}$ is the so called commutant.
(5) Every vector $0 \neq \xi \in H$ is cyclic.
(3) - (5) are equivalent conditions.
(ii) If $\operatorname{dim} H<\infty$, then $\pi=\pi_{1} \oplus \cdots \oplus \pi_{n}$, where the $\pi_{i}(1 \leq i \leq n)$ are irriducible.
(iii) $f: A \rightarrow \mathbb{C}$ is called a pure state, if it holds: " $0 \leq g \leq f \Rightarrow \exists \lambda \in[0,1]: g=$ $\lambda f^{\prime \prime}$. If $f$ is a pure state, then $\left(\pi_{f}, H_{f}, \xi_{f}\right)$ is irreducible.
(iv) For all $0 \neq x \in A$ there is an irriducible representation $\pi: A \rightarrow H$ such that $\|\pi(x)\|=\|x\|$.

## 6 von Neumann algebras

Motivation 6.1: In Definition 1.1 we have introduced several locally convex topologies on $B(H)$. A *-subalgebra $A \subseteq B(H)$ that is closed with respect to the operator norm topology is a $C^{*}$-algebra; in fact, each $C^{*}$-algebra arises in this way for a suitable Hilbert space $H$ (see Corollary 5.16). What happens, if $A$ is instead required to be closed with respect to the weak - or the strong operator topology?

Let in the following $(H,\langle\cdot, \cdot\rangle)$ be a complex Hilbert space.
Lemma 6.2: Let $\varphi: B(H) \rightarrow \mathbb{C}$ be a linear functional. Then the following are equivalent:
(i) $\varphi$ is continuous with respect to the weak operator topology.
(ii) $\varphi$ is continuous with respect to the strong operator topology.
(iii) There is $n \in \mathbb{N}$ and there are $\xi_{1}, \ldots, \xi_{n}, \eta_{1}, \ldots, \eta_{n} \in H$ such that

$$
\varphi(x)=\sum_{i=1}^{n}\left\langle x \xi_{i}, \eta_{i}\right\rangle \quad \text { for all } x \in H
$$

Proof: Exercise!
Theorem 6.3: Let $C \subseteq B(H)$ be a convex set. Then $C$ is weakly closed if and only if $C$ is strongly closed.

The proof of this theorem relies on Lemma 6.2 and the following variant of the Hahn-Banach separation theorem, which can be derived from Theorem 2.7 from the Functional Analysis I lecture notes.

Theorem 6.4 (Hahn-Banach): Let $(X, \mathfrak{T})$ be a locally convex topological $\mathbb{C}$-vector space. Suppose that $C \subseteq X$ is a closed convex subset and $x_{0} \in X \backslash C$. Then there exists $\varphi \in(X, \mathfrak{T})^{* 1}$ and $\gamma \in \mathbb{R}$ such that

$$
\operatorname{Re}(\varphi(x))<\gamma<\operatorname{Re}\left(\varphi\left(x_{0}\right)\right) \quad \text { for all } x \in C
$$

Proof (of Theorem 6.3): " $\Rightarrow$ " is clear, since $\mathfrak{T}_{\text {WOT }} \subseteq \mathfrak{T}_{\text {SOT }}$ (see (Remark 1.2)).
" $\Leftarrow$ ": Suppose that $C$ is strongly closed. We assume $C \subsetneq \mathrm{cl}_{\mathrm{WOT}}(C)$, i. e. there is $x_{0} \in \operatorname{cl}_{\text {WOT }}(C) \backslash C$. Thus by Theorem 6.4, we find $\varphi \in\left(B(H), \mathfrak{T}_{\text {SOT }}\right)^{*}$ and $\gamma \in \mathbb{R}$, such that

$$
\operatorname{Re}(\varphi(x))<\gamma<\operatorname{Re}\left(\varphi\left(x_{0}\right)\right) \quad \text { for all } x \in C
$$

[^7]Due to Lemma 6.2, $\varphi$ is also weakly continuous and thus we infer that

$$
\operatorname{Re}(\varphi(x))<\gamma<\operatorname{Re}\left(\varphi\left(x_{0}\right)\right) \quad \text { for all } x \in \operatorname{cl}_{\text {WOT }}(C)
$$

which contradicts $x_{0} \in \operatorname{cl}_{\text {WOt }}(C)$. Hence $C=\operatorname{cl}_{\text {WOT }}(C)$, i. e. $C$ is weakly closed.
Recall that in any topological vector space, the closure of a convex set is again convex. Thus, we deduce from Theorem 6.3:

Corollary 6.5: Let $C \subseteq B(H)$ be a convex set. Then $\operatorname{cl}_{\mathrm{WOT}}(C)=\operatorname{cl}_{\text {SOT }}(C)$.
Proof: - $\mathrm{cl}_{\mathrm{WOT}}(C)$ is convex and weakly closed, thus by Theorem 6.3 it holds that $\mathrm{cl}_{\mathrm{WOT}}(C)$ is strongly closed and thus $\mathrm{cl}_{\text {WOт }}(C) \supseteq \operatorname{cl}_{\text {SOT }}(C)$.

- $\mathrm{cl}_{\mathrm{SOT}}(C)$ is convex and strongly closed, thus by Theorem 6.3 it holds that $\operatorname{cl}_{\mathrm{SOT}}(C)$ is weakly closed and thus $\mathrm{cl}_{\mathrm{SOT}}(C) \supseteq \operatorname{cl}_{\text {WOT }}(C)$.

Thus, since each ${ }^{*}$-subalgebra $A$ of $B(H)$ is in particular convex, we see that $\operatorname{cl}_{\text {WOt }}(A)=\operatorname{cl}_{\text {SOT }}(A)$. For unital $A$, i.e. $1=\operatorname{id}_{H} \in A$, we can handle this huge analytic object by purely algebraic means.

Definition 6.6: Let $S \subseteq B(H)$ be any subset. We call
(i) $S^{\prime}:=\{y \in B(H) \mid \forall x \in S: x y=y x\}$ the commutant of $S$,
(ii) $S^{\prime \prime}:=\left(S^{\prime}\right)^{\prime}$ the bicommutant - or double commutant of $S$.

Lemma 6.7: Let $S \subseteq B(H)$ be any subset.
(i) $S^{\prime}$ is a weakly (and strongly) closed unital subalgebra of $B(H)$.
(ii) If $S=S^{*}$ (i.e. if $x \in S$, then $x^{*} \in S$ ), then $S^{\prime}$ is a weakly closed (and strongly) closed unital *-subalgebra of $B(H)$.
(iii) We have $S \subseteq S^{\prime \prime}$ and $S^{\prime \prime \prime}:=\left(S^{\prime \prime}\right)^{\prime}=S^{\prime}$. If $T \subseteq B(H)$ is another subset, then

$$
S \subseteq T \quad \Longrightarrow \quad T^{\prime} \subseteq S^{\prime}
$$

Proof: Exercise!
Lemma 6.7 tells us that each *-algebra $A \subseteq B(H)$ sits inside the weakly (and strongly) closed unital *-subalgebra $A^{\prime \prime} \subseteq B(H)$; therefore

$$
A \subseteq \operatorname{cl}_{\mathrm{SOT}}(A) \subseteq \operatorname{cl}_{\mathrm{WOT}}(A) \subseteq A^{\prime \prime}
$$

We can say more, if $A$ is unital.
Theorem 6.8 (von Neumann's bicommutant theorem): Let $M \subseteq B(H)$ be a unital *-subalgebra. Then the following statements are equivalent:
(i) $M=M^{\prime \prime}$
(algebraic condition)
(ii) $M$ is weakly closed.
(iii) $M$ is strongly closed.
(analytic condition)
(analytic condition)

Definition 6.9: Let $M \subseteq B(H)$ be a unital *-subalgebra. If $M$ satisfies the equivalent conditions in Theorem 6.8, then $M$ is called a von Neumann algebra (acting on $H$ ).

The theory of von Neumann algebras began with a series of groundbreaking papers "On rings of operators" by Francis J. Murray and John von Neumann that appeared 1936 to 1943.

Theorem 6.10: Let $A \subseteq B(H)$ be a unital *-subalgebra. Then $A^{\prime \prime} \subseteq \operatorname{cl}_{\text {SOT }}(A)$ and thus (Lemma 6.7) we have the equalities

$$
\operatorname{cl}_{\mathrm{SOT}}(A)=\operatorname{cl}_{\mathrm{WOT}}(A)=A^{\prime \prime}
$$

Proof: Take $y \in A^{\prime \prime}$. We have to show that $y \in \operatorname{cl}_{\text {SOT }}(A)$, which means that

$$
\forall \varepsilon>0, \xi_{1}, \ldots, \xi_{n} \in H \exists x \in A:\left\|y \xi_{i}-x \xi_{i}\right\|<\varepsilon \quad \forall i=1, \ldots, n
$$

(1) Consider the case $n=1$. Let $\varepsilon>0$ and $\xi \in H$ be be given. Put

$$
H_{0}:=\operatorname{cl}(A \xi):=\operatorname{cl}(\{x \xi \mid x \in A\}) \subseteq H
$$

and let $p$ be the orthogonal projection onto $H_{0}$. Then $p \in A^{\prime}$ : We clearly have $a H_{0} \subseteq H_{0}$ for all $a \in A$ and also $a H_{0}^{\perp} \subseteq H_{0}^{\perp}$ for all $a \in A$, since for $\zeta \in H_{0}^{\perp}$, $\eta \in H$ and $a \in A$ we have

$$
\langle a \zeta, \eta\rangle=\left\langle\zeta, a^{*} \eta\right\rangle=0
$$

Thus for all $\zeta \in H_{0}^{\perp}, \eta \in H_{0}$ and $a \in A$, we get that

$$
a p(\eta+\zeta)=a \eta=p(a \eta)=p(a \eta+a \zeta)=p a(\eta+\zeta)
$$

thus $a p=p a$ for all $a \in A$.
Now, since $y \in A^{\prime \prime}$, we see that $p y=y p$, which implies

$$
y H_{0}=y p H_{0}=p y H_{0} \subseteq H_{0} .
$$

Because $1 \in A$ implies $\xi=1 \xi \in H_{0}$, we get $y \xi_{0} \in H_{0}=\operatorname{cl}(A \xi)$, i. e. there is $x \in A$ such that $\|y \xi-x \xi\|<\varepsilon$.
(2) The case of $n$ vectors $\xi_{1}, \ldots, \xi_{n} \in H$ will be reduced to (1) by a matrix trick. Consider $H^{n}=\bigoplus_{i=1}^{n} H \ni \xi:=\left(\xi_{1}, \ldots, \xi_{n}\right)^{\top}$ and a mapping

$$
\begin{aligned}
\pi: A & \longrightarrow B(H)^{n} \\
a & \longmapsto \pi(a),
\end{aligned}
$$

where $\pi(a)=\operatorname{diag}(a, \ldots, a)$, i.e. $\pi(a)\left(\eta_{1}, \ldots, \eta_{n}\right)^{\top}:=\left(a \eta_{1}, \ldots, a \eta_{1}\right)^{\top}$. Then $\pi(A)$ is a unital *-subalgebra of $B\left(H^{n}\right)$ and we have $\pi(y) \in \pi(A)^{\prime \prime}$; note that $\pi(A)^{\prime}=\left\{\left(a_{i, j}\right)_{1 \leq i, j \leq n} \mid \forall i, j=1, \ldots, n: a_{i, j} \in A^{\prime}\right\}$.

Thus, applying (1) to $\pi(y) \in \pi(A)^{\prime \prime}$ and $\xi \in H^{n}$ yields $x \in A$ such that

$$
\|\pi(y) \xi-\pi(x) \xi\|<\varepsilon
$$

i. e. $\left\|y \xi_{i}-x \xi_{i}\right\|<\varepsilon$ for $i=1, \ldots, n$.

Proof (of Theorem 6.8): "(i) $\Rightarrow$ (ii)": is the statement that was shown in Lemma 6.7, "(ii) $\Rightarrow$ (iii)" is a consequence of $\mathfrak{T}_{\text {WOT }} \subseteq \mathfrak{T}_{\text {SOT }}$ (as stated in Remark 1.2) and "(iii) $\Rightarrow(\mathrm{i})$ " is the statement Theorem 6.10.

Definition 6.11: Let $S \subseteq B(H)$ be any subset. Put $S^{*}:=\left\{x^{*} \mid x \in S\right\}$. Then $\mathrm{vN}(S):=\left(S \cup S^{*}\right)^{\prime \prime}$ is called the von Neumman algebra generated by $S$; due to Exercise 1 (d), Sheet 5, this is the smallest von Neumann algebra $M \subseteq B(H)$ that contains $S$.

Remark 6.12: (i) Every von Neumann algebra $M \subseteq B(H)$ is also closed with respect to the operator norm topology (see Remark 1.2) and thus a $C^{*}$-algebra. Their general theories, however, are wildly different.
(ii) Obvious examples of von Neumann algebras are $B(H)$ and $\mathbb{C} 1 \subseteq B(H)$. It is less obvious that there are other non-trivial examples.
(iii) For any subset $S=S^{*} \subseteq B(H), S^{\prime}$ is a von Neumann algebra.
(iv) Von Neumann algebras are closed under the measurable functional calculus (see Theorem 11.5 in the Functional Analysis I lecture notes): If $M \subseteq B(H)$ is a von Neumann algebra and $x \in M$ a normal operator, then the continuous functional calculus

$$
\Phi: C(\operatorname{Sp}(x)) \longrightarrow C^{*}(x, 1) \subseteq B(H)
$$

admits a unique extension

$$
\tilde{\Phi}: B_{b}(\operatorname{Sp}(x)) \longrightarrow W^{*}(x, 1)=\mathrm{vN}(x) \subseteq M \subseteq B(H)
$$

such that $\tilde{\Phi}$ is a *-homomorphism with

- $\|\tilde{\Phi}(f)\| \leq\|f\|_{\infty}$ for all $f \in B_{b}(\operatorname{Sp}(x))$,
- $\left(f_{n} \rightarrow f\right.$ pointwise, $f$ bounded $) \Rightarrow\left(\tilde{\Phi}\left(f_{n}\right) \rightarrow \tilde{\Phi}(f)\right.$ in WOT $)$.

In particular, we know that all spectral projections $E_{x}(B)=\tilde{\Phi}\left(\chi_{B}\right)$ for Borel subsets $B \subseteq \operatorname{Sp}(x)$ belong to $\mathrm{vN}(x)$ and hence to $M$. Therefore, von Neumann algebras contain - in contrast to $C^{*}$-algebras - many projections; note that there are $C^{*}$-algebras that contain no projections except 0 and 1 .

It follows that $\operatorname{cl}_{\mathrm{ONT}}(\langle\{p \in M \mid p$ projection $\}\rangle)=M$ : Indeed, if $x \in M$ is selfadjoint, we may approximate $\operatorname{id}_{\mathrm{Sp}(x)} \in C(\mathrm{Sp}(x))$ uniformly by step functions.

Lemma 6.13: Let $M \subseteq B(H)$ be a von Neumann algebra and $x \in B(H)$. Then $x \in M$ holds if and only if $u x=x u$ for all unitary $u \in M^{\prime}$.

Proof: " $\Rightarrow$ ": This is trivial, since $M=M^{\prime \prime}$.
" $\Leftarrow "$ : Every element in a $C^{*}$-algebra $C$ is a linear combination of (at most four) unitaries.

Proof: Let $a=a^{*} \in C$ with $\|a\| \leq 1$. Then we have the decomposition $a=\frac{1}{2}\left(u+u^{*}\right)$ where $u:=a+\mathrm{i} \sqrt{1-a^{2}}$. Each $a$ can be written as $a=\operatorname{Re}(a)+\operatorname{iIm}(a)$ where

$$
\operatorname{Re}(a)=\frac{1}{2}\left(a+a^{*}\right) \quad \text { and } \quad \operatorname{Im}(a)=\frac{1}{2 \mathrm{i}}\left(a-a^{*}\right)
$$

Applying this to $M=M^{\prime}$, yields that $y x=x y$ holds for all $y \in M^{\prime}$, thus $x \in\left(M^{\prime}\right)^{\prime}=M^{\prime \prime}=M$.

Corollary 6.14 (Polar decomposition in von Neumann algebras): Let $M \subseteq B(H)$ be a von Neumann algebra and $x \in M$. Consider the polar decomposition $x=v|x|$, $|x|=\left(x^{*} x\right)^{1 / 2}$, of $x$ in $B(H)$, where $v \in B(H)$ is a partial isometry with the property $\operatorname{ker}(v)=\operatorname{ker}(x)$ (see (Theorem 0.11)). Then $v \in M$ (and clearly $|x| \in M$ ).

Proof: Let $u \in M^{\prime}$ be a unitary. Then:

$$
\begin{aligned}
x=u x u^{*} & =u v|x| u^{*} \\
& =u v u^{*}|x| \quad(|x| \in M \text { and } u y=y u \text { for all } y \in M) \\
& =u v u^{*}|x| .
\end{aligned}
$$

Now $u v u^{*}$ is a partial isometry with

$$
\operatorname{ker}\left(u v u^{*}\right)=u \operatorname{ker}(v)=u \operatorname{ker}(x)=\operatorname{ker}\left(u x u^{*}\right)=\operatorname{ker}(x)
$$

Thus, the uniqueness of the polar decomposition yields that $v=u v u^{*}$ and thus that $v u=u v$. By Lemma 6.13, we get that $v \in M$.

## 7 The Kaplansky density theorem

Motivation 7.1: Let $A \subseteq B(H)$ be a ${ }^{*}$-subalgebra. Consider the strong closure $B:=\operatorname{cl}_{\text {SOT }}(A) \subseteq B(H)$ of $A$. Then, for each $x \in B$ we find a net $\left(x_{\lambda}\right)_{\lambda \in \Lambda}$ in $A$ such that $x_{\lambda} \rightarrow_{\text {SOT }} x$. In general, there is no reason, why $\left(\left\|x_{\lambda}\right\|\right)_{\lambda \in \Lambda}$ should be bounded, but can we make a better choice such that $\sup _{\lambda \in \Lambda}\left\|x_{\lambda}\right\| \leq\|x\|$ ?

Theorem 7.2: Let $f: \mathbb{R} \rightarrow \mathbb{C}$ be a continuous and bounded function. Then $f$ is SOT-continuous, i.e., if $x_{\lambda} \rightarrow_{\mathrm{SOT}} x$ in $B(H)$ and $x_{\lambda}=x_{\lambda}^{*}, x=x^{*}$, then $f\left(x_{\lambda}\right) \rightarrow_{\text {SOT }} f(x)$.

Proof: (1) Consider first the case $f \in C_{0}(\mathbb{R}) \subset C_{b}(\mathbb{R})$. Put

$$
\mathcal{F}_{0}:=\left\{f \in C_{0}(\mathbb{R}) \mid f \text { is SOT-continuous }\right\}
$$

Then $\mathcal{F}_{0} \subseteq C_{0}(\mathbb{R})$ is a subalgebra of $C_{0}(\mathbb{R})$ that is closed with respect to the supnorm (Check!). Furthermore, since $x \mapsto x^{*}$ is continous with respect to SOT on the set of normal operators (note that $x \in B(H)$ is normal if and only if $\|x \xi\|=\left\|x^{*} \xi\right\|$ for all $\xi \in H$ ), we have that $\mathcal{F}_{0}$ forms a ${ }^{*}$-subalgebra of $C_{0}(\mathbb{R})$. Consider now

$$
\begin{aligned}
g: \mathbb{R} & \longrightarrow \mathbb{C} \\
t & \longmapsto \frac{t}{1+t^{2}}
\end{aligned}
$$

Clearly, $g \in C_{0}(\mathbb{R})$. Claim: $g \in \mathcal{F}_{0}$.
Proof: Take a net $\left(x_{\lambda}\right)_{\lambda \in \Lambda}$ of selfadjoint operators in $B(H)$ that converges to $x=x^{*} \in B(H)$ with respect to SOT. Then

$$
\begin{aligned}
& g\left(x_{\lambda}\right)-g(x) \\
& =x_{\lambda}\left(1+x_{\lambda}^{2}\right)^{-1}-x\left(1+x^{2}\right)^{-1} \\
& =\left(1+x_{\lambda}^{2}\right)^{-1}\left[x_{\lambda}\left(1+x^{2}\right)-\left(1-x_{\lambda}^{2}\right) x\right]\left(1+x^{2}\right)^{-1} \\
& =\left(1+x_{\lambda}^{2}\right)^{-1}\left[\left(x_{\lambda}-x\right)+x_{\lambda}\left(x-x_{\lambda}\right) x\right]\left(1+x^{2}\right)^{-1}
\end{aligned}
$$

so that for each $\xi \in H$ we have:

$$
\begin{aligned}
& \left\|g\left(x_{\lambda}\right) \xi-g(x) \xi\right\| \\
& \leq\left\|\left(1+x_{\lambda}^{2}\right)^{-1}\left(x_{\lambda}-x\right)\left(1+x^{2}\right)^{-1} \xi\right\|+\left\|\left(1+x_{\lambda}^{2}\right)^{-1} x_{\lambda}\left(x_{\lambda}-x\right) x\left(1+x^{2}\right)^{-1} \xi\right\| \\
& \leq \underbrace{\left\|\left(1+x_{\lambda}^{2}\right)^{-1}\right\|\left\|\left(x_{\lambda}-x\right)\left(1+x^{2}\right)^{-1} \xi\right\|}_{\leq 1}+\underbrace{\left\|\left(1+x_{\lambda}^{2}\right)^{-1} x_{\lambda}\right\|}_{\rightarrow 1} \underbrace{\left\|\left(x_{\lambda}-x\right) x\left(1+x^{2}\right)^{-1} \xi\right\|}_{\rightarrow 0}
\end{aligned}
$$

because of the properties of the functional calculus and $\left(1+x^{2}\right)^{-1} \xi, x\left(1+x^{2}\right)^{-1} \xi \in H$ are fixed, i.e. $g\left(x_{\lambda}\right) \rightarrow_{\text {SOT }} g(x)$.

Analogously, $f \in \mathcal{F}_{0}$ where $f: \mathbb{R} \rightarrow \mathbb{C}, t \mapsto \frac{1}{1+t^{2}}$. Note that $\{f, g\}$ separates the points of $\mathbb{R}$ and $f(t)>0$ for all $t \in \mathbb{R}$ and thus, by the Stone-Weierstraß theorem for $C_{0}(\mathbb{R})$, we get that $\mathcal{F}_{0}=C_{0}(\mathbb{R})$.
(2) Cosider now the case $f \in C_{b}(\mathbb{R})$. Put

$$
\mathcal{F}:=\{f \in C(\mathbb{R}) \mid f \text { is SOT continuous }\} .
$$

Note that for $h_{1}, h_{2} \in \mathcal{F}$, where one $h_{i}$ is bounded, it holds that $h_{1} h_{2}$ is bounded. Thus in particular: If $h \in \mathcal{F}$ is bounded, then $\operatorname{id} h \in \mathcal{F}$.

Take $h \in C_{b}(\mathbb{R})$. Then, with $f, g$ constructed like in (1), we have $h f, h g \in \mathcal{F}$. Since $f+\mathrm{id} g=1$, we may deduce that

$$
h=h(f+\mathrm{id} g)=h f+\mathrm{id}(h g) \in \mathcal{F}
$$

with the above arguments.
Let $S \subseteq B(H)$ be any subset. In the following we will denote with

$$
S_{s a}:=\left\{x \in S \mid x^{*}=x\right\}
$$

the selfadjoint part of $S$.
Theorem 7.3 (Kaplansky density theorem): Let $A \subseteq B(H)$ be $a^{*}$-algebra. Consider $B:=\operatorname{cl}_{\text {SOT }}(A) \subseteq B(H)$. Then the following statements holds true:
(i) $\operatorname{cl}_{\text {SOT }}\left(A_{s a}\right)=B_{s a}$,
(ii) $\operatorname{cl}_{\text {SOT }}\left(\left\{x \in A_{s a} \mid\|x\| \leq 1\right\}\right)=\left\{x \in B_{s a} \mid\|x\| \leq 1\right\}$,
(iii) $\operatorname{cl}_{\text {SOT }}(\{x \in A \mid\|x\| \leq 1\})=\{x \in B \mid\|x\| \leq 1\}$.

Proof: (i) Since $A_{s a}$ is convex, we have by Corollary 6.5 that $\operatorname{cl}_{\text {SOT }}\left(A_{s a}\right)=$ $\operatorname{cl}_{\text {WOT }}\left(A_{s a}\right)$. Now, since $x \mapsto x^{*}$ is weakly continuous, we get that $\mathrm{cl}_{\mathrm{WOT}}\left(A_{s a}\right) \subseteq$ $B_{s a}$. Conversely, if $x \in B_{s a} \subseteq B$ is given, then we find a net $\left(x_{\lambda}\right)_{\lambda \in \Lambda}$ in $A$ such that $x_{\lambda} \rightarrow_{\text {SOT }} x$. Then $x_{\lambda} \rightarrow_{\text {wot }}$ and thus $x_{\lambda}^{*} \rightarrow_{\text {WOT }} x^{*}=x$. Hence $A_{s a} \ni$ $\operatorname{Re}\left(x_{\lambda}\right) \rightarrow_{\text {WOT }} x$ and thus $x \in \operatorname{cl}_{\text {WOT }}\left(A_{s a}\right)$. In summary we have $\operatorname{cl}_{\text {SOT }}\left(A_{s a}\right)=$ $\operatorname{cl}_{\text {WOT }}\left(A_{s a}\right)=B$.
(ii) Take $x \in B_{s a}$ with $\|x\| \leq 1$. Due to (i), we find a net $\left(x_{\lambda}\right)_{\lambda \in \Lambda}$ in $A_{s a}$ with $x_{\lambda} \rightarrow_{\text {SOT }} x$. Consider $f \in C_{0}(\mathbb{R})$ that is defined by

$$
f(t)= \begin{cases}t, & t \in[-1,1] \\ \frac{1}{t}, & t \notin[-1,1]\end{cases}
$$

From Theorem 7.2, it follows that $\left(\operatorname{cl}_{\mathrm{ONT}}(A)\right)_{s a} \ni f\left(x_{\lambda}\right) \rightarrow_{\mathrm{SOT}} f(x)=x$, where we used that $\operatorname{Sp}(x) \subseteq[-1,1]$, and $\left\|f\left(x_{\lambda}\right)\right\| \leq 1$ since $\|f\|_{\infty} \leq 1$. Now, we have: For all $\lambda \in \Lambda$ and all $n \in \mathbb{N}$ there is $y_{n, \lambda} \in A_{s a}$ such that

$$
\left\|y_{\lambda, n}\right\| \leq 1 \quad \text { and } \quad\left\|f\left(x_{\lambda}\right)-y_{\lambda, n}\right\|<\frac{1}{n}
$$

Thus, $\left(y_{\lambda, n}\right)_{\lambda \in \Lambda, n \in \mathbb{N}}$ is a net in $A_{s a}$ that converges to $x$ in the strong operator topology and satisfies $\left\|y_{\lambda, n}\right\| \leq 1$. This shows " $\supseteq$ ". We leave " $\subseteq$ " as an exercise.

7 The Kaplansky density theorem
(iii) Take $x \in B$ with $\|x\| \leq 1$. Consider

$$
y:=\left(\begin{array}{cc}
0 & x \\
x^{*} & 0
\end{array}\right) \in M_{2}(B) \subseteq B(H \oplus H)
$$

In fact, $y \in M_{2}(B)_{s a}$. Note that a net $\left(y_{\lambda}\right)_{\lambda \in \Lambda}$ with

$$
y_{\lambda}=\left(\begin{array}{ll}
y_{\lambda}^{11} & y_{\lambda}^{12} \\
y_{\lambda}^{21} & y_{\lambda}^{22}
\end{array}\right) \in M_{2}(B)
$$

is strongly convergent to

$$
y=\left(\begin{array}{ll}
y^{11} & y^{12} \\
y^{21} & y^{22}
\end{array}\right) \in M_{2}(B)
$$

if and only if $y_{\lambda}^{i j} \rightarrow_{\text {SOT }} y^{i j}$ for all $1 \leq i, j \leq 2$. Thus, $M_{2}(B)=\operatorname{cl}_{\text {SOT }}\left(M_{2}(A)\right)$. Due to (ii), we find a net $\left(y_{\lambda}\right)_{\lambda \in \Lambda}$ in $M_{2}(A)_{s a}$ such that

$$
\left(\begin{array}{cc}
y_{\lambda}^{11} & y_{\lambda}^{12} \\
y_{\lambda}^{21} & y_{\lambda}^{22}
\end{array}\right)=y_{\lambda} \rightarrow_{\text {SOT }} y=\left(\begin{array}{cc}
0 & x \\
x^{*} & 0
\end{array}\right)
$$

and $\left\|y_{\lambda}\right\| \leq 1$. Hence $x_{\lambda}:=y_{\lambda}^{12} \rightarrow_{\text {SOT }} x$ and $\left\|x_{\lambda}\right\| \leq 1$. The latter follows from

$$
\left\langle y_{\lambda}\binom{0}{\xi},\binom{\eta}{0}\right\rangle=\left\langle x_{\lambda} \xi, \eta\right\rangle
$$

for all $\xi, \eta \in H$. This shows the inclusion " $\supseteq$ ". Again, we leave " $\subseteq$ " as an exercise.

Remark 7.4: (i) Note that (ii) in Theorem 7.3 is not a trivial consequence of (iii) as $x \mapsto x^{*}$ is not continuous with respect to SOT (see (Remark 1.3)).
(ii) One can show that in the situation of Theorem 7.3 also

$$
\operatorname{cl}_{\text {SOT }}(\{x \in A \text { positive } \mid\|x\| \leq 1\})=\{x \in B \text { positive } \mid\|x\| \leq 1\}
$$

and, if $A^{\prime}$ is a $C^{*}$-algebra:

$$
\operatorname{cl}_{\mathrm{SOT}}\left(\left\{x \in A^{\prime} \text { unitary }\right\}\right)=\{x \in B \mid x \text { unitary }\} .
$$

Corollary 7.5: Let $A \subseteq B(H)$ be $a^{*}$-algebra. Then, for each $x \in B:=\operatorname{cl}_{\mathrm{SOT}}(A)$, we find a net $\left(x_{\lambda}\right)_{\lambda \in \Lambda}$ in $A$ such that

$$
x_{\lambda} \rightarrow \text { SOT } x \quad \text { and } \quad \sup _{\lambda \in \Lambda}\left\|x_{\lambda}\right\| \leq\|x\| \text {. }
$$

If $x$ is selfadjoint, then $\left(x_{\lambda}\right)_{\lambda \in \Lambda}$ can be chosen to consist of selfadjoint operators $x_{\lambda}$.

If the underlying Hilbert space $(H,\langle\cdot, \cdot\rangle)$ is separable, then there exists even a sequence $\left(x_{n}\right)_{n \in \mathbb{N}}$ in $A$ with the above properties.

Proof: Exercise!

## 8 Abelian von Neumann algebras

Motivation 8.1: Let $A$ be a commutative unital $C^{*}$-algebra. According to the Gelfand-Naimark Theorem (Theorem 10.16 from the Functional Analysis I lecture notes), the Gelfand transformation

$$
\chi: A \longrightarrow C(\operatorname{Spec}(A)), \quad x \longmapsto \hat{x}
$$

where $\hat{x}(\varphi):=\varphi(x)$, gives an isometric *-isomorphism, i.e., $A$ is of the form $C(K)$ for some compact Hausdorff space $K$. Therefore, the theory of $C^{*}$-algebras can be seen as a kind of "non-commutative topology".

Analogously, von Neumann algebra theory is considered as a kind of "noncommutative measure theory"; in order to understand this point of view, we study here abelian (i.e., commutative) von Neumann algebras.

Remark 8.2: Since every abelian von Neumann algebra is particular a commutative $C^{*}$-algebra, it is isomorphic to $C(K)$ for some compact Hausdorff space $K$; this space, however, is extremally disconnected (i.e., the closure of every open set is again open). This highlights that von Neumann algebras are rather exceptional among $C^{*}$-algebras.

Definition 8.3: Let $M \subseteq B(H)$ be a von Neumann algebra and let $0 \neq \xi \in H$ be given. We say that
(i) $\xi$ is cyclic for $M$, if $M \xi=\{x \xi \mid x \in M\}$ is dense in $H$,
(ii) $\xi$ is separating for $M$, if $x \xi \neq 0$ for all $0 \neq x \in M$.

Theorem 8.4: Let $M \subseteq B(H)$ be a von Neumann algebra and let $0 \neq \xi \in H$ be given. Then $\xi$ is cyclic for $M$ if and only if $\xi$ is separating for $M^{\prime}$.

Proof: " $\Rightarrow$ ": Let $\xi$ be cyclic for $M$. Suppose that $x \xi=0$ for some $x \in M^{\prime}$. Then for all $y \in M$ we have

$$
x y \xi=y x \xi=y 0=0 .
$$

Since $M \xi$ is dense in $H$, we infer that even $x \eta=0$ for all $\eta \in H$, hence $x=0$. Thus, $\xi$ is separation for $M^{\prime}$.
" $\Leftarrow "$ : Let $\xi$ be separating for $M^{\prime}$. Assume that $M \xi$ was not dense in $H$. Then the orthogonal projection $p$ onto $\operatorname{cl}(M \xi)^{\perp}$ was non-zero but we had

$$
0=\langle p \xi, \xi\rangle=\langle p \xi, p \xi\rangle=\|p \xi\|^{2}, \quad\left(p \xi \in M \xi^{\perp}\right)
$$

i. e., $p \xi=0$ and $p \in M^{\prime}$ (which follows from step (1) in the proof of Theorem 6.10) which was a contradiction.

Corollary 8.5: If $M \subseteq B(H)$ is an abelian von Neumann algebra and $\xi \in H$ cyclic for $M$, then $\xi$ is also separating for $M$.

Proof: If $\xi$ is cyclic for $M$, then $\xi$ is separating for $M^{\prime}$ due to Theorem 8.4. Now, because being separating passes to von Neumann subalgebras, the assertion follows since $M \subseteq M^{\prime}$.

Definition 8.6: A von Neumann algebra $M \subseteq B(H)$ is said to be separable, if the underlying Hilbert space $H$ is separable (see Remark 0.3 (vii)).

Theorem 8.7: Let $M \subseteq B(H)$ be a separable abelian von Neumann algebra, then there exists a separating vector $0 \neq \xi \in H$ for $M$.

Proof: By Zorns Lemma, there exists a maximal family of non-zero unit vectors $\left(\xi_{\lambda}\right)_{\lambda \in \Lambda}$ such that $M \xi_{\lambda} \perp M \xi_{\mu}$ for all $\lambda, \mu \in \Lambda$ where $\lambda \neq \mu$. Since $H$ is separable, that family is countable, say $\left(\xi_{n}\right)_{n=1}^{R}, R \in \mathbb{N} \cup\{\infty\}$. Its maximality implies that

$$
H_{0}:=\left\{x \xi_{n} \mid n \in \mathbb{N}, n \leq R, x \in M\right\} \subseteq H
$$

is dense in $H$. Denote by $p_{n}$ the orthogonal projection onto $\operatorname{cl}\left(M \xi_{n}\right)$. Then (see step (1) in the proof of Theorem 6.10) we have $p_{n} \in M^{\prime}$. Put

$$
\xi:=\sum_{n=1}^{R} \frac{1}{2^{n}} \xi_{n} \in H \backslash\{0\}
$$

(note that $\left(\xi_{n}\right)_{n=1}^{R}$ is orthonormal). Then $\xi$ is separating for $M$ : Take $x \in M$ such that $x \xi=0$. Then for all $n \in \mathbb{N}, n \leq R$ we have

$$
0=p_{n}(x \xi)=x p_{n} \xi=x \frac{1}{2^{n}} \xi^{n}=\frac{1}{2^{n}} x \xi_{n}
$$

so for all $n \in \mathbb{N}$ we have $x \xi_{n}=0$, i. e., for all $y \in M$ and $n \in \mathbb{N}, n \leq R$ it holds $0=y\left(x \xi_{n}\right)=x\left(y \xi_{n}\right)$. Thus for all $\eta \in H_{0}$ it holds $x \eta=0$ and as $H_{0}$ is dense in $H$, we get $x=0$.

Definition 8.8: An abelian von Neumann algebra $M \subseteq B(H)$ is called maximal, if $M \subseteq N \subseteq B(H)$ for an abelian von Neumann algebra $N \subseteq B(H)$ implies that $M=N$.

Lemma 8.9: An abelian von Neumann algebra $M \subseteq B(H)$ is maximal if and only if $M^{\prime}=M$ holds.

Proof: Exercise.
Corollary 8.10: Let $M \subseteq B(H)$ be a separable maximal abelian von Neumann algebra. Then there exists a cyclic vector $0 \neq \xi \in H$ for $M$.

Proof: By Theorem 8.7, there is a separating vector $0 \neq \xi \in H$ for $M$. Note that $M=M^{\prime}$ due to Lemma 8.9. Thus, Theorem 8.4 implies that $\xi$ is also cyclic for $M$.

Example 8.11: Let $K$ be a compact Hausdorff space and let $\mu$ be a finite Borel measure on $K$. For each function $f \in L^{\infty}(K, \mu)$, we define $M_{f} \in B\left(L^{2}(K, \mu)\right)$ by $M_{f} g:=f g$ for all $g \in L^{2}(K, \mu)$. We then have that $\left\|M_{f}\right\|=\|f\|_{L^{\infty}(K, \mu)}$ for all $f \in L^{\infty}(K, \mu)$. Consider the $C^{*}$-algebra

$$
A:=\left\{M_{f} \mid f \in C(K)\right\}
$$

Then we have the statements:
(i) $M:=A^{\prime}=\left\{M_{f} \mid f \in L^{\infty}(K, \mu)\right\}$, which is a von Neumann algebra acting on $L^{2}(K, \mu)$.
(ii) $A^{\prime \prime}=A^{\prime}$ and thus $M^{\prime}=M$, i.e., $M$ is maximal abelian.
(iii) The constant function 1 is cyclic and separating for $M$.

Theorem 8.12: Let $M \subseteq B(H)$ be an abelian von Neumann algebra and suppose that there is a cyclic vector $0 \neq \xi \in H$ for $M$. Then, for any SOT-dense unital $C^{*}$-subalgebra $A \subseteq M$ there exists a finite Radon measure $\mu$ on $K:=\operatorname{Spec}(A)$ with $\operatorname{supp}(\mu)=K$ and an unitary $U: L^{2}(K, \mu) \rightarrow H$ such that

$$
U^{*} M U=\left\{M_{f} \mid f \in L^{\infty}(K, \mu)\right\}\left(\cong L^{\infty}(K, \mu)\right) \subseteq B\left(L^{2}(K, \mu)\right)
$$

Remark 8.13: A finite Borel measure $\mu$ on a Hausdorff space $X$ is a finite Radon measure, if it is inner regular, i. e., if for each Borel set $B \in \mathfrak{B}(X)$ we have

$$
\mu(B)=\sup _{K \subseteq B \text { compact }} \mu(K)
$$

Those Radon measures are automatically outer regular, i.e., we have for each $B \in \mathfrak{B}(X)$ that

$$
\mu(B)=\inf _{B \subseteq U \text { open }} \mu(U)
$$

We define the support of a finite Radon measure by

$$
\operatorname{supp}(\mu):=X \backslash V
$$

where $V:=\bigcup\{U \subseteq X$ open $\mid \mu(U)=0\}$.
Proof: (1) By the Riesz representation theorem, we find a Radon measure $\mu$ on the compact Hausdorff space $K:=\operatorname{Spec}(A)$ such that

$$
\left\langle\chi^{-1}(f) \xi, \xi\right\rangle=\int_{K} f(x) d \mu(x)
$$

for all $f \in C(K)$, where $\chi^{-1}: C(K) \rightarrow A$ is the inverse Gelfand transform, which is positive. Define

$$
U_{0}: C(K) \longrightarrow H, \quad f \longmapsto \chi^{-1}(f) \xi
$$

Then, for all $f \in C(K)$, we have

$$
\left\|U_{0} f\right\|^{2}=\left\langle\chi^{-1}(f) \xi, \chi^{-1}(f) \xi\right\rangle=\left\langle\chi^{-1}\left(|f|^{2}\right) \xi, \xi\right\rangle=\int_{K}|f(x)|^{2} d \mu(x)=\|f\|_{L^{2}(K, \mu)}^{2}
$$

so that $U_{0}$ extends to an isometry $U: L^{2}(K, \mu) \rightarrow H$.
(2) Claim: $U\left(L^{2}(K, \mu)\right)=H$, i. e., $U$ is a unitary. Note that $A \xi=U_{0}(C(K)) \subseteq$ $U\left(L^{2}(K, \mu)\right)$ and $\operatorname{cl}(A \xi)=\operatorname{cl}(M \xi)=H$ (as $A \subseteq M$ is SOT-dense and $\xi$ is cyclic), thus $U\left(L^{2}(K, \mu)\right)=H$.
(3) Claim: $\operatorname{supp}(\mu)=K$. Otherwise, there would be some $f \in C(K)$ such that

$$
0=\int_{K}|f(x)|^{2} d \mu(x)=\left\|\chi^{-1}(f) \xi\right\|^{2} \quad \Longrightarrow \quad \chi^{-1}(f) \xi=0
$$

which contradicts the fact that $\xi$ is also separating for $M$ due to Corollary 8.5.
(4) Claim: For all $f \in C(K) \subseteq L^{\infty}$ it holds: $U^{*} \chi^{-1}(f) U=M_{f}$. Take any $g \in C(K) \subseteq L^{2}(K, \mu)$. Then

$$
U^{*} \chi^{-1}(f) U g=U^{*} \chi^{-1}(f) \chi^{-1}(g) \xi=U^{*} \chi^{-1}(f g) \xi=f g=M_{f} g
$$

Since $C(K)$ is dense in $L^{2}(K, \mu)$, we conclude that $U^{*} \chi^{-1}(f) U=M_{f}$. Thus, we see that $U^{*} A U=\left\{M_{f} \mid f \in C(K)\right\}$.
(5) Claim: $U^{*} M U=\left\{M_{f} \mid f \in L^{\infty}(K, \mu)\right\}$. Consider the ${ }^{*}$-isomorphism

$$
\operatorname{Ad}(U): B\left(L^{2}(K, \mu)\right) \longrightarrow B(H), \quad x \longmapsto U x U^{*}
$$

with inverse given by $\operatorname{Ad}\left(U^{*}\right): B(H) \rightarrow B\left(L^{2}(K, \mu)\right)$. Observe that $\operatorname{Ad}(U)$ and $\operatorname{Ad}\left(U^{*}\right)=(\operatorname{Ad}(U))^{-1}$ are obviously SOT-continuous. Thus, since cl $\mathrm{l}_{\text {SOT }}(A)=M$ and also $\operatorname{cl}_{\text {SOT }}\left(\left\{M_{f} \mid f \in C(K)\right\}\right)=\left\{M_{f} \mid f \in L^{\infty}(K, \mu)\right\}$ (where this equality is shown in Example 8.11) we infer from $\operatorname{Ad}\left(U^{*}\right)(A)=\left\{M_{f} \mid f \in C(K)\right\}$ (which holds by (4)), that

$$
\operatorname{Ad}\left(U^{*}\right)(M)=\left\{M_{f} \mid f \in L^{\infty}(K, \mu)\right\}
$$

as desired.

Remark 8.14: In Corollary 8.10, we have seen that if a von Neumann algebra is maximal abelian then it has a cyclic vector $\xi \in H$ - under the hypothesis that $H$ is separable. Theorem 8.12 (if combined with Example 8.11) says that the converse is also true - even without the separability hypothesis.

Theorem 8.15: Let $M \subseteq B(H)$ be a separable abelian von Neumann algebra. Then there exists a compact Hausdorff space $K$ with a finite Radon measure $\mu$ such that $M$ and $L^{\infty}(K, \mu)$ are ${ }^{*}$-isomorphic.

Proof: By Theorem 8.7 there exists a separating vector $\xi$ for $M$. Consider the sub Hilbert space $H_{0}:=\operatorname{cl}(M \xi) \subseteq H$. Then

$$
\varphi: M \longrightarrow B\left(H_{0}\right),\left.\quad x \longmapsto x\right|_{H_{0}}
$$

is a well-defined (since $x$ maps $H_{0}$ to itself) and injective (since $\xi_{0} \in H_{0}$ is separating for $M)^{*}$-homomorphism that is unital. Thus $\varphi$ is isometric (Proposition 4.10) so that $\varphi(M)$ is a $C^{*}$-algebra on $H_{0}$. Now consider $M_{0}=\operatorname{cl}_{\text {SOT }}(\varphi(M)) \subseteq B\left(H_{0}\right)$. Then $\xi$ is cyclic for $M_{0}$, because $M_{0} \xi \supseteq \varphi(M) \xi=M \xi$ and $\operatorname{cl}(M \xi)=H_{0}$, hence we may apply Theorem 8.12 which gives $U^{*} M_{0} U=\left\{M_{f} \mid f \in L^{\infty}(K, \mu)\right\}$ and thus a ${ }^{*}$-isomorphism $\psi_{0}: M_{0} \rightarrow L^{\infty}(K, \mu)$ with

$$
\psi_{0}^{-1}: L^{\infty}(K, \mu) \longrightarrow M_{0}, \quad f \longmapsto U M_{f} U^{*}
$$

We are done, if we can prove that infact $M_{0}=\varphi(M)$. To see this, take $y \in M_{0}$. By Corollary 7.5 , we find a net $\left(y_{\lambda}\right)_{\lambda \in \Lambda}$ in $\varphi(M)$ such that $y_{\lambda} \rightarrow y$ in the strong operator topology and $\left\|y_{\lambda}\right\| \leq\|y\|$ for all $\lambda \in \Lambda$. Thus, we find a net $\left(x_{\lambda}\right)_{\lambda \in \Lambda}$ in $M$ such that $y_{\lambda}=\varphi\left(x_{\lambda}\right)$ and

$$
\left\|x_{\lambda}\right\| \leq\left\|\varphi\left(x_{\lambda}\right)\right\|=\left\|y_{\lambda}\right\| \leq\|y\|
$$

for all $\lambda \in \Lambda$. Since $\{x \in M \mid\|x\| \leq r\}$ is compact with respect to the weak operator topology for all $r>0$, there is a subnet $\left(x_{\lambda(\gamma)}\right)_{\gamma \in \Gamma}$ of $\left(x_{\lambda}\right)_{\lambda \in \Lambda}$ such that $x_{\lambda(\gamma)} \rightarrow$ wот $x$ for some $x \in M(\|x\| \leq\|y\|)$. Thus, since $\varphi$ is WOT-continuous, $\varphi\left(x_{\lambda(\gamma)}\right) \rightarrow_{\text {wот }} \varphi(x)$ and $\varphi\left(x_{\lambda(\gamma)}\right) \rightarrow_{\text {WOт }} y$ (since $\varphi\left(x_{\lambda}\right) \rightarrow_{\text {SOT }} y$ ).

Thus by the uniqueness of the limit $y=\varphi(x) \in \varphi(M)$ which is what we wanted to show.

Remark 8.16: (i) We have used $\{x \in B(H) \mid\|x\| \leq r\}=: B(H)_{r}$ and hence $\{x \in M \mid\|x\| \leq r\}=M \cap B(H)_{r}$ are compact with respect to the weak operator topology for all $r>0$. This can be seen as follows: Consider the map

$$
\iota: B(H)_{r} \longrightarrow \prod_{\zeta, \eta \in H}\{z \in \mathbb{C}| | z \mid<r\|\zeta\|\|\eta\|\}=: K, \quad x \longmapsto(\langle x \zeta, \eta\rangle)_{\zeta, \eta \in H}
$$

It is clearly injecitve and moreover continuous if $B(H)_{r}$ is endowed with the weak operator topology and $K$ with the product topology; in fact, one can show that $\iota: B(H)_{r} \rightarrow \iota\left(B(H)_{r}\right)$ is a homeomorphism. By Tychonoff's theorem, $K$ is compact; thus it suffices to verify that $\iota\left(B(H)_{r}\right)$ is closed. If $\left(x_{\lambda}\right)_{\lambda \in \Lambda}$ is a net in $B(H)_{r}$ for which $\left(\iota\left(x_{\lambda}\right)\right)_{\lambda \in \Lambda}$ converges to some point $z=\left(z_{\zeta, \eta}\right)_{\zeta, \eta \in H}$, one finds a linear map $x: H \rightarrow H$ such that $\langle x \zeta, \eta\rangle=z_{\zeta, \eta}$ for all $\zeta, \eta$ in $H$. Thus, it follows $\|x\| \leq r$ since

$$
|\langle x \zeta, \eta\rangle|=\left|z_{\zeta, \eta}\right| \leq r\|\zeta\|\|\eta\| .
$$

Hence $x \in B(H)_{r}$ with $\iota(x)=z$.
(ii) If $H$ is separable and $x \in B(H)$ is normal, we may apply Theorem 8.12 to $A=C^{*}(x, 1)$ and $M=\mathrm{vN}(x)$, which yields a *-isomorphism $L^{\infty}(\operatorname{Spec}(A), \mu) \rightarrow M$ for some finite Radon measure $\mu$. Since $\operatorname{Spec}(A) \cong \operatorname{Sp}(x)$ due to Lemma 10.22 (from the Functional Analysis I lecture notes), we obtain a *-ismorphism

$$
\tilde{\Phi}: L^{\infty}(\mathrm{Sp}(x), \nu) \longrightarrow M
$$

with the induced measure $\nu$ on $\operatorname{Sp}(x)$, that agrees with the measurable functional calculus; see (Remark 6.12) (iv).
(iii) Let $M \subseteq B(H)$ be a separable abelian von Neumann algebra. One can show that there exists a selfadjoint element $x=x^{*} \in M$ such that $M=\mathrm{vN}(x)$.
(iv) From (iii), one can deduce that for each separable abelian von Neumann algebra $M \subseteq B(H)$ a countable (possibly empty) set $I$ exists such that $M$ is *-isomorphic to either $\ell^{\infty}(I)$ or $L^{\infty}\left([0,1], \lambda^{1}\right)^{1}$.
(v) Theorem 8.15 stays valid of the separability condition is removed. In that case one has to decompose $H$ as $H=\bigoplus_{\lambda \in \Lambda} \operatorname{cl}\left(M \xi_{\lambda}\right)$ for an orthonormal family $\left(\xi_{\lambda}\right)_{\lambda \in \Lambda}$ like in the proof of Theorem 8.7. Let $p_{\lambda}$ be the orthogonal projection onto $H_{\lambda}:=\operatorname{cl}\left(M \xi_{\lambda}\right)$. We know $p_{\lambda} \in M^{\prime}$; one can show that $M p_{\lambda} \subseteq B\left(H_{\lambda}\right)$ is an abelian von Neumann algebra with cyclic vector $\xi_{\lambda}$. Thus, by Theorem 8.15, $M p_{\lambda}$ is ${ }^{*}$-isomorphic to $L^{\infty}\left(K_{\lambda}, \mu_{\lambda}\right)$ for some probability space $\left(K_{\lambda}, \mu_{\lambda}\right)$. Then $M$ is ${ }^{*}$-isomorphic to

$$
\bigoplus_{\lambda \in \Lambda} L^{\infty}\left(K_{\lambda}, \mu_{\lambda}\right) \cong L^{\infty}\left(\prod_{\lambda \in \Lambda} K_{\lambda}, \prod_{\lambda \in \Lambda} \mu_{\lambda}\right)
$$

[^8]
## 9 Projections, factors and their type classification

Motivation 9.1: We have already seen (see Remark 6.12 (iv)), that von Neumann algebras contain typically "many" projections, i. e., operators $e$ that satisfy $e^{2}=$ $e=e^{*}$ (see Definition $0.10(\mathrm{vi})$ ). Thus, by studying projections, we gain some deeper understanding of the structure of von Neumann algebras.

Definition 9.2: Let $M \subseteq B(H)$ be a von Neumann algebra and let $e, f \in M$ be two projections. We say that
(i) $e$ and $f$ are equivalent $(e \sim f)$, if there exists a partial isometry $u \in M$ (see Definition 0.10 (vii)), such that

$$
u^{*} u=e \quad \text { and } \quad u u^{*}=f
$$

This is also called Murray-von-Neumann equivalence.
(ii) $e$ is subordinate to $f(e \precsim f)$, if there exists a projection $g \in M$ such that $e \sim g \leq f$.

Remark 9.3: (i) The partial isometry that provides the equivalence must belong to $M$; " $M$ knows that the projections are equivalent". Thus, both $\sim$ and $\precsim$ are relative to $M$. In fact, we can have $f \sim e$ in $B(H)$ but $f \nsim e$ in $M$ (namely if $u \notin M$ ).
(ii) For $M=B(H)$, we have the following:

$$
\begin{aligned}
& e \sim f \Leftrightarrow \operatorname{dim} e H=\operatorname{dim} f H \\
& e \precsim f \Leftrightarrow \operatorname{dim} e H \leq \operatorname{dim} f H
\end{aligned}
$$

(iii) $g \leq f$ means $f-g$ is positive. Note that

$$
f-g \geq 0 \Leftrightarrow g f=g \Leftrightarrow f g=g \Leftrightarrow g H \subseteq f H
$$

(iv) " $\sim$ " is an equivalence relation, i. e.,

- $e \sim e$
- $e \sim f \Leftrightarrow f \sim e$
- $(e \sim f \wedge f \sim g) \Rightarrow e \sim g$
" " is easily seen to be a preorder, i.e.,
- $e \precsim e$,
- $(e \precsim f \wedge f \precsim g) \Rightarrow e \precsim g$.

We will see that "ゐ" is in fact a partial order, i.e., it satsifies also

- $(e \precsim f \wedge f \precsim e) \Rightarrow e \sim f$.
(antisymmetry)
Under a further condition on $M$, " $\precsim$ " is even a total order on

$$
\mathcal{P}(M):=\{p \in M \mid p \text { projection }\}
$$

i. e., any two $e, f \in \mathcal{P}(M)$ are comparable in the sense that $e \precsim f$ or $f \precsim e$ must hold.

Theorem 9.4: (i) Let $\left(e_{i}\right)_{i \in I}$ be a family of mutually orthogonal projections in $B(H)$ (i.e., we have $e_{i} e_{j}=0$ for all $i, j \in I$ with $i \neq j$ ). Then $\sum_{i \in I} e_{i}$ converges in the strong operator topology to a projection $e$.
(ii) Let $\left(e_{i}\right)_{i \in I},\left(f_{i}\right)_{i \in I}$ be two families of mutually orthogonal projections in a von Neumann algebra $M \subseteq B(H)$. Then:

- $e_{i} \sim f_{i} \forall i \in I \Rightarrow \sum_{i \in I} e_{i} \sim \sum_{i \in I} f_{i}$,
- $e_{i} \precsim f_{i} \forall i \in I \Rightarrow \sum_{i \in I} e_{i} \precsim \sum_{i \in I} f_{i}$.

Proof: (i) Put $K_{i}:=e_{i} H$ and $K:=\bigoplus_{i \in I} K_{i} \subseteq H$. Let $e$ be the projection onto $K$. We want to show that

$$
\left(\sum_{i \in F} e_{i}\right)_{F \complement_{\mathrm{fin}} I} \rightarrow_{\mathrm{SOT}} e
$$

(1) For $\xi \in K^{\perp}$, we have $\xi \in K_{i}^{\perp}$ for all $i \in I$. Thus, $\sum_{i \in F} e_{i} \xi=0$ for all $F \subseteq_{\text {fin }} I$, thus $\left(\sum_{i \in F} e_{i} \xi\right)_{F \complement_{\text {fin }} I}$ converges to $0=e \xi$ in $H$.
(2) For $\xi \in K$, we find by definition of $K$ a summable family $\left(\xi_{i}\right)_{i \in I}$ with $\xi_{i} \in K_{i}$ for $i \in I$, that has the sum $\xi$, i. e., by Remark 5.23 (ii) from the Functional Analysis I lecture notes we get

$$
\left(\sum_{i \in F} \xi_{i}\right)_{F \complement_{\operatorname{fin}} I}=\left(\sum_{i \in F} e_{i} \xi\right)_{F \complement_{\operatorname{fin}} I} \quad\left(\text { as } \xi_{i}=e_{i} \xi\right)
$$

converges to $\xi=e \xi$ in $H$.
Thus, for each $\xi \in K \oplus K^{\perp}=H$, the net $\left(\sum_{i \in F} e_{i} \xi\right)_{F \complement_{\text {fin }} I}$ converges to $e \xi$ in $H$.
(ii) For each $i \in I$, we find a partial isometry $u_{i} \in M$ such that $u_{i}^{*} u_{i}=e_{i}$ and $u_{i} u_{i}^{*}=f_{i}$. Define $u \in B(H)$ by

$$
\left.u\right|_{K_{i}}:=\left.u_{i}\right|_{K_{i}} \quad \text { and }\left.\quad u\right|_{K^{\perp}}:=0 .
$$

Then $u$ is a partial isometry with $u^{*} u=\sum_{i \in I} e_{i}$ and $u u^{*}:=\sum_{i \in I} f_{i}$. Moreover, for each element $y \in M^{\prime}$, we have

$$
\left.y u\right|_{K_{i}}=\left.y u_{i}\right|_{K_{i}}=\left.u_{i} y\right|_{K_{i}}=\left.u_{i} y\right|_{K_{i}}=\left.u y\right|_{K_{i}} \quad\left(\text { as } y \in M^{\prime} \text { and } y e_{i}=e_{i}\right)
$$

for all $i \in I$, thus $y u=u y$; this means $u \in M^{\prime \prime}=M$.

For the second part: For each $i \in I$ we find a projection $g_{i} \in M$ such that $e_{i} \sim g_{i} \leq f_{i}$. Then (via the previous result):

$$
\sum_{i \in I} e_{i} \sim \sum_{i \in I} g_{i} \leq \sum_{i \in I} f_{i},
$$

i. e., $\sum_{i \in I} \precsim \sum_{i \in I} f_{i}$.

Theorem 9.5 (Cantor-Bernstein-Schröder theorem for projections): Let $M \subseteq B(H)$ be a von Neumann algebra and let e, $f \in M$ be two projections. Then:

$$
e \sim f \Longleftrightarrow e \precsim f \text { and } f \precsim e .
$$

Thus, $\precsim$ is a partial order on $\mathcal{P}(M)$ (or, more precisely, on $\mathcal{P}(M) / \sim$ ).
Proof: " $\Rightarrow$ " is obvious, since $e \sim f \leq f$ and $f \sim e \leq e$.
Now for " $\Leftarrow$ ": By assumption, we have

$$
\begin{aligned}
& e \precsim f \Rightarrow \exists u \in M \text { partial isometry : } e=u^{*} u \text { and } u u^{*} \leq f \\
& f \precsim e \Rightarrow \exists v \in M \text { partial isometry : } f=v^{*} v \text { and } v v^{*} \leq e
\end{aligned}
$$

Note that $\left.u\right|_{e H}: e H \rightarrow f H$ and $\left.v\right|_{f H}: f H \rightarrow e H$ are both isometries (see Definition 0.10 (vii)). Define:

- $e_{0}:=e-v^{*} v \leq e, f_{0}:=u e_{0} e^{*} \leq u u^{*} \leq f$ - note that

$$
f_{0}^{2}=u e_{0} u^{*} u e_{0} u^{*}=u e_{0} e e_{0} u^{*}=u e_{0} e^{*}=f_{0} .
$$

- $e_{n}:=v f_{n-1} v^{*} \leq v v^{*} \leq e, f_{n}:=u e_{n} u^{*} \leq u u^{*} \leq f$.
- $e_{\infty}:=e-\sum_{n \in \mathbb{N}} e_{n}, f_{\infty}:=f-\sum_{n \in \mathbb{N}} f_{n}$

Note that $\left(e_{n}\right)_{n \in \mathbb{N}}$ are mutually orthogonal projections in $M$ and $\left(f_{n}\right)_{n \in \mathbb{N}}$ are mutually orthogonal projections in $M$. Furhtermore note that the convergence of the series $\sum_{n \in \mathbb{N}} e_{n}, \sum_{n \in \mathbb{N}} f_{n}$ is guaranteed by Theorem 9.4 (i).
(1) Claim: $e_{n} \sim f_{n}$ for all $n \in \mathbb{N}$.

Indeed:

- $\left(u e_{n}\right)^{*}\left(u e_{n}\right)=e_{n} u^{*} u e_{n}=e_{n} e e_{n}=e_{n}^{2}=e_{n}$,
- $\left(u e_{n}\right)\left(u e_{n}\right)^{*}=u e_{n} e_{n} u^{*}=u e_{n} u^{*}=f_{n}$
where $u e_{n}$ is a partial isometry as $\left.u e_{n}\right|_{e_{n} H}: e_{n} H \rightarrow e H$ is isometric and $\left.u e_{n}\right|_{\left(e_{n} H\right)^{\perp}}=0$ (since $\left.\left(e_{n} H\right)^{\perp}=\left(1-e_{n}\right) H\right)$ that belongs to $M$.
(2) Claim: $e_{\infty} \sim f_{\infty}$.

Consider $v_{N}=v^{*}\left(e-\sum_{n=0}^{N} e_{n}\right)$. Then

$$
v_{N} v_{N}^{*}=v^{*}\left(e-\sum_{n=0}^{N} e_{n}\right) v=v^{*}\left(e-e_{0}\right) v-\sum_{n=1}^{N} v^{*} e_{n} v
$$

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$$
=v^{*} v v^{*} v-\sum_{n=1}^{N}\left(v^{*} v\right) f_{n-1}\left(v^{*} v\right)=f-\sum_{n=0}^{N-1} f_{n}
$$

and

$$
v_{N}^{*} v_{N}=\left(e-\sum_{n=0}^{N} e_{n}\right) v v^{*}\left(e-\sum_{n=0}^{N} e_{n}\right)=e-\sum_{n=0}^{N} e_{n},
$$

since $e-\sum_{n=0}^{N} e_{n}=\left(e-e_{0}\right)-\sum_{n=1}^{N} e_{n} \leq e-e_{0}=v v^{*}$. Thus, as $N \rightarrow \infty$, $v_{N} \rightarrow_{\text {SOT }} v_{\infty}:=v^{*} e_{\infty}$ (Theorem 9.4 (i)) and

$$
v_{\infty}^{*} v \text { WOT } \leftarrow v_{N}^{*} v_{N}=e-\sum_{n=0}^{N} e_{n} \rightarrow_{\text {SOT }} e_{\infty} \quad \text { as } N \rightarrow \infty
$$

thus $v_{\infty}^{*} v_{\infty}=e_{\infty}$ and analogously $v_{\infty} v_{\infty}^{*}=f_{\infty}$. Thus, by Theorem 9.4 (ii) it holds

$$
e=e_{\infty}+\sum_{n \in \mathbb{N}} e_{n} \sim f_{\infty}+\sum_{n \in \mathbb{N}} f_{n}=f
$$

Theorem 9.6: Let $M \subseteq B(H)$ be a von Neumann algebra. For $x \in M$, we denote by $\operatorname{supp}(x)$ the orthogonal projection in $B(H)$ onto $\operatorname{ker}(x)^{\perp}=\operatorname{cl}\left(\operatorname{im}\left(x^{*}\right)\right)$. Then $\operatorname{supp}(x) \in M$ and moreover $\operatorname{supp}(x) \sim \operatorname{supp}\left(x^{*}\right)$.

Proof: Consider the polar decomposition $x=v|x|$ of $x$ from Corollary 6.14. Then $v \in M$ is a partial isometry and satisfies (see Theorem 0.11) $v^{*} v=$ projection onto $\operatorname{ker}(x)^{\perp}=\operatorname{supp}(x)$ and $v v^{*}=\operatorname{projection}$ onto $\operatorname{cl}\left(\operatorname{im}\left(x^{*}\right)\right)=\operatorname{supp}\left(x^{*}\right)$. Thus, $\operatorname{supp}(x), \operatorname{supp}\left(x^{*}\right) \in M$ and $\operatorname{supp}(x) \sim \operatorname{supp}\left(x^{*}\right)$.

Definition 9.7: Let $M \subseteq B(H)$ be a von Neumann algebra.
(i) We call $Z(M):=M \cap M^{\prime}$ the center of $M$,
(ii) If $Z(M)=\mathbb{C} 1$, then $M$ is called a factor.

Remark 9.8: Every separable von Neumann algebra $M \subseteq B(H)$ can be decomposed as a direct integral

$$
M \cong \int_{X}^{\oplus} M_{x} d \mu(x)
$$

over some standard measure space $(X, \mu)$, corresponding to the decomposition

$$
H \cong \int_{X}^{\oplus} H_{x} d \mu(x)
$$

of the underlying Hilbert space $H$ by a measurable field of Hilbert spaces $\left(H_{x}\right)_{x \in X}$, where $M_{x} \subseteq B\left(H_{x}\right)$ is a factor for $\mu$-almost all $x \in X$.

In this sense, factors are seen as the building blocks of von Neumann algebras. That theory, however, is technically extensive and we do not go into details here.

Lemma 9.9: Let $M \subseteq B(H)$ be a factor. Then, for any two non-zero projections $e, f \in M$, there are non-zero projections $e_{1}, f_{1} \in M$ such that $e_{1} \leq e, f_{1} \leq f$ and $e_{1} \sim f_{1}$.

Proof: (1) Claim: fMe $\neq\{0\}$.
Consider $K:=\operatorname{cl}(M e H) \subseteq H$, which is invariant under $M$. Thus, the projection $p$ onto $K$ is in $M^{\prime}$. Furthermore, we have that $K$ is invariant under $M^{\prime}$, i. e., for all $x \in M^{\prime}$ it holds $x K \subseteq K$ and thus $p x=x p$ (since also $x K^{\perp} \subseteq K^{\perp}$ ). Therefore, $p \in M^{\prime \prime}=M$. As $p \in M^{\prime}$ and $p \in M^{\prime \prime}=M$, we have $p \in M \cap M^{\prime}=Z(M)=\mathbb{C} 1$, it must hold $p=0$ or $p=1$. Since Kneq $\{0\}$ we infer $p=1$. Thus, $\operatorname{cl}(M e H)=K=H$. Now since $f \neq 0$, we conclude $f \operatorname{cl}(M e H) \neq\{0\}$, i. e., $f M e \neq\{0\}$.
(2) By (1), we find $z \in M$ such that $x:=f z e \neq 0$. We have: $f x=x$ and $x e=x$ (i. e., $e^{*} x=x^{*}$ ). From $f x=x$ we deduce $f \cdot \operatorname{supp}\left(x^{*}\right)=\operatorname{supp}\left(x^{*}\right)$ and from $x e=x$ we conclude $e \cdot \operatorname{supp}(x)=\operatorname{supp}(x)$. Thus, by Remark 9.3

$$
f_{1}:=\operatorname{supp}\left(x^{*}\right) \leq f \quad \text { and } \quad e_{1}:=\operatorname{supp}(x) \leq e
$$

and by Theorem 9.6, we also have that $e_{1} \sim f_{1}$.

Theorem 9.10: Let $M \subseteq B(H)$ be a factor. If $e, f \in M$ are projections, then $e \precsim f$ or $f \precsim e$ holds. So, any two projections in $\mathcal{P}(M)$ are comparable.

Proof: Consider families $\left(e_{i}, f_{i}\right)_{i \in I}$ such that

- $\left(e_{i}\right)_{i \in I}$ are mutually orthogonal projections in $M$,
- $\left(f_{i}\right)_{i \in I}$ are mutually orthogonal projections in $M$,
- For all $i \in I: e_{i} \leq e, f_{i} \leq f$ and $e_{i} \sim f_{i}$

They are partially ordered by inclusion; thus, by Zorn's Lemma, there is a maximal family $\left(e_{i} f_{i}\right)_{i \in I}$. Then:

$$
e^{\prime}:=\sum_{i \in I} e_{i} \leq e \quad \text { and } \quad f^{\prime}:=\sum_{i \in I} f_{i} \leq f
$$

and $e^{\prime} \sim f^{\prime}$ by Theorem 9.4.
Claim: $e^{\prime}=e$ or $f^{\prime}=f$.
Indeed, if both $e-e^{\prime} \neq 0$ and $f-f^{\prime} \neq 0$, we find by Lemma 9.9 non-zero projections $g \leq e-e^{\prime}$ and $h \leq f-f^{\prime}$ such that $g \sim h$. Hence, we can enlarge $\left(e_{i} f_{i}\right)_{i \in I}$ by $(g, h)$ which contradictis the maximality of $\left(e_{i} f_{i}\right)_{i \in I}$. Thus " $e=e^{\prime} \sim$ $f^{\prime} \leq f \Rightarrow e \precsim f^{\prime}$ " or " $f=f^{\prime} \sim e^{\prime} \leq e \Rightarrow f \precsim e$ ".

Definition 9.11: Let $M \subseteq B(H)$ be a von Neumann algebra.
(i) A projection $0 \neq e \in M$ is called minimal, if it holds " $(f \in \mathcal{P}(M), f \leq e) \Rightarrow$ ( $f=0$ or $f=e$ )".

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(ii) A projection $e \in M$ is called finite, of " $(f \in \mathcal{P}(M), e \sim f \leq e) \Rightarrow f=e$ ".

Remark 9.12: (i) Consider $M=B(H)$. Then $e$ is minimal if and only if $\operatorname{dim} e H=1$ and $e$ is finite if and only if $\operatorname{dim} e H<\infty$.
(ii) In general $e$ is minimal if and only if $e M e=\mathbb{C} e$ and if $e$ is minimal, then $e$ is finite - the converse is false however.

Indeed: If $e$ is minimal, then $e \sim f \leq e$ implies $f=0$ or $f=e$, where $f=0$ is not possible as $f \sim e \neq 0$.

Definition 9.13 (Murray-von-Neumann, ~1930): Let $M \subseteq B(H)$ be a factor. We say the $M$ is of
(i) type $I$, if $M$ contains a minimal projection.
(ii) type $I I$, if $M$ contains no minimal projection but a finite projection.
(iii) type III, if $M$ contains no finite projection.

## 10 Type I factors and tensor products

Factors of type I are isomorphic to $B(H)$ for some suitable Hilbert space $H$; at the same time, their commutants are also of this particular form. This result is the goal of this chapter, for which tensor products of von Neumann algebras are needed.

Remark 10.1: (i) Let $\left(H_{1},\langle\cdot, \cdot\rangle_{1}\right)$ and $\left(H_{2},\langle\cdot, \cdot\rangle_{2}\right)$ be (complex) Hilbert spaces. On their algebraic tensor product (over $\mathbb{C}$ )

$$
H_{1} \otimes H_{2}=\left\{\sum_{i=1}^{n} \xi_{i} \otimes \eta_{i}: n \in \mathbb{N}, \xi_{i} \in H, \eta_{i} \in H\right\}
$$

we may introduce an inner product $\langle\cdot, \cdot\rangle$, which is uniquely determined

$$
\left\langle\xi_{1} \otimes \eta_{1}, \xi_{2} \otimes \eta_{2}\right\rangle:=\left\langle\xi_{1}, \xi_{2}\right\rangle_{2}\left\langle\eta_{1}, \eta_{2}\right\rangle_{2}
$$

for all $\xi_{1}, \xi_{2} \in H_{1} ; \eta_{1}, \eta_{2} \in H_{2}$. This yields a pre-Hilbert space, whose completion will be denoted by $H_{1} \hat{\otimes} H_{2}$ and is called the Hilbert space tensor product of $H_{1}$ and $\mathrm{H}_{2}$.
(ii) If $\left(\xi_{i}\right)_{i \in I}$ and $\left(\eta_{j}\right)_{j \in J}$ are orthonormal bases of $H_{1}$ and $H_{2}$ respectively, then $\left(\xi_{i} \otimes \eta_{j}\right)_{(i, j) \in I \times J}$ is an orthonormal basis of $H_{1} \hat{\otimes} H_{2}$. Thus

$$
H_{1} \hat{\otimes} H_{2} \cong \bigoplus_{i \in I} \underset{\substack{\xi_{i} \otimes H_{2}}}{\xi_{\cong}}
$$

(iii) If $x \in B\left(H_{1}\right)$ and $y \in B\left(H_{2}\right)$ are given, then

$$
(x \otimes y)(\xi \otimes \eta):=(x \xi) \otimes(y \eta) \quad \text { for all } \xi \in H_{1}, \eta \in H_{2}
$$

defines a linear operator $x \otimes y: H_{1} \otimes H_{2} \rightarrow H_{1} \otimes H_{2}$ that extends uniquely to an operator $x \hat{\otimes} y \in B\left(H_{1} \hat{\otimes} H_{2}\right)$ with

$$
\|x \hat{\otimes} y\|=\|x\|\|y\| .
$$

(iv) Let $M_{1} \subseteq B\left(H_{1}\right)$ and $M_{2} \subseteq B\left(H_{2}\right)$ be von Neumann algebras. We define their (von Neumann algebra) tensor product by

$$
M_{1} \bar{\otimes} M_{2}=\operatorname{cl}_{\mathrm{WOT}}\left(\left\langle\left\{x \hat{\otimes} y \mid x \in M_{1}, y \in M_{2}\right\}\right\rangle\right) \subseteq B\left(H_{1} \hat{\otimes} H_{2}\right)
$$

It is a deep theorem (in fact, it was one of the first applications of the TomitaTakesaki theory) that

$$
\left(M_{1} \bar{\otimes} M_{2}\right)^{\prime}=M_{1}^{\prime} \bar{\otimes} M_{2}^{\prime} .
$$

Furthermore one can show that if both $M_{1}$ and $M_{2}$ are factors, then also $M_{1} \bar{\otimes} M_{2}$ is a factor.

Lemma 10.2: Let $\left(H_{1},\langle\cdot, \cdot\rangle_{1}\right)$ and $\left(H_{2},\langle\cdot, \cdot\rangle_{2}\right)$ be complex Hilbert spaces. Let $M \subseteq$ $B\left(H_{2}\right)$ be a von Neumann algebra. Then
(i) $\iota: B\left(H_{2}\right) \rightarrow B\left(H_{1} \hat{\otimes} H_{2}\right), y \mapsto \mathrm{id}_{H_{1}} \hat{\otimes} y$ is WOT-continuous and we have

$$
\iota(M)=\iota(M)^{\prime \prime}=\left(\operatorname{Cid}_{H_{1}}\right) \hat{\otimes} M
$$

i. e., $\operatorname{id}_{H_{1}} \otimes M=\iota(M) \subseteq B\left(H_{1} \hat{\otimes} H_{2}\right)$ is a von Neumann algebra.
(ii) We have that

$$
\left(B\left(H_{1}\right) \hat{\otimes} M\right)^{\prime}=\operatorname{id}_{H_{1}} \otimes M^{\prime} \quad \text { and } \quad\left(\operatorname{id}_{H_{1}} \otimes M\right)^{\prime}=B\left(H_{1}\right) \bar{\otimes} M^{\prime}
$$

In particular, since $B\left(H_{2}\right)^{\prime}=\operatorname{Cid}_{H_{2}}$,

$$
\left(\operatorname{id}_{H_{1}} \otimes B\left(H_{2}\right)\right)^{\prime}=B\left(H_{1}\right) \otimes \operatorname{id}_{H_{2}} \quad \text { and } \quad B\left(H_{1}\right) \bar{\otimes} B\left(H_{2}\right)=B\left(H_{1} \hat{\otimes} H_{2}\right)
$$

Proof: Fix an orthonormal basis $\left(\xi_{i}\right)_{i \in I}$ of $H_{1}$. We define

- $v_{i, j} \in B\left(H_{1}\right)$ for $i, j \in I$ via $v_{i, j} \xi:=\left\langle\xi, \xi_{j}\right\rangle_{1} \xi_{i}$
- $V_{i} \in B\left(H_{2}, H_{1} \hat{\otimes} B_{2}\right)$ for $i \in I$ via $V_{i, \eta}:=\xi_{i} \otimes \eta$.
(1) Claim: $V_{i} V_{j}^{*}=v_{i, j} \hat{\otimes} \mathrm{id}_{H_{2}}$ for all $i, j \in I$.

Proof (of (1)): First, we check that $V_{j}^{*}(\xi \otimes \eta)=\left\langle\xi, \xi_{j}\right\rangle \eta$. Indeed, for all $\tilde{\eta} \in H_{2}$, $\left\langle V_{j} \tilde{\eta}, \xi \otimes \eta\right\rangle=\left\langle\xi_{j} \otimes \tilde{\eta}, \xi \otimes \eta\right\rangle=\left\langle\xi_{j}, \xi\right\rangle_{1}\langle\tilde{\eta}, \eta\rangle_{2}=\left\langle\tilde{\eta},\left\langle\xi, \xi_{j}\right\rangle_{1} \eta\right\rangle_{2}=\left\langle\tilde{\eta}, V_{j}^{*}(\xi \otimes \eta)\right\rangle$.
Thus we get

$$
V_{i} V_{j}^{*}(\xi \otimes \eta)=V_{i}\left(\left\langle\xi, \xi_{j}\right\rangle_{1} \eta\right)=\left\langle\xi, \xi_{j}\right\rangle_{1} \xi_{1} \otimes \eta=v_{i, j} \xi \otimes \eta
$$

(2) Claim: $V_{i}^{*} V_{j}=\delta_{i, j} \operatorname{id}_{H_{2}}$.

Proof (of (2)): Indeed we have

$$
V_{i}^{*} V_{j} \eta=V_{i}^{*}\left(\xi_{j} \otimes \eta\right)=\left\langle\xi_{j}, \xi_{i}\right\rangle_{1} \eta=\delta_{i, j} \eta .
$$

(3) Claim: Let $A:=\left\{V_{i} \times V_{j}^{*} \mid i, j \in I, x \in M\right\}$, then $A^{\prime}=\operatorname{id}_{H_{1}} \otimes M^{\prime}$.

Proof (of (3)): Note that $V_{i} \times V_{j}^{*}=v_{i, j} \hat{\otimes} x$, thus " $\supseteq$ " is clear. Conversely, take $w \in A^{\prime}$. Then, in particular

$$
w V_{i} V_{j}^{*}=V_{i} V_{j}^{*} w
$$

for all $i, j \in I$. Assertion 2 implies that $y:=V_{i}^{*} w V_{i}=V_{j}^{*} w V_{j}$ for all $i, j \in I$ and $V_{i}^{*} w V_{j}=0$ for all $i, j \in I$ where $i \neq j$; hence

$$
w=\left(\sum_{i \in I} V_{i} V_{i}^{*}\right) w\left(\sum_{j \in I} V_{j} V_{j}^{*}\right)=\sum_{i \in I} V_{j} y V_{i}^{*}=\sum_{i \in I} v_{i, i} \hat{\otimes} y=\operatorname{id}_{H_{1}} \otimes y
$$

We want $y \in M^{\prime}$ : Take $x \in M$, then, for any $i \in I$,

$$
y x \stackrel{(2)}{=} V_{i}^{*} w\left(V_{i} \times V_{i}^{*}\right) V_{i}=V_{i}^{*}\left(V_{i} \times V_{i}^{*}\right) w V_{i}=x V_{i}^{*} w V_{i}=x y
$$

(4) Claim: $\left(B\left(H_{1}\right) \bar{\otimes} M\right)^{\prime}=\operatorname{id}_{H_{1}} \otimes M^{\prime}$.

Proof (of (4)): We clearly have $A \subseteq B\left(H_{1}\right) \bar{\otimes} M$; thus

$$
\left(B\left(H_{1} \bar{\otimes} M\right)^{\prime} \subseteq A^{\prime} \stackrel{(3)}{=} \operatorname{id}_{H_{1}} \otimes M^{\prime} \subseteq\left(B\left(H_{1}\right) \bar{\otimes} M\right)^{\prime}\right.
$$

and hence $\left(B\left(H_{1}\right) \bar{\otimes} M\right)^{\prime}=\operatorname{id}_{H_{1}} \otimes M^{\prime}$.
(5) Claim: $\left(\mathrm{id}_{H_{1}} \otimes M\right)^{\prime}=B\left(H_{1}\right) \bar{\otimes} M^{\prime}$.

Proof (of (5)): Applying (4) to $M^{\prime}$ instead of $M$, we get that

$$
\left(\mathrm{id}_{H_{1}} \otimes M\right)^{\prime}=\left(B\left(H_{1}\right) \bar{\otimes} M^{\prime}\right)^{\prime \prime}=B\left(H_{1}\right) \bar{\otimes} M^{\prime}
$$

Now, (4) and (5) prove (ii). For proving (i), note that by (ii) $\iota(M)=\operatorname{id}_{H_{1}} \otimes$ $M=\left(B\left(H_{1}\right) \bar{\otimes} M\right)^{\prime}$ so that $\iota(M)$ is weakly closed (see (Lemma 6.7) (i)). Thus $\iota(M)=\operatorname{cl}_{\text {WOT }}(M)=\iota(M)^{\prime \prime}$ and furthermore $\operatorname{cl}_{\text {WOT }}(M)=\left(\operatorname{Cid}_{H_{1}}\right) \bar{\otimes} M$.

Example 10.3: If $\left(H_{1},\langle\cdot, \cdot\rangle_{1}\right)$ and $\left(H_{2},\langle\cdot, \cdot\rangle_{2}\right)$ are complex Hilbert spaces, then $\operatorname{id}_{H_{1}} \otimes B\left(H_{2}\right) \subseteq B\left(H_{1} \hat{\otimes} H_{2}\right)$ is a factor of type I. Indeed:

$$
Z\left(\operatorname{id}_{H_{1}} \otimes B\left(H_{2}\right)\right)=\left(\operatorname{id}_{H_{1}} \otimes B\left(H_{1}\right)\right) \cap\left(\operatorname{id}_{H_{1}} \otimes B\left(H_{2}\right)\right)^{\prime}=\operatorname{Cid}_{H_{1} \hat{\otimes} H_{2}}
$$

and $\operatorname{id}_{H_{1}} \hat{\otimes} e$ for any projection $e \in B\left(H_{2}\right)$ with $\operatorname{dim} e H=1$ is minimal.
We prove now that in fact any factor of type I is of this particular form.
Theorem 10.4: Let $M \subseteq B(H)$ be a factor of type I. Then there are Hilbert spaces $H_{1}, H_{2}$ and a unitary $U: H \rightarrow H_{1} \hat{\otimes} H_{2}$ such that

$$
U M U^{*}=\operatorname{id}_{H_{1}} \otimes B\left(H_{2}\right) \quad \text { and } \quad U M^{\prime} U^{*}=B\left(H_{1}\right) \otimes \operatorname{id}_{H_{2}}
$$

Thus, we have a WOT-continuous *-isomorphism

$$
\Phi: B\left(H_{2}\right) \longrightarrow M, \quad x \longmapsto U^{*}\left(\mathrm{id}_{H_{1}} \otimes x\right) U
$$

Proof: Consider a minimal projection $e \in M$. Then, by Exercise 2 of Sheet 8, we find a family $\left(e_{i}\right)_{i \in I}$ of mutually orthogonal projections in $M$ such that

$$
e_{i} \sim e \forall i \in I \quad \text { and } \quad 1=r+\sum_{i \in I} e_{i}
$$

for some projection $r \in M$ with $r \precsim e$ but $r \nsim e$.
Claim: $r=0$. Take a partial isometry $u \in M$ such that $u^{*} u=r$ and $u u^{*} \leq e$. By minimality of $e$, either $u u^{*}=0$ or $u u^{*}=e$; the latter would give $r \sim e$, which is excluded, thus $u u^{*}=0$, which then gives $u=0$ and hence $r=0$.

Thus $1=\sum_{i \in I} e_{i}$, i. e.

$$
H=\bigoplus_{i \in I} e_{i} H
$$

## 10 Type I factors and tensor products

Take now partial isometries $u_{i} \in M$ such that $u_{i}^{*} u_{i}=e_{i}$ and $u_{i} u_{i}^{*}=e$ for all $i \in I$. Then

$$
U: H \longrightarrow e H \hat{\otimes} \ell^{2}(I) \quad\left(=H_{1} \hat{\otimes} H_{2}\right), \quad \xi \longmapsto \sum_{i \in I}\left(u_{i} \xi\right) \otimes \delta_{i}
$$

where $\left(\delta_{i}\right)_{i \in I}$ is an orthonormal basis of $\ell^{2}(I)$, is a surjective isometry, hence a unitary. Note that

$$
U^{*}(\eta \otimes \delta)=\sum_{j \in I}\left\langle\delta, \delta_{j}\right\rangle u_{j}^{*} \eta
$$

Now, take any $x \in M$ and put $\tilde{x}:=U x U^{*}$. Then

$$
\tilde{x}(\eta \otimes \delta)=U \sum_{j \in I}\left\langle\delta, \delta_{j}\right\rangle x u_{j}^{*} \eta=\sum_{i, j \in I}\left\langle\delta, \delta_{j}\right\rangle \underbrace{\left(u_{i} x u_{j}^{*} \eta\right)}_{=: x_{i, j} \in M} \otimes \delta_{i} .
$$

Now, since $e x_{i, j} e=u_{i} u_{i}^{*}\left(u_{i} x u_{j}^{*}\right) u_{j} u_{j}^{*}=u_{i} x u_{j}^{*}=x_{i, j}$ and $e M e=\mathbb{C} e$ by (Remark 9.12 ) (ii), it follows that $x_{i, j}=\lambda_{i, j} e$ for some $\lambda_{i, j} \in \mathbb{C}$. Thus, we get

$$
\begin{aligned}
\tilde{x}(\eta \otimes \delta) & =\sum_{i, j \in I}\left\langle\delta, \delta_{j}\right\rangle \lambda_{i, j} e \eta \otimes \delta_{i} \\
& =\sum_{i, j \in I}\left\langle\delta, \delta_{j}\right\rangle \lambda_{i, j} \eta \otimes \delta_{i} \quad \quad(\text { as } e \eta \in e H) \\
& =\eta \otimes \Lambda \delta
\end{aligned}
$$

where $\Lambda \in B\left(\ell^{2}(I)\right)$ is defined by $\Lambda \delta:=\sum_{i, j \in I}\left\langle\delta, \delta_{i}\right\rangle \lambda_{i, j} \delta_{i}$, (i. e., $\left.\Lambda \hat{=}\left(\lambda_{i, j}\right)_{i, j \in I}\right)$; finally

$$
U x U^{*}=\tilde{x}=\operatorname{id}_{e H} \hat{\otimes} \Lambda \in B\left(e H \hat{\otimes} \ell^{2}(I)\right)
$$

i. e., $U M U^{*} \subseteq \operatorname{id}_{e H} \otimes B\left(\ell^{2}(I)\right)$.

We also have the other inclusion " $\supseteq$ ", since each $\Lambda \in B\left(\ell^{2}(I)\right)$ induces an operator

$$
x=U^{*}\left(\operatorname{id}_{e H} \otimes \Lambda\right) U=\sum_{i, j \in I}\left\langle\Lambda \delta_{i}, \delta_{j}\right\rangle u_{j}^{*} u_{i} \in M
$$

(as $u_{j}^{*} u_{i} \in M$ ) such that $U x U^{*}=\operatorname{id}_{e H} \hat{\otimes} \Lambda$. Thus in total we get

$$
U M U^{*}=\operatorname{id}_{e H} \otimes B\left(\ell^{2}(I)\right)
$$

and (Lemma 10.2) (ii) yields that

$$
U M^{\prime} U^{*}=\left(U M U^{*}\right)^{\prime}=\left(\operatorname{id}_{e H} \otimes \ell^{2}(I)\right)^{\prime}=B(e H) \otimes \operatorname{id}_{\ell^{2}(I)}
$$

Definition 10.5: Let $M \subseteq B(H)$ be a factor of type I. We say that
(i) $M$ is of type $I_{n}$, if $M \cong B\left(H_{2}\right)$ where $\operatorname{dim} H_{2}=n$.
(ii) $M$ is of type $I_{\infty}$, if $M \cong B\left(H_{2}\right)$ where $\operatorname{dim} H_{2}=\infty$.

## 11 Group von Neumann algebras and type II factors

We present a general construction that produces von Neumann algebras starting from discrete groups. This will show, that there are indeed "interesting" examples of von Neumann algebras. In particular, we will prove existence of factors of type II.

Throughout the following, let $G$ be a discrete group. We associate with $G$ the Hilbert space

$$
\ell^{2}(G):=\left\{\xi: G \longrightarrow \mathbb{C}: \sum_{x \in G}|\xi(x)|^{2}<\infty\right\}
$$

with the inner product given by $\langle\xi, \eta\rangle:=\sum_{g \in G} \xi(x) \overline{\eta(x)}$.
Definition 11.1: (i) The left regular representation of $G$ is given by

$$
\lambda: G \longrightarrow B\left(\ell^{2}(G)\right), \quad g \longmapsto \lambda_{g}
$$

where $\left(\lambda_{g} \xi\right)(x):=\xi\left(g^{-1} x\right)$.
(ii) The right regular representation of $G$ is given by

$$
\rho: G \longrightarrow B\left(\ell^{2}(G)\right), \quad g \longmapsto \rho_{g}
$$

where $\left(\rho_{g} \xi\right)(x):=\xi(x g)$.
Remark 11.2: (i) For each $y \in G$, we define $\delta_{y} \in \ell^{2}(G)$ by

$$
\delta_{y}(x):= \begin{cases}1, & \text { if } x=y \\ 0, & \text { else }\end{cases}
$$

Then $\left(\delta_{y}\right)_{y \in G}$ forms an orthonormal basis of $\ell^{2}(G)$. Thus, we may write for any $\xi \in \ell^{2}(G)$ :

$$
\xi=\sum_{y \in G} \xi(y) \delta_{y}
$$

and $\lambda_{g}$ and $\rho_{g}$ are given by continuous and linear extension of $\lambda_{g} \delta_{y}=\delta_{g y}$ and $\rho_{g} \delta_{y}=\delta_{y g^{-1}}$. Since $\lambda_{g} \lambda_{g^{-1}}=1=\lambda_{g^{-1}} \lambda_{g}$ and $\lambda_{g}^{*}=\lambda_{g^{-1}}$, all $\lambda_{g}$ are unitaries on $\ell^{2}(G)$; the same holds for $\rho_{g}$.

Definition 11.3: (i) For $\xi, \eta \in \ell^{2}(G)$, we define their convolution by

$$
\begin{equation*}
\xi * \eta: G \longrightarrow \mathbb{C}, \quad(\xi * \eta)(x):=\sum_{g \in G} \xi(g) \eta\left(g^{-1} x\right)=\sum_{g \in G} \xi\left(x g^{-1}\right) \eta(g) \tag{11.1}
\end{equation*}
$$

In fact, $|(\xi * \eta)(x)| \leq\|\xi\|\|\eta\|$ for all $x \in G$.
(ii) Let $\xi \in \ell^{2}(G)$. We put $D_{\xi}:=\left\{\eta \in \ell^{2}(G) \mid \xi * \eta \in \ell^{2}(G)\right\}$ and define the unbounded linear operator

$$
L_{\xi}: D_{\xi} \longrightarrow \ell^{2}(G), \quad \eta \longmapsto \xi * \eta .
$$

Analogously, $D_{\xi}^{\prime}:=\left\{\eta \in \ell^{2}(G) \mid \eta * \xi \in \ell^{2}(G)\right\}$ and

$$
R_{\xi}: D_{\xi}^{\prime} \longrightarrow \ell^{2}(G), \quad \eta \longmapsto \eta * \xi
$$

Lemma 11.4: For $\xi \in \ell^{2}(G)$ both $L_{\xi}$ and $R_{\xi}$ have closed graph in $\ell^{2}(G) \oplus \ell^{2}(G)$. In particular: If $D_{\xi}=\ell^{2}(G)$, then $L_{\xi} \in B\left(\ell^{2}(G)\right)$, and analogously, if $D_{\xi}^{\prime}=\ell^{2}(G)$, then $R_{\xi} \in B\left(\ell^{2}(G)\right)$.

Proof: Let $\left(\eta_{n}\right)_{n \in \mathbb{N}}$ be a sequence in $\ell^{2}(G)$ such that $\xi_{n} \rightarrow \eta \in \ell^{2}(G)$ and that $L_{\xi} \eta_{n} \rightarrow \zeta \in \ell^{2}(G)$. Then, for each $x \in G$, we have

$$
\begin{aligned}
|\zeta(x)-(\xi * \eta)(x)| & =\lim _{n \rightarrow \infty}\left|\left(\xi * \eta_{n}\right)(x)-(\xi * \eta)(x)\right| \\
& =\lim _{n \rightarrow \infty}\left|\xi *\left(\eta_{n}-\eta\right)(x)\right| \leq \lim _{n \rightarrow \infty}\|\xi\|\left\|\eta_{n}-\eta\right\|=0
\end{aligned}
$$

thus we get $\xi * \eta=\zeta \in \ell^{2}(G)$, i. e., $\eta \in D_{\xi}$ and $L_{\xi} \eta=\zeta$. By the closed graph theorem (Theorem 4.16 from the Functional Analysis I lecture notes), we see that $L_{\xi} \in B\left(\ell^{2}(G)\right)$ if $D_{\xi}=\ell^{2}(G)$.

Definition 11.5: A vector $\xi \in \ell^{2}(G)$ is called

- left-convolver if $\xi * \ell^{2}(G) \subseteq \ell^{2}(G)$ (i.e., $D_{\xi}=\ell^{2}(G)$ ),
- right-convolver if $\ell^{2}(G) * \xi \subseteq \ell^{2}(G)$ (i. e., $D_{\xi}^{\prime}=\ell^{2}(G)$ ).

We define

$$
\begin{aligned}
& L(G):=\left\{L_{\xi} \mid \xi \in \ell^{2}(G) \text { is left-convolver }\right\} \subseteq B\left(\ell^{2}(G)\right), \\
& R(G):=\left\{R_{\xi} \mid \xi \in \ell^{2}(G) \text { is right-convolver }\right\} \subseteq B\left(\ell^{2}(G)\right) .
\end{aligned}
$$

Remark 11.6: (i) For each $g \in G, \delta_{g} \in \ell^{2}(G)$ is both a left- and right-convolver. We have (see Eq. (11.1)) $\delta_{g} * \xi=\lambda_{g} \xi$, i. e., $L_{\delta_{g}}=\lambda_{g}$ and analogously $\xi * \delta_{g}=\rho_{g^{-1}} \xi$, i. e., $R_{\delta_{g-1}}=\rho_{g}$; thus, $\lambda(G) \subseteq L(G)$ and $\rho(G) \subseteq R(G)$.
(ii) For $\xi \in \ell^{2}(G)$, we define $\bar{\xi} \in \ell^{2}(G)$ by $\bar{\xi}(x):=\overline{\xi\left(x^{-1}\right)}$. If $\xi \in \ell^{2}(G)$ is a left-convolver, then so is $\bar{\xi}$ and we have that $L_{\bar{\xi}}=L_{\xi}^{*}$. Similarly, $R_{\bar{\xi}}=R_{\xi}^{*}$ for right-convolvers $\xi$.
(iii) Since convolution is associative, we have $L_{\xi * \eta}=L_{\xi} L_{\eta}$ and $R_{\xi * \eta}=R_{\eta} R_{\xi}$.

Remark 11.6 shows that $L(G)$ and $R(G)$ are unital *-subalgebras of $B\left(\ell^{2}(G)\right)$ that contain $\lambda(G)$ and $\rho(G)$ respectively. We next show, that actually both are von Neumann algebras and commutants of each other.

Theorem 11.7: Let $G$ be a discrete group. Then both $L(G)$ and $R(G)$ are von Neumann algebras acting on $\ell^{2}(G)$ and it holds true that

$$
L(G)=R(G)^{\prime}=\rho(G)^{\prime} \quad \text { and } \quad R(G)=L(G)^{\prime}=\lambda(G)^{\prime}
$$

and therefore $\rho(G)^{\prime \prime}=R(G)$ and $\lambda(G)^{\prime \prime}=L(G)$.
We call $L(G)$ and $R(G)$ the left - and right group von Neumann algebra of $G$ repectively.

Proof: (1) Claim: $\rho(G)^{\prime} \subseteq L(G)$.
Proof (of (1)): Take $T \in \rho(G)^{\prime}$ and put $\xi:=T \delta_{e} \in \ell^{2}(G)$. Then, for all $g \in G$,

$$
\begin{align*}
\xi * \delta_{g} & =\rho_{g^{-1}} \xi  \tag{11.1}\\
& =\rho_{g^{-1}} T \delta_{e}=T \rho_{g^{-1}} T \delta_{e}=T \rho_{g^{-1}} \delta_{e}=T \delta_{g}
\end{align*}
$$

and hence, by linearity, we get

$$
\xi * \eta=T \eta
$$

for all $\xi \in\left\langle\left\{\delta_{g} \mid g \in G\right\}\right\rangle \subseteq \ell^{2}(G)$. If $\eta \in \ell^{2}(G)$ is given, we find a sequence $\left(\eta_{n}\right)_{n \in \mathbb{N}}$ in $\left\langle\left\{\delta_{g} \mid g \in G\right\}\right\rangle \subseteq D_{\xi}$ such that $\eta_{n} \rightarrow \eta$ and $L_{\xi} \eta_{n}=T \eta_{n} \rightarrow T \eta$ (due to the continuity of $T$ ). Thus, by Lemma 11.4, $\eta \in D_{\xi}$ and $L_{\xi} \eta=T \eta$. Therefore $D_{\xi}=\ell^{2}(G)$, i. e., $\xi$ is a left-convolver and $T=L_{\xi} \in L(G)$.
(2) Claim: $L(G)=R(G)^{\prime}=\rho(G)^{\prime}$.

Proof (of (2)): We clearly have $L(G) \subseteq R(G)^{\prime} \subseteq \rho(G)^{\prime}$ and by (1) we also have $\rho(G)^{\prime} \subseteq L(G)$. In particular, $L(G)$ is weakly closed (by (Lemma 6.7) (i)) and hence a von Neumann algebra.
(3) Similarly, $R(G)$ is a von Neumann algebra and $R(G)=L(G)^{\prime}=\lambda(G)^{\prime}$.
(4) Using (2) and (3) and von Neumanns bicommutant theorem (Theorem 6.8), we infer that

$$
\rho(G)^{\prime \prime} \stackrel{(2)}{=} R(G)^{\prime \prime} \stackrel{(3)}{6.8} R(G) \quad \text { and } \quad \lambda(G)^{\prime \prime} \stackrel{(3)}{=} L(G)^{\prime \prime} \stackrel{(2)}{6.8} L(G) \text {. }
$$

Theorem 11.8: Let $G$ be a discrete group. Then

$$
\tau: L(G) \longrightarrow \mathbb{C}, \quad x \longmapsto\left\langle x \delta_{e}, \delta_{e}\right\rangle
$$

is a positive linear functional with $\tau(\mathrm{id})=1$ (i.e., a state), which is WOTcontinuous and moreover

- faithful, i. e., " $\tau\left(x^{*} x\right)=0 \Rightarrow x=0 "$,
- tracial, i. e., $\tau(x y)=\tau(y x)$ for all $x, y \in L(G)$.

Proof: We only have to verify, that $\tau$ is faithful and tracial; the other assertions are obvious.
(1) Claim: $\tau$ is faithful.

Proof (of (1)): Take $x=L_{\xi} \in L(G)$ for some left-convolver $\xi \in \ell^{2}(G)$. Then:

$$
\|\xi\|^{2}=\left\|L_{\xi} \delta_{e}\right\|^{2}=\left\langle L_{\xi} \delta_{e}, L_{\xi} \delta_{e}\right\rangle=\left\langle L_{\xi}^{*} L_{\xi} \delta_{e}, \delta_{e}\right\rangle=\tau\left(x^{*} x\right)
$$

Thus, if $\tau\left(x^{*} x\right)=0$ implies $\xi=0$ and hence $x=L_{\xi}=0$.
(2) Claim: $\tau$ is tracial.

Proof (of (2)): By the WOT-continuity of $\tau$, it suffices to verify that $\tau(x y)=\tau(y x)$ for all $x, y$ in the strongly dense subalgebra $\left\langle\left\{\lambda_{g} \mid g \in G\right\}\right\rangle$ (due to Theorem 11.7). By the linearity, we only have to check that $\tau\left(\lambda_{g} \lambda_{h}\right)=\tau\left(\lambda_{h} \lambda_{g}\right)$ for all $g, h \in G$. This can be seen as follows:

$$
\begin{aligned}
\tau\left(\lambda_{g} \lambda_{h}\right)=\left\langle\lambda_{g} \lambda_{h} \delta_{e}, \delta_{e}\right\rangle=\left\langle\delta_{g h} \delta_{e}, \delta_{e}\right\rangle & = \begin{cases}1, & g h=e \\
0, & \text { else. }\end{cases} \\
& =\left\{\begin{array}{ll}
1, & h g=e, \\
0, & \text { else. }
\end{array}=\left\langle\lambda_{h} \lambda_{g} \delta_{e}, \delta_{l}\right\rangle=\tau\left(\lambda_{h} \lambda_{g}\right)\right.
\end{aligned}
$$

which concludes the proof of (2).
As $\tau$ is indeed tracial and faithful, we have proven Theorem 11.8.
Which (left) group von Neumann algebras are factors?
Definition 11.9: A group $G$ is called infinite conjugacy class group (abbreviated i.c.c. group), if each non-trivial (i. e., $h \neq e$ ) conjugacy class $\left\{g h g^{-1} \mid g \in G\right\}$ is infinite.

Theorem 11.10: Let $G$ be a discrete group. Then $L(G)$ is a factor if and only if $G$ is an infinite conjugacy class group.

Proof: " $\Rightarrow$ ": Suppose $G$ is not an infinite conjugacy class group, i. e., we find a non-trivial element $e \neq h \in G$ such that $h^{G}:=\left\{g h g^{-1} \mid g \in G\right\}$ is finite. We define $x:=\sum_{k \in h^{G}} \lambda_{k} \in L(G)$. Since $g h^{G} h^{-1}=h^{G}$ for all $g \in G$, we have

$$
\lambda_{g} x \lambda_{g^{-1}}=\sum_{k \in h^{G}} \lambda_{g k g^{-1}}=\sum_{k \in h^{G}} \lambda_{k}=x
$$

and thus $\lambda_{g} x=x \lambda_{g}$. By Theorem 11.7, we have that $x \in \lambda(G)^{\prime}=L(G)^{\prime}$; hence $x \in L(G) \cap L(G)^{\prime}=Z(L(G))$. Now, observe that $x \notin \mathbb{C} 1$ as $\left(\lambda_{g}\right)_{g \in G}$ are linearly
independent in $L(G)$ and $\lambda_{e}=1$ - Indeed, if $g_{1}, \ldots, g_{n} \in G$ and $\alpha_{1}, \ldots, \alpha_{n} \in \mathbb{C}$ with $\sum_{i=1}^{n} \alpha_{i} \lambda_{i}=0$ and $g_{1}, \ldots, g_{n}$ mutually different are given, then

$$
0=\left(\sum_{i=1}^{n} \alpha_{i} \lambda_{g_{i}}\right) \delta_{e}=\sum_{i=1}^{n} \alpha_{i} \delta_{g_{i}}
$$

and since $\left(\delta_{g}\right)_{g \in G}$ is an orthonormal basis of $\ell^{2}(G)$, we infer $\alpha_{1}=\cdots=\alpha_{n}=0$. Thus $Z(L(G)) \supsetneq \mathbb{C} 1$, i. e., $L(G)$ is not a factor.
" $\Leftarrow$ ": Suppose that $G$ is an infinite conjugacy class group. Take any $x \in Z(L(G))$; since $x \in L(G)$, we find a left-convolver $\xi \in \ell^{2}(G)$ such that $x=L_{\xi}$ and since also $x \in L(G)^{\prime}$, it satisfies for $h \in G$ :

$$
L_{\xi}=x=\lambda_{h^{-1}} x \lambda_{h}=\lambda_{h}^{*} x \lambda_{h}=L_{\delta_{h}-1 * \xi},
$$

and if applied to $\delta_{e}, \xi=\delta_{h^{-1}} * \xi * \delta_{h}$. Thus, for any $g \in G$ :

$$
\xi(g)=\left(\left(\delta_{h^{-1}} * \xi\right) * \delta_{h}\right)(g) \stackrel{\text { Eq. }}{=}\left(\text { (11.1) }_{=}\left(\delta_{h^{-1}} * \xi\right)\left(g h^{-1}\right) \stackrel{\text { Eq. }}{=}(11.1)\right.
$$

Thus, $\xi: G \rightarrow \mathbb{C}$ is constant on $\left\{h g h^{-1} \mid h \in G\right\}$ for each $g \in G$. Since $G$ is assumed to be an infinite conjugacy class group and $\xi \in \ell^{2}(G), \xi \equiv 0$ is enforced on $\bigcup_{g \in G \backslash\{e\}}\left\{h g h^{-1} \mid h \in G\right\}=G \backslash\{e\}$. Hence: $\xi=\xi(e) \delta_{e}$ and thus $x=\xi(e) 1 \in \mathbb{C} 1$.

What can be said about the type of such factors?
Remark 11.11: (i) A von Neumann algebra $M \subseteq B(H)$ is said to be finite if $1 \in M$ is a finite projection, or equivalently, if every isometry in $M$ is a unitary (i. e., $v^{*} v=1 \Rightarrow v v^{*}=1$ ). According to Exercise 2 (a), Sheet 7, each projection in a finite von Neumann algebra is finite.
(ii) Let $M \subseteq B(H)$ be a von Neumann algebra, that has a faithful tracial state $\tau: M \rightarrow \mathbb{C}$, which is moreover normal, i. e., WOT-continuous on

$$
M_{1}:=\{x \in M \mid\|x\| \leq 1\}
$$

Then $M$ is finite (In fact, the converse is also true, this requires work).
(iii) Let $M \subseteq B(H)$ be a factor of type I . Then $M$ is finite if and only if $M$ is of type $\mathrm{I}_{n}$ for some $n \in \mathbb{N}$, which follows from Theorem 10.4 and (Remark 9.12) (i). In particular, a factor of type I is finite if and only if it is finite dimensional.

Definition 11.12: Let $M \subseteq B(H)$ be a factor of type II. We say that
(i) $M$ is of type $I I_{1}$, if 1 is a finite projection.
(ii) $M$ is of type $I I_{\infty}$, if 1 is not a finite projection (but there are finite projections).

Remark 11.13: Let $M \subseteq B(H)$ be a factor. If $M$ has a faithful normal tracial state $\tau: M \rightarrow \mathbb{C}$ and is not finite dimensional, then $M$ is of type $\mathrm{II}_{1}$. Indeed: 1 is finite due to Remark 11.11 (ii) but not of type I due to Remark 11.11 (iii), hence of type $\mathrm{II}_{1}$.

Corollary 11.14: Let $G$ be a non-trivial discrete infinite conjugacy class group. Then $L(G)$ is a factor of type $I I_{1}$.

Proof: By Theorem 11.10, $L(G)$ is a factor. Due to Theorem 11.8, it admits a faithful normal tracial state $\tau: L(G) \rightarrow \mathbb{C}$. Since $G$ is non-trivial and an infinite conjugacy class group, it contains infinitely many elements; thus $L(G)$ is not finite dimensional. The assertion now follows from Remark 11.13.

Example 11.15: (i) Consider $G=S_{\infty}:=\bigcup_{n \in \mathbb{N}} S_{n}$, i. e., the group of all permutations of $\mathbb{N}$ that move only finitely many points. $S_{\infty}$ is an infinity conjugacy class group. Indeed, if $\sigma \in S_{\infty}, \sigma \neq \mathrm{id}$ is given, then we find $i \neq j$ such that $\sigma(i)=j$. Consider the transposition $\pi_{r}:=(i, r)$ for $r>\max \{i, j\}$, then we have

$$
\left(\pi_{r} \sigma \pi_{r}^{-1}\right)(r)=\left(\pi_{r} \sigma \pi_{r}\right)(r)=j
$$

i. e., all $\pi_{r} \sigma \pi_{r}^{-1}$ are different. Thus $L\left(S_{\infty}\right)$ is a type $\mathrm{II}_{1}$ factor. It is the so-called hyperfinite $I I_{1}$-factor $\mathcal{R}$.
(ii) Consider $G=\mathbf{F}_{n}$, the free group with $n$ generators; if $n \geq 2$, then $\mathbf{F}_{n}$ is an infinite conjugacy class group. Thus, $L\left(\mathbf{F}_{n}\right)$ is a type $\mathrm{II}_{1}$ factor; a so-called free group factor.

Murray and von Neumann have shown, that $L\left(\mathbf{F}_{n}\right) \not \not \mathcal{R}$; for that purpose they have introduced the so-called property $\Gamma$. It is still an open problem, whether $L\left(\mathbf{F}_{n}\right) \cong L\left(\mathbf{F}_{m}\right)$ for $n, m \geq 2, n \neq m$. This problem motivated Voiculesceu (arround 1985) to develop "free probability theory"; this has led to discover, in particular, amazing connections between such operator algebraic questions and random matrix theory. By using that bridge, Dykema and Radulescu (1994) have shown independently that

- either $L\left(\mathbf{F}_{n}\right) \cong L\left(\mathbf{F}_{m}\right)$ for all $n, m \geq 2$,
- or $L\left(\mathbf{F}_{n}\right) \neq L\left(\mathbf{F}_{m}\right)$ for all $n, m \geq 2, n \neq m$.

What other infinite conjugacy class groups $G$ satisfy $L(G) \cong \mathcal{R}$ ?
Remark 11.16: (i) A separable von Neumann algebra $M \subseteq B(H)$ is called hyperfinite, if there is an increasing sequence $\left(A_{n}\right)_{n \in \mathbb{N} \backslash\{0\}}$ of finite dimensional *-subalgebras of $M$ such that

$$
M=\operatorname{cl}_{\mathrm{WOT}}\left(\bigcup_{n \in \mathbb{N}} A_{n}\right) .
$$

For $M=L\left(S_{\infty}\right)$, one can choose $A_{n}:=L\left(S_{n}\right)$.
(ii) It was proven by Murray and von Neumann that up to isomorphism, there is a unique hyperfinite $\mathrm{II}_{1}$-factor; it is denoted by $\mathcal{R}$.
(iii) One can show that $\mathcal{R}$ embeds into any factor of type $\mathrm{II}_{1}$. On the other hand, the Connes-embedding conjecture problem asks whether every type $\mathrm{II}_{1}$-factor on a separable Hilbert space can be embedded into the ultrapower $\mathcal{R}^{\omega}$ of $\mathcal{R}$ by a free ultra-filter $\omega$; roughly speaking, this asks for matricial approximations.

Definition 11.17: A countable discrete group is called amenable, if there is a state $m: \ell^{\infty}(G) \rightarrow \mathbb{C}$ such that for all $f \in \ell^{\infty}(G)$ and $g \in G$ it holds

$$
m(f)=m(\lambda g f)
$$

where $(\lambda g f)(x):=f\left(g^{-1} x\right)$ for all $x \in G ; m$ is called a left-invariant mean on $G$.
Remark 11.18: (i) If $G$ is a finite group, then

$$
m(f):=\frac{1}{|G|} \sum_{g \in G} f(g)
$$

defines a left-invariant mean on $G$.
(ii) If $G$ can be written as $G=\bigcup_{n \in \mathbb{N}} G_{n}$ for an increasing sequence $\left(G_{n}\right)_{n \in \mathbb{N}}$ of subgroups of $G$, then it holds: If $G_{n}$ is amenable for all $n \in \mathbb{N}$, then so is $G$. This shows that $S_{\infty}=\bigcup_{n \in \mathbb{N}} S_{n}$ is amenable.

As a consequence of "Connes' tour de force" (1976) about injective type $\mathrm{II}_{1}$ factors, we have:

Theorem 11.19: Let $G$ be a countable non-trivial infinite conjugacy class group. Then $L(G) \cong \mathcal{R}$ if and only if $G$ is amenable.

## 12 The trace construction on finite factors

In this Chapter, we briefly discuss how on finite factors, i.e., factors that finite in the sense of (Remark 11.11) (i), a unique faithful normal tracial state can be constructed; see (Remark 11.11) (ii). The main results read as follows:

Theorem 12.1 (Existence of the trace): Let $M \subseteq B(H)$ be a factor. Then the following are equivalent:
(i) $M$ is finite.
(ii) $M$ has a (norm-continuous) tracial state $\tau: M \rightarrow \mathbb{C}$.
(iii) $M$ has a normal tracial state $\tau: M \rightarrow \mathbb{C}$.

Theorem 12.2 (Uniqueness of the trace): Let $M \subseteq B(H)$ be a finite factor. Then there is a unique (norm-continuous) tracial state $\tau: M \rightarrow \mathbb{C}$. Moreover, $\tau$ is automatically normal and faithful.

Every finite factor $M$ is either of type $\mathrm{I}_{n}$ for some $n \in \mathbb{N}$ or of type $\mathrm{II}_{1}$; if $M$ is of type $\mathrm{I}_{n}$, then $M \cong B\left(\mathbb{C}^{n}\right)=M_{n}(\mathbb{C})$. Then the trace $\tau$ on $M$ comes from

$$
\operatorname{tr}_{n}: M_{n}(\mathbb{C}) \longrightarrow \mathbb{C}, \quad\left(a_{i, j}\right)_{1 \leq i, j \leq n} \longmapsto \frac{1}{n} \sum_{i=1}^{n} a_{i, i}
$$

Thus, it suffices to treat the type $\mathrm{II}_{1}$ case.
Theorem 12.3 (Dimension function): Let $M \subseteq B(H)$ be a factor of type $I I_{1}$. Then there exists a function

$$
\Delta: \mathcal{P}(M) \longrightarrow[0,1]
$$

the so-called dimension function of $M$, with the following properties:
(i) $\Delta(1)=1$ and for all $p \in \mathcal{P}(M)$ it holds $\Delta(p)=0$ if and only if $p=0$.
(ii) For all $p, q \in \mathcal{P}(M)$ it holds $p \precsim q$ if and only if $\Delta(p) \leq \Delta(q)$. In particular it holds $p \sim q$ if and only if $\Delta(p)=\Delta(q)$.
(iii) For all $p, q \in \mathcal{P}(M), p q=0$, it holds $\Delta(p+q)=\Delta(p)+\Delta(q)$.
(iv) $\Delta$ is completely additive, i. e., for each family $\left(p_{i}\right)_{i \in I}$ of mutually orthogonal projections in $M$, it holds $\Delta\left(\sum_{i \in I} p_{i}\right)=\sum_{i \in I} \Delta\left(p_{i}\right)$.

Proof (Sketch): (1) "Halving lemma": For each $p \in \mathcal{P}(M)$, there are projections $p_{0}, p_{1} \in \mathcal{P}(M)$ such that $p_{0} \sim p_{1}$ and $p_{0}+p_{1}=p$; in other words, " $p_{0}$ and $p_{1}$ halve $p$ ".

Proof (of (1)): (a) Since $p$ is not minimal, we find $q \in \mathcal{P}(M)$ such that $q \leq p$, $q \neq 0$, and $q \neq p$. Since $M$ is a factor, Claim (1) in the proof of (Lemma 9.9) shows, that $q M(p-q) \neq\{0\}$. Take $x \in M$ with $y:=q x(p-q) \neq 0$ and consider its polar decomposition $y=u|y|$, then $u \neq 0$ and $u^{*} u+u u^{*} \leq p$; clearly $u^{*} u \perp u u^{*}$. Indeed, note that $u^{*} u$ is the projection onto $\operatorname{ker}(y)^{\perp} \subseteq(p-q) H \subseteq(1-q) H$ and $u u^{*}$ is the projection onto $\operatorname{cl}(\operatorname{im}(y)) \subseteq q H$.
(b) Consider families $\left(p_{i}, q_{i}\right)_{i \in I}$ such that

- $\left(p_{i}\right)_{i \in I}$ are mutually orthogonal projections in $M$,
- $\left(q_{i}\right)_{i \in I}$ are mutually orthogonal projections in $M$,
- For all $i \in I$ it holds $p_{i} \sim q_{i}$,
- $\sum_{i \in I} p_{i}+\sum_{i \in I} q_{i} \leq p$.

By Zorn's Lemma, there is a maximal family $\left(p_{i}, q_{i}\right)_{i \in I}$ of this kind; now put $p_{0}:=\sum_{i \in I} p_{i}$ and $q_{0}:=\sum_{i \in I} q_{i}$. Using (a), it follows due to maximality of $\left(p_{i}, q_{i}\right)_{i \in I}$ that $p_{0}+q_{0}=p$.

Note: The halving lemma is true in any diffuse factor; a von Neumann algebra is said to be diffuse, if it contains no minimal projections.
(2) "Fundamental projections": There is a sequence $\left(p_{n}\right)_{n \in \mathbb{N} \backslash\{0\}}$ of mutually orthogonal projections on $M$ such that for all $n \in \mathbb{N}$ it holds

$$
p_{n} \sim 1-\sum_{i=1}^{n} p_{i}
$$

One constructs the sequence $\left(p_{n}\right)_{n \in \mathbb{N}}$ by iterating (1). We have the following:
(a) If for $p \in \mathcal{P}(M)$ it holds $p \precsim p_{n}$ for all $n \in \mathbb{N}$, then $p=0$.
(b) $\sum_{n \in \mathbb{N}} p_{n}=1$ and therefore $p_{n} \sim \sum_{i \geq n+1} p_{i}$.

Indeed: $p:=1-\sum_{i \in \mathbb{N}} p_{i} \leq 1-\sum_{i=1}^{n} p_{i} \sim p_{n}$ for all $n \in \mathbb{N}$, by (a) now $p=0$ follows.
(c) For all $0 \neq p \in \mathcal{P}(M)$ there is $n \in \mathbb{N}$ such that $p_{n} \precsim p$.

Define

$$
\mathcal{F} \mathcal{P}(M):=\left\{p \in \mathcal{P}(M) \mid \exists n \in \mathbb{N}: p \sim p_{n}\right\}
$$

the set of fundamental projections in $M$
Fundamental projections (or, more precisely, equivalence classes thereof) play for projections the role of dyadic rationals for numbers in $[0,1]$. In fact, we have: For each $0 \neq p \in M$, there exists a unique increasing sequence $n_{1}<n_{2}<n_{3}<\ldots$ in $\mathbb{N}$ and a sequence $\left(p_{k}^{\prime}\right)_{k \in \mathbb{N}}$ of mutually orthogonal projections in $M$ such that
(a) $p_{k}^{\prime} \sim p_{k}$ for all $k \in \mathbb{N}$,
(b) $p=\sum_{k \in \mathbb{N}} p_{k}^{\prime}$.

We define

$$
\Delta(p):=\sum_{k \in \mathbb{N}} 2^{-n_{k}}
$$

i. e., $\Delta\left(p_{n}\right)=2^{-n}$. One can show - but this requires work - that this yields the desired dimension function.

Proof (of Theorem 12.1): (1) "Radon-Nikodym-trick": Let $\varphi, \psi: \mathcal{P}(M) \rightarrow[0, \infty)$ be completely additive maps with $\varphi, \psi \not \equiv 0$. Suppose that $\varphi$ is faithful, i. e., $\varphi(e) \neq 0$ whenever $e \neq 0$. Then
$\forall \varepsilon>0 \exists p \in \mathcal{F} \mathcal{P}(M), \theta>0 \forall p \in \mathcal{P}(M), q \leq p: \theta \varphi(q) \leq \psi(q) \leq \theta(1+\varepsilon) \varphi(q)$.
(2) Let $\psi: M \rightarrow \mathbb{C}$ be a positive linear functional and $\varepsilon>0$ such that

$$
\forall p \in \mathcal{P}(M): \Delta(q) \leq \psi(q) \leq(1+\varepsilon) \Delta(q) .
$$

Then for all positive $x \in M$ and unitaries $u \in M$,

$$
\psi\left(u x u^{*}\right) \leq(1+\varepsilon) \psi(x)
$$

i. e., $\psi$ is an $\varepsilon$-trace; equivalently, $\psi$ satisfies $\psi\left(x x^{*}\right) \leq(1+\varepsilon) \psi\left(x^{*} x\right)$ for all $x \in M$.
(3) Using (1) and (2), one can show that for every $\varepsilon>0$ a normal $\varepsilon$-trace $\psi_{\varepsilon}$ exists with

$$
\forall q \in \mathcal{P}(M):(1+\varepsilon)^{-1} \Delta(q) \leq \psi_{\varepsilon}(q) \leq(1+\varepsilon)^{2} \Delta(q)
$$

(4) Choose a decreasing sequence $\left(\varepsilon_{n}\right)_{n \in \mathbb{N}}$ in $(0, \infty)$ with $\varepsilon_{n} \downarrow 0$ as $n \rightarrow \infty$; then for all $p \in \mathcal{P}(M)$ it holds

$$
\lim _{n \rightarrow \infty} \psi_{\varepsilon_{n}}(q)=\Delta(q)
$$

and in fact, $\left(\psi_{\varepsilon_{n}}\right)_{n \in \mathbb{N}}$ converges in norm to a linear functional $\psi: M \rightarrow \mathbb{C}$ which turns out to be a faithful normal tracial state. This shows: "(i) $\Rightarrow$ (iii)" holds also for $M$ being of type $\mathrm{II}_{1}$.
"(iii) $\Rightarrow$ (ii)" is trivial and "(ii) $\Rightarrow$ (i)" is an exercise.

Proof (of Theorem 12.2): (1) "Dixmier's averaging theorem": Let $M \subseteq B(H)$ be a factor and let $x \in M$ be given. Then
(a) For all $\varepsilon>0$ there are unitaries $u_{1}, \ldots, u_{n} \in M$ and $\alpha \in \mathbb{C}$ such that

$$
\left\|\frac{1}{n} \sum_{i=1}^{n} u_{i} x u_{i}^{*}-\alpha 1\right\|<\varepsilon .
$$

(b) We put $K_{x}:=\operatorname{cl}_{\|\cdot\| \|}\left(\operatorname{conv}\left\{u x u^{*} \mid u \in M\right.\right.$ unitary $\left.\}\right)$; then $K_{x} \cap \mathbb{C} 1 \neq \varnothing$.

Note: (b) follows from (a) as $\varepsilon \downarrow 0$.
(2) If $\tau: M \rightarrow \mathbb{C}$ is a norm-continuous tracial state on the factor $M$, then $K_{x} \cap \mathbb{C} 1=\{\tau(x) 1\}$. Thus, there is at most one norm-continuous tracial state on $M$; hence the one found in Theorem 12.1 the unique one, which is even normal and faithful.

Remark 12.4: Consider $M=M_{n}(\mathbb{C})$. Then

$$
\operatorname{tr}_{n}(x) 1=\int_{U_{n}(\mathrm{C})} u x u^{*} d u
$$

where $U_{n}(\mathbb{C})=\left\{u \in M_{n}(\mathbb{C}) \mid u\right.$ unitary $\}$ and " $d u$ " stands for the Haar probability measure on $U_{n}(\mathbb{C})$. This fact is generalised by (1) (b) in the proof of Theorem 12.2.

## 13 The standard representation of tracial von Neumann algebras

Motivation 13.1: So far, we have studied von Neumann algebras (mostly) on their "own" Hilbert space that accompanies them by definition. However, especially for tracial von Neumann algebras, i.e., von Neumann algebras that are equipped with a faithful normal tracial state (which is unique in the case of a factor due to Theorem 12.2), there are other, sometimes better behaved, representations on other Hilbert spaces.

Definition 13.2: Let $M \subseteq B(H)$ be a von Neumann algebra and let $K$ be another complex Hilbert space.
(i) A unital *-homomorphism $\pi: M \rightarrow B(K)$ is called a representation of $M$ on $K$ (see Definition 5.12).
(ii) A representation $\pi: M \rightarrow B(K)$ is said to be normal, if it's restriction to $M_{1}=\{x \in M \mid\|x\| \leq 1\}$ is WOT-continuous.
(iii) A representation $\pi: M \rightarrow B(K)$ is said to be faithful, if it is injective.

Theorem 13.3: Let $\pi: M \rightarrow B(K)$ be a normal representation of $M$ on $K$. Then $\pi(M) \subseteq B(K)$ is a von Neumann algebra.

We do not give a proof but we point out that special cases appeared before in

- the proof of Theorem 8.15,
- in (Lemma 10.2) (i).

Theorem 13.4: Let $M \subseteq B(H)$ be a von Neumann algebra and let $\varphi: M \rightarrow \mathbb{C}$ be a state. Consider the cyclic representation $\left(\pi_{\varphi}, H_{\varphi}, \xi_{\varphi}\right)$ obtained by the GNS construction with respect to $\varphi$, see Theorem 5.15. Then $\varphi$ is normal (in the sense of (Remark 11.11) (ii)) if and only if $\pi_{\varphi}: M \rightarrow B\left(H_{\varphi}\right)$ is normal.

Proof: Since $\varphi(x)=\left\langle\pi_{\varphi}(x) \xi_{\varphi}, \xi_{\varphi}\right\rangle$ for all $x \in M, \pi_{\varphi}$ being normal enforces $\varphi$ to be normal.

Conversely, if $\varphi$ is normal, then

$$
x \longmapsto\left\langle\pi_{\varphi}(x) \hat{a}, \hat{b}\right\rangle_{\varphi}=\varphi\left(b^{*} x a\right)
$$

is WOT-continuous on $M_{1}$ for all $a, b \in M$, where $\hat{a}$ denotes the class of $a$ in $M / N_{\varphi} \subseteq H_{\varphi} ;$ note that $\xi_{\varphi}=\hat{1}$, so that $\hat{a}=\pi_{\varphi}(a) \xi_{p}$. Now, since $\xi_{\varphi}$ is cyclic, i. e., $\pi_{\varphi}(M) \xi_{\varphi}$ is dense in $H_{\varphi}$, we infer that $x \mapsto\left\langle\pi_{\varphi}(x), \eta\right\rangle_{\varphi}$ is WOT-continuous on $M_{1}$ for all $\xi, \eta \in H_{\varphi}$. Thus, $\pi_{\varphi}$ is normal.

Definition 13.5: Let $(M, \tau)$ be a tracial von Neumann algebra. We call the representation $\pi_{\tau}: M \rightarrow B\left(H_{\tau}\right)$ the standard representation of $(M, \tau)$ and we write

$$
L^{2}(M, \tau):=H_{\tau}\left(=\operatorname{cl}_{\langle\cdot, \cdot\rangle_{\tau}}(M) \text { as } N_{\tau}=\{0\}\right)
$$

The standard representation is normal (Theorem 13.4) and moreover faithful, since $\tau(x)=\left\langle\pi_{\tau}(x) \xi_{\tau}, \xi_{\tau}\right\rangle_{\tau}$ for all $x \in M$, thus we may identify $M$ with its image $\pi_{\tau}(M) \subseteq B\left(L^{2}(M, \tau)\right)$, in which case $M$ is said to be in standard form. Note that $\xi_{\tau}=\hat{1}$ is both cyclic and separating for $M$.

Remark 13.6: (i) A von Neumann algebra is said to be separable, if it has a faithful normal representation on a separable Hilbert space. This generalises (Definition 8.6), without changing the conclusions.
(ii) Let $(M, \tau)$ be a tracial von Neumann algebra. Then $\left(M_{1}, d\right)$ with

$$
d(x, y):=\|\hat{x}-\hat{y}\|_{\tau}
$$

is a complete metric space and the induced topology conincides with the strong operator topology from $B\left(L^{2}(M, \tau)\right)$. For a tracial von Neumann algebra ( $M, \tau$ ), the following are equivalent:

- $M$ is separable,
- $M_{1}$ is separable contains a SOT-dense sequence,
- $L^{2}(M, \tau)$ is a separable Hilbert space.

Example 13.7: Consider the tracial von Neumann algebra $\left(M_{n}(\mathbb{C}), \mathrm{tr}_{n}\right)$ acting on $\mathbb{C}^{n}$. It is in standard form on $L^{2}\left(M_{n}(\mathbb{C}), \operatorname{tr}_{n}\right)=\mathbb{C}^{n^{2}} \cong \mathbb{C}^{n} \otimes \mathbb{C}^{n}$ with

$$
\begin{aligned}
\pi_{\operatorname{tr}_{n}}\left(M_{n}(\mathbb{C})\right) & =M_{n}(\mathbb{C}) \otimes 1 \subseteq B\left(\mathbb{C}^{n} \otimes \mathbb{C}^{n}\right) \\
\text { and } \pi_{\operatorname{tr}_{n}}\left(M_{n}(\mathbb{C})\right)^{\prime} & =1 \otimes M_{n}(\mathbb{C}) \subseteq B\left(\mathbb{C}^{n} \otimes \mathbb{C}^{n}\right),
\end{aligned}
$$

where $\pi_{\operatorname{tr}_{n}}(x)=x \otimes 1$.
Definition 13.8: Consider a tracial von Neumann algebra $(M, \tau)$. The antilinear unitary involution $J: L^{2}(M, \tau) \rightarrow L^{2}(M, \tau)$ which is the extension to $L^{2}(M, \tau)$ of the antiunitary isometry

$$
J \hat{x}:=\widehat{x^{*}}
$$

for all $x \in M$, the canonical conjugation operator on $L^{2}(M, \tau)$. It satisfies $J^{2}=1$ and $\langle J \xi, \eta\rangle_{\tau}=\langle J \eta, \xi\rangle_{\tau}$ for all $\xi, \eta \in L^{2}(M, \tau)$.

Theorem 13.9: Let $(M, \tau)$ be a tracial von Neumann algebra in standard form on $L^{2}(M, \tau)$. Then $J M J=M^{\prime}$ on $B\left(L^{2}(M, \tau)\right)$ and $\tau_{M^{\prime}}(x):=\left\langle x \xi_{\tau}, \xi_{\tau}\right\rangle_{\tau}$ defines a faithful normal tracial state on $M^{\prime}$.

Remark 13.10: For $x \in M$ we have $L_{x}:=\pi_{\tau}(x) \in B\left(L^{2}(M, \tau)\right)$ with $L_{x} \hat{y}=\hat{x y}$ for all $y \in M$. Similarly, we have $R_{x} \in B\left(L^{2}(M, \tau)\right)$ with $R_{x} \hat{y}=\hat{y} x$ for all $y \in M$. Then $J L_{x} J=R_{x^{*}}$. Hence $M^{\prime}=J M J=\left\{R_{x} \mid x \in M\right\}$.

Definition 13.11: Let $M$ be a von Neumann algebra.
(i) A (left) $M$-module is a Hilbert space $H$ that comes with a normal representation $\pi: M \rightarrow B(H)$. We write $(H, \pi)$ or ${ }_{M} H$.
(ii) Let $\left(H, \pi_{H}\right)$ and $\left(K, \pi_{K}\right)$ be two $M$-modules. A (left) modular operator from $H$ to $K$ is an operator $T \in B(H, K)$ such that for all $x \in M$ it holds

$$
T \pi_{H}(x)=\pi_{K}(x) T
$$

In other words: The diagram

commutes. The space of all modular operators from $H$ to $K$ will be denoted by ${ }_{M} B(H, K)$.
(iii) We say that ${ }_{M} H$ and ${ }_{M} K$ are isomorphic (or equivalent), ${ }_{M} H \cong{ }_{M} K$, if there exists a unitary operator in ${ }_{M} B(H, K)$.

Example 13.12: Let $(M, \tau)$ be a tracial von Neumann algebra.
(i) The standard representation of $(M, \tau)$ gives a "canonical" $M$-module $L^{2}(M, \tau)$; see (Definition 13.5).
(ii) We can produce "larger" $M$-modules by amplifications: $H=L^{2}(M, \tau) \hat{\otimes} K$ for some (arbitrary) Hilbert space $K$; indeed

$$
\pi: M \longrightarrow B(H), \quad \pi(x)(\xi \otimes \eta):=(x \xi) \otimes \eta
$$

defines a normal representation. For example

$$
L^{2}(M, \tau) \hat{\otimes} \mathbb{C}^{n} \cong \bigoplus_{i=1}^{n} L^{2}(M, \tau)
$$

and $\pi(x)\left(\xi_{1}, \ldots, \xi_{n}\right):=\left(x \xi_{1}, \ldots, x \xi_{n}\right)$.
(iii) For an $M$-module $\left(H, \pi_{H}\right)$, we can construct a "smaller" $M$-module by "cutting down" with any projection $p \in{ }_{M} B(H)$, i. e., we take the Hilbert space $p H$ with the normal representation

$$
\pi_{p H}: M \longrightarrow B(p H),\left.\quad x \longmapsto \pi_{H}(x)\right|_{p H} \in B(p H)
$$

Note that ${ }_{M} B H=\pi_{H}(M)^{\prime}$.
By combining these constructions, we obtain in fact all separable $M$-modules of a separable type $\mathrm{II}_{1}$ factor $M$.

Theorem 13.13: Let $M$ be a separable factor of type $I_{1}$ in standard form on $L^{2}(M, \tau)$ for its unique faithful normal tracial state $\tau: M \rightarrow \mathbb{C}$. For every separable $M$-module ${ }_{M} H$, there exists an isometry $v \in{ }_{M} B\left(H, L^{2}(M, \tau) \otimes \ell^{2}(\mathbb{N})\right)$. In particular, $p:=v v^{*} \in{ }_{M} B\left(L^{2}(M, \tau) \hat{\otimes} \ell^{2}(\mathbb{N})\right)$ and

$$
{ }_{M} H \cong{ }_{M}\left(p\left(L^{2}(M, \tau) \hat{\otimes} \ell^{2}(\mathbb{N})\right)\right)
$$

This defines a bijection between the set of equivalence classes consisting of isomorphic separable $M$-modules and the set

$$
\mathcal{P}\left(\left(M \otimes \operatorname{id}_{\ell^{2}(\mathbb{N})}\right)^{\prime}\right) / \sim .
$$

Remark 13.14: We have to deal with the type $\mathrm{II}_{\infty}$ factor

$$
{ }_{M} B\left(L^{2}(M, \tau) \hat{\otimes} \ell^{2}(\mathbb{N})\right)=\left(M \otimes \mathrm{id}_{\ell^{2}(\mathbb{N}}\right)^{\prime}=M^{\prime} \bar{\otimes} B\left(\ell^{2}(\mathbb{N})\right)
$$

where $M^{\prime}$ is of type $\mathrm{II}_{1}$ and $B\left(\ell^{2}(\mathbb{N})\right)$ is of type $\mathrm{II}_{\infty}$. It carries a faithful normal semi-finite tracial weight

$$
\tau_{M^{\prime}} \bar{\otimes} \operatorname{Tr}:\left(M^{\prime} \bar{\otimes} B\left(\ell^{2}(\mathbb{N})\right)_{+} \longrightarrow[0, \infty]\right.
$$

where $\operatorname{Tr}: B\left(\ell^{2}(\mathbb{N})\right)_{+} \rightarrow[0, \infty]$ is defined by $\operatorname{Tr}(x):=\sum_{n \in \mathbb{N}}\left\langle x e_{n}, e_{n}\right\rangle$ for any orthonormal basis $\left(e_{n}\right)_{n \in \mathbb{N}}$ of $\ell^{2}(\mathbb{N})$; see Exercise 4 of Sheet 5 , where Hilbert$S c h m i d t$ operators were introduced, i. e., $x \in B\left(\ell^{2}(\mathbb{N})\right)$ with

$$
\|x\|_{2}=\operatorname{Tr}\left(x^{*} x\right)^{\frac{1}{2}}<\infty .
$$

Definition 13.15: Let $M$ be a von Neumann algebra and define the set of positive elements $M_{+}:=\{x \in M \mid x$ positive $\}$ of $M$. A map $\operatorname{Tr}: M_{+} \rightarrow[0, \infty]$ is called a tracial weight, if

- $\operatorname{Tr}(x+y)=\operatorname{Tr}(x)+\operatorname{Tr}(y)$ for all $x, y \in M_{*}$,
- $\operatorname{Tr}(\lambda x)=\lambda \operatorname{Tr}(x)$ for all $x \in M_{+}, \lambda \geq 0,{ }^{1}$
- $\operatorname{Tr}\left(x^{*} x\right)=\operatorname{Tr}\left(x x^{*}\right)$ for all $x \in M$.

A tracial weight $\operatorname{Tr}: M_{+} \rightarrow[0, \infty]$ is called
(i) semi-finite, if for all $0 \neq x \in M_{+}$there is $0 \neq y \in M_{+}$such that $y \leq x$ and $\operatorname{Tr}(y)<\infty$ hold,
(ii) normal, if for every increasing net $\left(x_{\lambda}\right)_{\lambda \in \Lambda}$ in $M^{+}$converging strongly to $x \in M_{+}$it holds $\operatorname{Tr}\left(x_{\lambda}\right) \rightarrow \operatorname{Tr}(x)$,
(iii) faithful, if for $x \in M_{+}$it holds: $" \operatorname{Tr}(x)=0 \Leftrightarrow x=0$ ".

[^9]Remark 13.16: (i) Every type $\mathrm{II}_{\infty}$ factor $M$ is of the form $N \bar{\otimes} B\left(\ell^{2}(I)\right)$ for some type $\mathrm{II}_{1}$ factor $N$ and an infinite set $I$ and thus admits a faithful normal semi-finite tracial weight $\operatorname{Tr}: M_{+} \rightarrow[0, \infty]$, which is unique up to scaling with $\lambda>0$.
(ii) In the situation of (i), we have

- $\operatorname{Tr}(\mathcal{P}(M))=[0, \infty]$,
- For all $p, q \in \mathcal{P}(M)$ it holds: " $p \precsim q$ if and only if $\operatorname{Tr}(p) \leq \operatorname{Tr}(q)$ ".

Definition 13.17: Let $M$ be a separable factor of type $\mathrm{II}_{1}$. For every separable $M$-module $(H, \pi)$, we define its $M$-dimension (or coupling constant) by

$$
\operatorname{dim}_{M} H:=\left(\tau_{M^{\prime}} \bar{\otimes} \operatorname{Tr}\right)\left(v v^{*}\right) \in[0, \infty]
$$

where $v$ is chosen like in (Theorem 13.13).
Remark 13.18: Combining (Theorem 13.13), (Remark 13.14) and (Remark 13.16) (ii), we conclude that $(H, \pi) \mapsto \operatorname{dim}_{M} H$ gives a bijection between classes of equivalent separable $M$-modules and the set $[0, \infty]$.
Theorem 13.19: In the situation of (Definition 13.17), we have that
(i) $\operatorname{dim}_{M} L^{2}(M, \tau)=1$,
(ii) $\operatorname{dim}_{M}\left(\bigoplus_{i \in I} H_{i}\right)=\sum_{i \in I} \operatorname{dim}_{M} H_{i}$, if $I$ is countable,
(iii) $\operatorname{dim}_{M} p H=\tau_{\pi(M)^{\prime}}(p) \operatorname{dim}_{M} H$ if $\pi(M)^{\prime} \subseteq B(H)$ is a type $I I_{1}$ factor.

Definition 13.20: Let $N \subseteq M$ be factors of type $\mathrm{I}_{1}$ (with $1_{M} \in N$, i. e., subfactors). We call

$$
[M: N]:=\operatorname{dim}_{N} L^{2}\left(M, \tau_{M}\right)
$$

the Jones index of $N$ in $M$ (note that $L^{2}\left(M, \tau_{M}\right)$ is a $M$-module and thus, in particular, a $N$-module).

Remark 13.21: Since $L^{2}\left(N, \tau_{N}\right) \subseteq L^{2}\left(M, \tau_{M}\right)$ as $\left.\tau_{M}\right|_{N}=\tau_{N}$, we have the decomposition $L^{2}\left(M, \tau_{M}\right)=L^{2}\left(N, \tau_{N}\right) \oplus L^{2}\left(N, \tau_{N}\right)^{\perp}$. Thus:

$$
\operatorname{dim}_{N} L^{2}\left(M, \tau_{M}\right)=\operatorname{dim}_{N} L^{2}\left(N, \tau_{N}\right)+\operatorname{dim}_{N} L^{2}\left(N, \tau_{N}\right)^{\perp} \geq 1+0=1
$$

i. e., $[M: N] \geq 1$ and it holds $[M: N]=1$ if and only if $M=N$.

Example 13.22: (i) Let $N$ be a type $\mathrm{II}_{1}$ factor. $N \cong 1 \otimes N \subseteq M_{k}(\mathbb{C}) \bar{\otimes} N=M$, then $[M: N]=k^{2}$.
(ii) Let $H \subseteq G$ be non-trivial countable discrete infinite conjugacy class groups. Then $[L(G): L(H)]=[G: H]$.

Theorem 13.23 (Jones, 1983): Let $N \subseteq M$ be factors of type $I I_{1}$ with $[M: N]<\infty$. Then

$$
[M: N] \in[4, \infty) \cup\left\{4 \cos ^{2}\left(\frac{\pi}{n+2}\right): n \in \mathbb{N}\right\}
$$

and all these values show up as indices of subfactors of the hyperfinite $I I_{1}$ factor $\mathcal{R}$.

## 14 Type III factors

At the time of von Neumann, factors of type III were more or less an enigma. This has changed totally with the work of Tomita (1967) and Takesaki (1970) on the modular theory for von Neumann algebras and the subsequent classification of type III factors by Connes (1973). We present briefly some of these pearls.

Let $M \subseteq B(H)$ be a von Neumann algebra and suppose that there exists a cyclic and separating vector $\xi \in H$. Consider the faithful normal state

$$
\varphi: M \longrightarrow \mathbb{C}, \quad x \longmapsto\langle x \xi, \xi\rangle .
$$

Definition 14.1: We define unbounded antilinear operators on $H$ by

$$
\begin{aligned}
S_{0}: H \supseteq D\left(S_{0}\right): & :=M \xi \longrightarrow H, & & x \xi \longmapsto x^{*} \xi \\
\text { and } F_{0}: H \supseteq D\left(F_{0}\right): & =M^{\prime} \xi \longrightarrow H, & & y \xi \longmapsto y^{*} \xi .
\end{aligned}
$$

They are well-defined since $\xi$ is separating for $M$ and $M^{\prime}$ (see Theorem 8.4).
Lemma 14.2: The operators $S_{0}$ and $F_{0}$ from Definition 14.1 are densely defined. We have $F_{0} \subseteq S_{0}^{*}$ and $S_{0} \subseteq F_{0}^{*}$, so that $S_{0}$ and $F_{0}$ are closable.

Proof: Since $\xi$ is cyclic for $M$, we see that $S_{0}$ is densely defined. Since $\xi$ is also separating for $M=\left(M^{\prime}\right)^{\prime}$, it is cyclic for $M^{\prime}$ by (Theorem 8.4), thus $F_{0}$ is densely defined. Take now $x \in M$ and $y \in M^{\prime}$. We have then:

$$
\left\langle S_{0}(x \xi), y \xi\right\rangle=\left\langle x^{*} \xi, y \xi\right\rangle=\langle\xi, x y \xi\rangle=\langle\xi, y x \xi\rangle=\left\langle y^{*} \xi, x \xi\right\rangle=\left\langle F_{0}(y \xi), x \xi\right\rangle
$$

which shows $S_{0}^{*} \supseteq F_{0}$ and $F_{0}^{*} \supseteq S_{0}$. We define:

- $S$ to be the closure of $S_{0}$,
- $F:=S^{*}=S_{0}^{*}$.
(One can show that $F$ is the closure of $F_{0}$, which gives a perfect symmetry in $M$ and $M^{\prime}$; this, however, is not necessary for the development of the theory).

Since $S_{0}=S_{0}^{-1}$, it follows that $S$ and $F$ are injective with dense range and $S=S^{-1}, F=F^{-1}$.

Definition 14.3: Put $\Delta:=S^{*} S=F S$, which is densely defined, positive and invertible with $\Delta^{-1}=S F$. If $S=J \Delta^{1 / 2}$ is the polar decomposition (!), then the operator $J=J_{\varphi}: H \rightarrow H$ is an invertible antilinear isometry satisfying $J=J^{-1}$, $J \Delta J=\Delta^{-1}$ and $F=J \Delta^{-1 / 2}$. We call $\Delta=\Delta_{\varphi}$ the modular operator for $(M, \varphi)$.

This generalises (Definition 13.8) to the non-tracial setting; the modular operator quantifies the failure of $\varphi$ being tracial and the resulting difference between $S$ and $J$.

Example 14.4: Consider the Hilbert space $H=M_{n}(\mathbb{C})$ with the inner product $\langle x, y\rangle=\operatorname{tr}_{n}\left(y^{*} x\right)$ and let $M=M_{n}(\mathbb{C}) \subseteq B(H)$ act by left-multiplication. Every positive linear functional $\varphi: M \rightarrow \mathbb{C}$ is of the form $\varphi(x)=\operatorname{tr}_{n}\left(h_{\varphi} x\right)$ with a unique density matrix $h_{\varphi} \geq 0$; if $\varphi$ is faithful, $h_{\varphi}$ is invertible and $\xi=h_{\varphi}^{1 / 2} \in H$ is a cyclic and separating vector with

$$
\langle x \xi, \xi\rangle=\left\langle x h_{\varphi}^{1 / 2}, h_{\varphi}^{1 / 2}\right\rangle=\operatorname{tr}_{n}\left(h_{\varphi} x\right)=\varphi(x)
$$

One can show that $\Delta_{\varphi} x=h_{\varphi} x h_{\varphi}^{-1}$ and $J_{\varphi} x=x^{*}$.
We even have an analogue of (Theorem 13.9):
Theorem 14.5 (Tomita-Takesaki theorem): In the situation described above, we have that

$$
J_{\varphi} M J_{\varphi}=M^{\prime} \quad \text { and } \quad \Delta_{\varphi}^{\mathrm{i} t} M \Delta_{\varphi}^{-\mathrm{i} t}=M \forall t \in \mathbb{R} .
$$

Note that for every $\alpha \in \mathbb{C}$, a closed operator $\Delta_{\varphi}^{\alpha}$ can be defined by functional calculus. In particular, $t \mapsto \Delta_{\varphi}^{\mathrm{i} t}$ is a strongly continuous one-parameter group of unitaries in $B(H)$.

Definition 14.6: The strongly continuous one-parameter group

$$
t \longmapsto \sigma_{t}^{\varphi}:=\left.\operatorname{Ad} \Delta_{\varphi}^{\mathrm{i} t}\right|_{M}, \quad \sigma_{t}^{\varphi}(x)=\Delta_{\varphi}^{\mathrm{i} t} x \Delta_{\varphi}^{-\mathrm{i} t}
$$

of automorphisms of $M$ is called the modular automorphism group of $(M, \varphi)$.
Proposition 14.7: The modular automorphism group $\left\{\sigma_{t}^{\varphi} \mid t \in \mathbb{R}\right\}$ of $(M, \varphi)$ satisfies the Kubo-Martin-Schwinger boundary condition with respect to $\varphi$ (abbreviated KMS boundary condition), i. e.,
(i) $\varphi=\varphi \circ \sigma_{t}^{\varphi}$ for all $t \in \mathbb{R}$,
(ii) For all $x, y \in M$, there is a bounded continuous function $F: \operatorname{cl}(\mathbb{S}) \rightarrow \mathbb{C}$ for $\mathbb{S}:=\{z \in \mathbb{C} \mid 0<\operatorname{Im}(z)<1\}$ which is holomorphic on $\mathbb{S}$ such that for all $t \in \mathbb{R}$ it holds:

$$
F(t)=\varphi\left(\sigma_{t}^{\varphi}(x) y\right) \quad \text { and } \quad F(t+\mathrm{i})=\varphi\left(y \sigma_{t}^{\varphi}(x)\right)
$$

(In fact, this determines $\sigma^{\varphi}$ uniquely).
Definition 14.8: Let $M$ be a separable factor. We put

$$
S(M):=\bigcap\left\{\operatorname{Sp}\left(\Delta_{\varphi}\right) \mid \varphi: M \rightarrow \mathbb{C} \text { is a faithful normal state }\right\} .
$$

For $\operatorname{Sp}\left(\Delta_{\varphi}\right)$ consider $M \cong \pi_{\varphi}(M) \subseteq B\left(H_{\varphi}\right)$ for the cyclic representation $\left(H_{\varphi}, \pi_{\varphi}, \xi_{\varphi}\right)$; see (Theorem 13.4).

Theorem 14.9 (Connes, 1973): Let $M$ be a separable factor. Then:
(i) If $0 \in S(M)$, then $M$ is of type III,
(ii) If $M$ is of type III, then $S(M) \backslash\{0\}$ is a closed subgroup of $\mathbb{R}_{+}=(0, \infty)$; thus, precisely one of the following cases occurs:

- $M$ is of type $I I_{0}$, i. e., $S(M)=\{0,1\}$,
- $M$ is of type $I I_{\lambda}$, i.e., $S(M)=\{0\} \cup\left\{\lambda^{n} \mid n \in \mathbb{Z}\right\}$ for some $0<\lambda<1$,
- $M$ is of type $I I I_{1}$, i. e., $S(M)=[0, \infty)$.

Example 14.10: Take $0<\lambda<1$. We define $\alpha:=\lambda /(1+\lambda) \in\left(0, \frac{1}{2}\right)$ and

$$
\varphi_{\lambda}: M_{2}(\mathbb{C}) \longrightarrow \mathbb{C}, \quad \varphi_{\lambda}\left(\left(\begin{array}{ll}
a_{1,1} & a_{1,2} \\
a_{2,1} & a_{2,2}
\end{array}\right)\right)=\alpha a_{1,1}+(1-\alpha) a_{2,2}
$$

The Powers factor (Powers, 1967) given by

$$
\mathcal{R}_{\lambda}:=\overline{\bigotimes_{i \in \mathbb{N}}}\left(M_{2}(\mathbb{C}), \varphi_{\lambda}\right)
$$

as an infinite tensor product, is of type $\mathrm{III}_{\lambda}$. This construction is analogous to Exercise 2 of Sheet 10.

## 15 Universal $C^{*}$-algebras

A $C^{*}$-algebra is a Banach algebra endowed with an involution and a norm, that satisfies the $C^{*}$-condition. This axiomatic definition allows for a universal construction:

Construction 15.1: Let $E=\left\{x_{i} \mid i \in I\right\}$ be a set of generators, $I$ an index set. Let $P(E)$ be the involutive $\mathbb{C}$-algebra of non-commutative polynomials in $E \cup E^{*}$, where $E^{*}:=\left\{x_{i}^{*} \mid i \in I\right\},\left(\lambda x_{i_{1}} \cdots x_{i_{k}}\right)^{*}:=\bar{\lambda} x_{i_{k}}^{*} \cdots x_{i_{1}}^{*}$. Hence

$$
P(E)=\left\langle\left\{x_{i_{1}}^{\alpha_{1}} \cdots x_{i_{k}}^{\alpha_{k}} \mid \alpha_{i} \in\left\{1,{ }^{*}\right\}\right\}\right\rangle .
$$

Let $R \subseteq P(E)$ be a set of relations. Let $J(R) \subseteq P(E)$ be the two-sided ideal in $P(E)$ generated by $R$. Put

$$
A(E, R):=P(E) / J(R)
$$

the universal involutive algebra with generators $E$ and relations $R$. For $x \in A(E, R)$, put

$$
\|x\|:=\sup \left\{p(x) \mid p \text { is a } C^{*} \text {-seminorm on } A(E, R)\right\}
$$

Here $p$ is a $C^{*}$-seminorm if and only if $p(\lambda x)=|\lambda| p(x), p(x+y) \leq p(x)+p(y)$, $p(x y) \leq p(x) p(y)$ and $p\left(x^{*} x\right)=p(x)^{2}$ for all $x, y \in A(E, R), \lambda \in \mathbb{C}$. If now $\|x\|<\infty$ for all $x \in A(E, R)$, put

$$
C^{*}(E \mid R):=\operatorname{cl}_{\|\cdot\|}(A(E, R) /\{x \in A(E, R) \mid\|x\|=0\}),
$$

the universal $C^{*}$-algebra with generators $E$ and relations $R$.
Remark 15.2: (i) The principle used in Construction 15.1 is more or less the same principle as in " $\{$ non-commutative polynomials in $x$ and 1$\}=C^{*}(x, 1) \subseteq A$ ".
(ii) "Relations" really means: If $x_{1}^{*} x_{1}-1 \in R$, then $x_{1}^{*} x_{1}=1$ in $C^{*}(E \mid R)$.
(iii) If $A$ is any given ${ }^{*}$-algebra, $C^{*}(A):=\mathrm{cl}_{\|\cdot\|}(A /\{x \mid\|x\|=0\})$, where $\|x\|:=\sup \left\{p(x) \mid p\right.$ is $C^{*}$-seminorm on $\left.A\right\}$; the so called enveloping $C^{*}$-algebra of $A$, if $\|x\|<\infty$ for all $x \in A(E, R)$.

Proposition 15.3: The universial $C^{*}$-algebras have the following universal property: Let $B$ be a $C^{*}$-algebra, $E^{\prime}:=\left\{y_{i} \in B \mid i \in I\right\} \subseteq B$ be a subset satisfying the relations $R$. Then there is a unique *-homomorphism $\varphi: C^{*}(E \mid R) \rightarrow B$ with $\varphi\left(x_{i}\right)=y_{i}$.

Proof: Consider the following diagram:


There is a so called replacement homomorphism $\varphi_{0}$, since $\varphi_{0}(R)=0$ and this replacement homomorphism is continuous (by the definition of $\|x\| ; p(x):=$ $\left\|\varphi_{0}(x)\right\|_{B}$ defines a $C^{*}$-seminorm and it holds $\left.\left\|\varphi_{0}(x)\right\| \leq\|x\|\right)$, thus extends to a *-homomorphism $\varphi: C^{*}(E \mid R) \rightarrow B$.

Lemma 15.4: If there is a constant $C>0$ such that $p\left(x_{i}\right)<C$ for all $C^{*}$-seminorms $p$ on $A(E, R)$ and all $i \in I$, then $C^{*}(E \mid R)$ exists (i. e., it holds $\|x\|<\infty$ for all $x \in A(E, R)$ ).

Proof: Let $x=x_{i_{1}}^{\alpha_{1}} \cdots x_{i_{k}}^{\alpha_{k}}\left(\alpha_{i} \in\left\{1,{ }^{*}\right\}\right)$ be a monomial and $p(x) \leq C^{k}$. Then this also holds for polynomials.

Proposition 15.5: Let $\mathbb{S}^{1}:=\{\lambda \in \mathbb{C}| | \lambda \mid=1\}$. Then

$$
C\left(\mathbb{S}^{1}\right) \cong C^{*}\left(u, 1 \mid u^{*} u=u u^{*}=1\right)
$$

Hence, $C\left(\mathbb{S}^{1}\right)$ is the universal $C^{*}$-algebra generated by a unitary.
Proof: $C^{*}(u, 1)=C^{*}(u, 1 \mid \ldots)$ has $E=\{1, u\}$ and the relations $u^{*} u-1, u u^{*}-1$, $1 u-u, u 1-u, 11-1$, etc. It exists, since $p(1)^{2}=p\left(1^{*} 1\right)=p(1) \leq 1$ and $p(u)^{2}=p\left(u^{*} u\right)=p(1) \leq 1$.

The identity function defined by $z(t)=t$ for $t \in \mathbb{S}^{1}$ and $1(t) \equiv 1$ for $t \in \mathbb{S}^{1}$ satisfy $z^{*} z=z z^{*}=1$. Hence, by the universal property, there is *-homomorphism

$$
\begin{aligned}
\varphi_{0}: C^{*}(u, 1) & \longrightarrow C\left(\mathbb{S}^{1}\right), \\
u & \longmapsto z, \\
1 & \longmapsto 1 .
\end{aligned}
$$

On the other hand $C^{*}(u, 1) \cong C(\operatorname{Sp}(u))$ by the theorem of Gelfand-Naimark (or rather the functional calculus) and $\operatorname{Sp}(u) \subseteq \mathbb{S}^{1}$. Hence we have the diagram


It holds by functional calculus

$$
\varphi \circ \varphi(u)=\psi(z)=u
$$

hence $\varphi$ is injective; in total we conclude that $\varphi$ is an isomorphism.
The intuition is: A "universal unitary $u$ " should have the maximal spectrum, so $\operatorname{Sp}(u)=\mathbb{S}^{1}$. Hence $C^{*}(u, 1) \cong C(\operatorname{Sp}(u))=C\left(\mathbb{S}^{1}\right)$.

Example 15.6: The universal $C^{*}$-algebra $C^{*}\left(x \mid x=x^{*}\right)$ does not exist! There are elements $y=y^{*} \in B$ with arbitrarily large norm, hence $p(x)=\|y\|$ can be large, hence $\|x\|=\infty$ for some $x \in A\left(x, x=x^{*}\right)$.

Proposition 15.7: For $n \geq 2$, the following are isomorphic:
(i) $M_{n}(\mathbb{C})$,
(ii) $C^{*}\left(e_{i, j}, i, i=1, \ldots, n \mid e_{i, j}^{*}=e_{j, i}, e_{i, j} e_{k, l}=\delta_{j, k} e_{i, l} \forall i, j, k, l\right)$,
(iii) $C^{*}\left(x_{1}, \ldots, x_{n} \mid x_{i}^{*} x_{j}=\delta_{i, j} x_{1}\right)$.

Proof: " $M_{n}(\mathbb{C}) \cong C^{*}\left(e_{i, j}, i, i=1, \ldots, n \mid e_{i, j}^{*}=e_{j, i}, e_{i, j} e_{k, l}=\delta_{j, k} e_{i, l} \forall i, j, k, l\right)$ ": Put $E_{i, j}:=\left(\delta_{i, k} \delta_{j, l}\right)_{1 \leq k, l \leq n} \in M_{n}(\mathbb{C})$. These $E_{i, j}$ satisfy the relations stated in the universal $C^{*}$-algebra in (ii), thus there is

$$
\begin{aligned}
\varphi: C^{*}\left(e_{i, j}, i, i=1, \ldots, n \mid e_{i, j}^{*}=e_{j, i}, e_{i, j} e_{k, l}=\delta_{j, k} e_{i, l} \forall i, j, k, l\right) & \longrightarrow M_{n}(\mathbb{C}), \\
e_{i, j} & \longmapsto E_{i, j}
\end{aligned}
$$

by the universal property, which is surjective. Since the universal $C^{*}$ algebra $C^{*}\left(e_{i, j}, i, i=1, \ldots, n \mid e_{i, j}^{*}=e_{j, i}, e_{i, j} e_{k, l}=\delta_{j, k} e_{i, l} \forall i, j, k, l\right)$ is $n^{2}$-dimensional, and so is $M_{n}(\mathbb{C}), \varphi$ is injective.

Note that $C^{*}\left(e_{i, j}, i, i=1, \ldots, n \mid e_{i, j}^{*}=e_{j, i}, e_{i, j} e_{k, l}=\delta_{j, k} e_{i, l} \forall i, j, k, l\right)$ exists, since $p\left(e_{i, j}\right)^{2}=p\left(e_{i, j}^{*} e_{i, j}\right)=p\left(e_{j, j}\right), p\left(e_{j, j}\right)^{2}=p\left(e_{j, j}\right) \in\{0,1\}$.

The other isomorphy is an exercise on Sheet 12:

such that $\varphi \circ \psi=\operatorname{id}_{C^{*}\left(x_{1}, \ldots, x_{n} \mid \ldots\right)}$ and $\psi \circ \varphi=\operatorname{id}_{C^{*}\left(e_{i, j}, 1 \leq i, j \leq n \mid \ldots\right)}$.
Remark 15.8: On Sheet $11 \pm \varepsilon$ in Functional Analysis I, we proved that $M_{n}(\mathbb{C})$ is simple. Hence, if $B$ is a $C^{*}$-algebra with $y_{1}, \ldots, y_{n} \in B$ and $y_{i}^{*} y_{j}=\delta_{i, j} y_{1}$, then $M_{n}(\mathbb{C}) \cong C^{*}\left(y_{1}, \ldots, y_{n}\right) \subseteq B$, because for

$$
\varphi: M_{n}(\mathbb{C}) \cong C^{*}\left(x_{1}, \ldots, x_{n}\right) \longrightarrow C^{*}\left(y_{1}, \ldots, y_{n}\right) \subseteq B
$$

we know that $\operatorname{ker}(\varphi) \triangleleft M_{n}(\mathbb{C})$; hence $\operatorname{ker}(\varphi)=\{0\}$.
Proposition 15.9: The following $C^{*}$-algebras are isomorphic:
(i) The compact operators $K(H)$ on a separable Hilbert space $H$,
(ii) $C^{*}\left(e_{i, j}, i, j \in \mathbb{N} \mid e_{i, j}^{*}=e_{j, i}, e_{i, j} e_{k, l}=\delta_{j, k} e_{i, l}\right)$,
(iii) $C^{*}\left(x_{i}, i \in \mathbb{N} \mid x_{i}^{*} x_{j}=\delta_{i, j} x_{1}\right)$

In (ii) and (iii) we may replace $\mathbb{N}$ by an arbitrary infinite countable set.

Proof: "(i) $\cong($ ii $)$ ": Let $\left(e_{n}\right)_{n \in \mathbb{N}}$ be an orthonormal basis of $H$. Define the operators $f_{i, j}\left(e_{n}\right):=\delta_{j, n} e_{i}$. Then $f_{i, j} \in K(H)$ since it is of finite rank. By the universal property, we find

$$
\varphi: C^{*}\left(e_{i, j}, i, j \in \mathbb{N} \mid \ldots\right) \longrightarrow K(H), \quad e_{i, j} \longmapsto f_{i, j}
$$

furthermore $\varphi$ is surjective, since the finite rank operators are dense in $K(H)$. It remains to be shown that $\varphi$ is injective. Put therefore

$$
M_{n}:=C^{*}\left(e_{i, j} \mid 1 \leq i, j \leq n\right) \subseteq C^{*}\left(e_{i, j}, i, j \in \mathbb{N} \mid \ldots\right)
$$

By (Proposition 15.7) and (Remark 15.8), we have $M_{n} \cong M_{n}(\mathbb{C})$. This shows that $\left.\varphi\right|_{M_{n}} \rightarrow K(H)$ is injective, since $\operatorname{ker}\left(\left.\varphi\right|_{M_{n}}\right) \triangleleft M_{n} \cong M_{n}(\mathbb{C})$; thus $\left.\varphi\right|_{M_{n}}$ is isometric. Now as $\varphi$ is isometric on the dense subset $\bigcup_{n \in \mathbb{N}} M_{n} \subseteq C^{*}\left(e_{i, j} i, j \in \mathbb{N} \mid \ldots\right), \varphi$ is isometric on $C^{*}\left(e_{i, j}, i, j \in \mathbb{N} \mid \ldots\right)$.
"(ii) $\cong(i i i) ":$ Exercise on Sheet 12 .
Remark 15.10: $K(H)$ is simple (this is an exercise on Sheet 12), so the analog of (Remark 15.8) holds true.

## 16 Example: Toeplitz algebra $\mathcal{T}$

Definition 16.1: We denote by $\mathcal{T}:=C^{*}\left(v, 1 \mid v^{*} v=1\right)$ the universal $C^{*}$-algebra generated by an isometry, the so called Toeplitz-algebra.

Theorem 16.2: (i) The ideal $\left\langle 1-v v^{*}\right\rangle$ in $\mathcal{T}$ (the ideal generated by $1-v v^{*}$ ) is isomorphic to $K(H)$ for some separable Hilbert space $H$ and we have

$$
\mathcal{T} /\left\langle 1-v v^{*}\right\rangle \cong C^{*}\left(\mathbb{S}^{1}\right)
$$

In other words, the sequence $0 \rightarrow K(H) \rightarrow \mathcal{T} \rightarrow C\left(\mathbb{S}^{1}\right) \rightarrow 0$ is exact.
(ii) The map

$$
\varphi: \mathcal{T} \longrightarrow C^{*}(S) \subseteq B\left(\ell^{2}(\mathbb{N}), \quad v \longmapsto S\right.
$$

where $S$ denotes the unilateral shift (defined by $S e_{n}=S e_{n+1}$ ), is an isomorphism.

Proof: (i) We want to prove the assertion in small steps:
(1) $1-v v^{*} \in \mathcal{T}$ is a projection.
(2) $\left\langle 1-v v^{*}\right\rangle=\operatorname{cl}\left(\operatorname{span}\left\{v^{k}\left(1-v v^{*}\right) v^{* l} \mid k, l \in \mathbb{N}\right\}\right)=: A \subseteq \mathcal{T}$.

Proof (of (2)): "?" is clear. " $\subseteq$ ": By construction, $A$ is a closed vector space. It holds $v x, v^{*} x \in A$ for all $x \in A$ since $v^{*}\left(1-v v^{*}\right)=0$, thus $y x, x y \in A$ for all $y \in T$ and $x \in A$, i. e., $A$ is an ideal.
(3) $f_{k, l}:=v^{k}\left(1-v v^{*}\right) v^{* l}$ satisfy the relations of $K(H)=C^{*}\left(e_{i, j}, i, j \in \mathbb{N}_{0} \mid \ldots\right)$, as we have the following:

$$
f_{i, j} f_{k, l}=v^{i}\left(1-v v^{*}\right) v^{* j} v^{k}\left(1-v v^{*}\right) v^{* l}=\delta_{j, k} v^{i}\left(1-v v^{*}\right) v^{* l} .
$$

Thus there is an injective map

$$
\eta: K(H) \longleftrightarrow \mathcal{T}, \quad e_{i, j} \longmapsto f_{i, j}
$$

with $\eta(K(H))=\left\langle 1-v v^{*}\right\rangle$ by (2).
(4) We have the following diagram:


But $\psi\left(\left\langle 1-v v^{*}\right\rangle\right)=0$, hence there is a map

$$
\beta: \mathcal{T} /\left\langle 1-v v^{*}\right\rangle \longrightarrow C\left(\mathbb{S}^{1}\right), \quad \pi(v) \longmapsto u
$$

thus $\alpha \circ \beta=\mathrm{id}, \beta \circ \alpha=\mathrm{id}$ thus $C\left(\mathbb{S}^{1}\right) \cong \mathcal{T} /\left\langle 1-v v^{*}\right\rangle$.
(ii) As for part (i), we want to show the assertion in small steps:
(1) $\varphi$ exists (since $S$ is an isometry) and is surjective. We need to show that $\|\varphi(x)\|=\|x\|$, but $\|x\|=\sup \{\|\rho(x)\| \mid \rho$ irreducible representation of $\mathcal{T}\}$ (refer to Remark 5.18 (iv)), hence we need to show $\|\rho(x)\| \leq\|\varphi(x)\|$ for all irreducible representations $\rho: \mathcal{T} \rightarrow B(H)$ (of course it holds $\|\varphi(x)\| \leq\|x\|$ ).
(2) Claim: Let $\rho: \mathcal{T} \rightarrow B(H)$ be irreducible, $\rho(p) \neq 0$ with $p:=1-v v^{*}$. Then $\rho$ and $\varphi$ are unitarily equivalent and then $\|\rho(x)\|=\left\|U \varphi(x) U^{*}\right\| \leq\|\varphi(x)\|$.
Proof (of (2)): We have: For all $x \in \mathcal{T}$ there is $\lambda \in \mathbb{C}$ such that $p x p=\lambda p$ (because if $x=v^{k} v^{* l}$ is a monomial with $k, l \in \mathbb{N}_{0}$, then

$$
p x p=\left(1-v v^{*}\right) v^{k} v^{* l}\left(1-v v^{*}\right)=\delta_{k, 0} \delta_{l, 0} p
$$

which transfers to polynomials and limits. Thus there is a state $f: \mathcal{T} \rightarrow \mathbb{C}, x \mapsto \lambda$ with $p x p=f(x) p$ for all $x \in \mathcal{T}$. Then for all $x \in \mathcal{T}$ we have

$$
\left\langle\pi_{f}(x) \xi_{f}, \xi_{f}\right\rangle=f(x)=\left\langle\varphi(x) e_{1}, e_{1}\right\rangle
$$

as $\left\langle\varphi(x) e_{1}, e_{1}\right\rangle=\left\langle S^{k} S^{* l} e_{1}, e_{1}\right\rangle=\delta_{0, l} \delta_{0, k}=f(x)$, if $x$ is a monomial, which again transfers to polynomials and limits. Therefore, $\pi_{f}$ and $\varphi$ are unitarily equivalent. Since $\rho(p) \neq 0$, we find $\xi \in H_{\rho}$ such that $\|\xi\|=1$ and $\rho(p) \xi=\xi$. Then

$$
\langle\rho(x) \xi, \xi\rangle=\langle\rho(x) \rho(p) \xi, \rho(p) \xi\rangle=\langle\rho(p x p) \xi, \xi\rangle=f(x)\langle\rho(p) \xi, \xi\rangle=f(x)
$$

hence also $\rho$ is unitarily equivalent to $f$.
(3) Let $\rho: \mathcal{T} \rightarrow B(H)$ be irreducible $\rho(p)=0$. Then $\rho(v)$ is a unitary. Moreover, $\left\langle 1-S S^{*}\right\rangle \triangleleft C^{*}(S)$ is isomorphic to $K(H)$. Hence $\pi: C^{*}(S) \rightarrow C^{*}(S) /\left\langle 1-S S^{*}\right\rangle$ satisfies that $\pi(S)$ is unitary. It remains to be shown that $\operatorname{Sp}(\pi(S))=\mathbb{S}^{1}$, as then we have the following diagram

and thus $\|\rho(x)\| \leq\|\alpha \circ \pi \circ \varphi(x)\| \leq\|\varphi(x)\|$.
Proof: Let $\lambda \in \mathbb{S}^{1}$. Consider $d(\lambda) e_{n}:=\lambda^{n} e_{n}$, then $d(\lambda) \in B(H)$ is unitary with $d(\lambda)^{*}=d(\bar{\lambda})$. We have $d(\lambda) S d(\lambda)^{*}=\lambda S$ thus

$$
\beta: C^{*}(S) \longrightarrow C^{*}(S), \quad x \longmapsto d(\lambda) x d(\lambda)^{*}
$$

is an automorphism with $\beta\left(\left\langle 1-S S^{*}\right\rangle\right) \subseteq\left\langle 1-S S^{*}\right\rangle$ and $\beta(S)=\lambda S$. Hence we have the map

$$
\dot{\beta}: C^{*}(S) /\left\langle 1-S S^{*}\right\rangle \xrightarrow{\cong} C^{*}(S) /\left\langle 1-S S^{*}\right\rangle, \quad \dot{\beta}(\pi(S))=\lambda \pi(S),
$$

which shows $\operatorname{Sp}(\pi(S))=\operatorname{Sp}(\dot{\beta}(\pi(S)))=\lambda \operatorname{Sp}(\pi(S))$ for all $\lambda \in \mathbb{S}^{1}$, thus finally $\operatorname{Sp}(\pi(S))=\mathbb{S}^{1}$.

16 Example: Toeplitz algebra $\mathcal{T}$

Remark 16.3: $L^{2}\left(\mathbb{S}^{1}\right)$ has the orthonormal basis $e^{n}:=\left(z \mapsto z^{n}\right)$ for $n \in \mathbb{Z}$. One can show that

$$
\mathcal{T} \longrightarrow\left\{T_{f}+k \mid f \in C\left(\mathbb{S}^{1}\right), k \in K\left(L^{2}\left(\mathbb{S}^{1}\right)\right)\right\}, \quad v \longmapsto T_{z}
$$

where $T_{f}:=P_{H^{2}} M_{f}, M_{f}(g):=f g$ for $g \in L^{2}\left(\mathbb{S}^{1}\right)$ and $H^{2}:=\operatorname{span}\left\{e_{n} \mid n \geq\right.$ $0\} \subseteq L^{2}\left(\mathbb{S}^{1}\right)$. Oftentimes, $\mathcal{T}$ is introduced in this way. $\left\{T_{f}+k \mid f \in C\left(\mathbb{S}^{1}\right), k \in\right.$ $\left.K\left(\overline{L^{2}}\left(\mathbb{S}^{1}\right)\right)\right\}$ are called the Toeplitz operators.

## 17 The (irrational) rotation algebra $A_{\vartheta}$

Definition 17.1: Let $\vartheta \in \mathbb{R}$. Put

$$
A_{\vartheta}:=C^{*}\left(u, v \text { unitaries } \mid u v=e^{2 \pi \mathrm{i} \vartheta} v u\right)
$$

and call $A_{\vartheta}$ rotation algebra. For $\vartheta \in \mathbb{R} \backslash \mathbb{Q}$, we call $A_{\vartheta}$ irrational rotation algebra.
Let in the following $\lambda:=e^{2 \pi \mathrm{i} \vartheta}$.
Remark 17.2: (i) $A_{\vartheta}$ exists, since $u, v$ are unitaries.
(ii) $A_{\vartheta=0} \cong C\left(\mathbb{T}^{2}\right)$, where $\mathbb{T}^{2}=\mathbb{S}^{1} \times \mathbb{S}^{1} \subseteq \mathbb{C}^{2}$ is the 2-torus, via

$$
u \longmapsto\left(\left(z_{1}, z_{2}\right) \mapsto z_{1}\right), \quad v \longmapsto\left(\left(z_{1}, z_{2}\right) \mapsto z_{2}\right) .
$$

The existence of this *-homomorphism is granted by the universal property and it is an isomorphism since $A_{\vartheta=0} \cong C\left(\operatorname{Spec}\left(A_{\vartheta=0}\right)\right) \cong C\left(\mathbb{T}^{2}\right)$ where the first isomorphy is ensured by the first fundamental theorem of $C^{*}$-algebras and the second isomorphy holds, because $\varphi \in \operatorname{Spec}\left(A_{\vartheta=0}\right)$ is determined by the values of $\varphi(u) \in \mathbb{S}^{1}$ and $\varphi(v) \in \mathbb{S}^{1}$.

Hence, for $\vartheta \neq 0$, we may view $A_{\vartheta}$ as a non-commutative function algebra on a "non-commutative torus $\mathbb{T}_{\vartheta}^{2}$ ".
(iii) $\vartheta \notin \mathbb{Q}$ behave very differently from $\vartheta \in \mathbb{Q}$; that's why mostly the irrational rotation algebras are studied (see Sheet 13 and Theorem 17.xx).
(iv) The non-commutative tori are examples of non-commutative manifolds in A. Connes "non-commutative geometry", an analog of differential geometry. Such a non-commutative differential geometry is supposed to be relevant for quantum physics ("At a deep and perhaps fundamental level, quantum field theory and non-commutative geometry are made of the same stuff.", Gracia-Bondia et. al, page 522), for instance for the quantum Hall effect.

Proposition 17.3: $A_{\vartheta}$ has some concrete representations (i.e., $A_{\vartheta} \neq 0$ ).
(i) Let $\tilde{S} \in B\left(\ell^{2}(\mathbb{Z})\right)$ be the bilateral shift and consider the map $d$ defined by $d(\lambda) e_{n}:=\lambda^{n} e_{n}$. We have

$$
\pi: A_{\vartheta} \longrightarrow B\left(\ell^{2}(\mathbb{Z}), \quad u \longmapsto d(\lambda), v \longmapsto \tilde{S} .\right.
$$

(ii) Let $\tilde{u}, \tilde{v}$ be the following operators:

$$
\begin{array}{ll}
\tilde{u}: L^{2}\left(\mathbb{S}^{1}\right) \longrightarrow L^{2}\left(\mathbb{S}^{1}\right), & (\tilde{u} f)(t):=f(\lambda t), \\
\tilde{v}: L^{2}\left(\mathbb{S}^{1}\right) \longrightarrow L^{2}\left(\mathbb{S}^{1}\right), & (\tilde{v} f)(t):=t f(t)
\end{array}
$$

We have

$$
\pi: A_{\vartheta} \longrightarrow B\left(L^{2}\left(\mathbb{S}^{1}\right)\right), \quad u \longmapsto \tilde{u}, v \longmapsto \tilde{v}
$$

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Proof: (i) $\tilde{S}$ is unitary and as $d(\lambda)^{*}=d(\bar{\lambda}), d(\lambda)$ is unitary as well. Now it holds

$$
d(\lambda) \tilde{S} e_{n}=d(\lambda) e_{n+1}=\lambda^{n+1} e_{n+1}=\lambda \tilde{S} d(\lambda) e_{n}
$$

hence $\pi$ exists.
(ii) Apply $\tilde{u} \tilde{v}$ to $e_{n}:=\left(z \mapsto z^{n}\right)$.

We ask: Are the representations in Proposition 17.3 isomorphic? Does

$$
C^{*}(\tilde{S}, d(\lambda)) \cong C^{*}(\tilde{u}, \tilde{v})
$$

hold? Are the maps $\pi$ in Proposition 17.3 injective? The answer to all these questions is "yes (if $\vartheta \notin \mathbb{Q}$ )", and the reason for this is that $A_{\vartheta}$ is simple (for $\vartheta \notin \mathbb{Q}$ ), which we now want to show.

Definition 17.4: Let $A$ be a unital $C^{*}$-algebra.
(i) If $B \subseteq A$ is a $C^{*}$-algebra such that $1_{A} \in B$ and $\varphi: A \rightarrow B$ is positive, linear, unital and $\varphi^{2}=\varphi$, we call $\varphi$ a (conditional) expectation. $\varphi$ is called faithful, if for all $a \geq 0$ with $\varphi(a)=0$ it holds $a=0$.
(ii) $\tau: A \rightarrow \mathbb{C}$ is a (normalised) trace, if $\tau$ is positive, linear and it holds $\tau(1)=1$, $\tau(a b)=\tau(b a)$ for all $a, b \in A . \tau$ is faithful, if for all $a \geq 0$ with $\varphi(a)=0$ it holds $a=0$.

Proposition 17.5: Let $\vartheta \notin \mathbb{Q}$.
(i) For $\xi, \mu \in \mathbb{S}^{1}$, the map

$$
\rho_{\xi, \mu}: A_{\vartheta} \longrightarrow A_{\vartheta}, \quad u \longmapsto \xi u, v \longmapsto \mu v
$$

is an automorphism (i.e., $a^{*}$-isomorphism).
(ii) The maps $\varphi_{1}, \varphi_{2}: A_{\vartheta} \rightarrow A_{\vartheta}$ defined via

$$
\varphi_{1}(x):=\int_{0}^{1} \rho_{1, e^{2 \pi \mathrm{i} t}}(x) d t, \quad \varphi_{2}:=\int_{0}^{1} \rho_{e^{2 \pi \mathrm{i} t}, 1}(x) d t
$$

are contractive (i.e., $\left\|\varphi_{i}\right\| \leq 1$ ), faithful expectations.
(iii) We have $\varphi_{1}\left(A_{\vartheta}\right) \subseteq C^{*}(u) \subseteq A_{\vartheta}, \varphi_{2}\left(A_{\vartheta}\right)=C^{*}(v) \subseteq A_{\vartheta},\left.\varphi_{1}\right|_{C^{*}(u)}=\operatorname{id}_{C^{*}(u)}$, $\left.\varphi_{2}\right|_{C^{*}(v)}=\operatorname{id}_{C^{*}(v)}$,
(iv) We have for any $a_{k, l} \in \mathbb{C}$ (where only finitely many $a_{k, l} \neq 0$ )

$$
\varphi_{1}\left(\sum_{k, l \in \mathbb{Z}} a_{k, l} u^{k} v^{l}\right)=\sum_{k \in \mathbb{Z}} a_{k, 0} u^{k}, \quad \varphi_{2}\left(\sum_{k, l \in \mathbb{Z}} a_{k, l} u^{k} v^{l}\right)=\sum_{l \in \mathbb{Z}} a_{0, l} v^{l}
$$

(v) For $\vartheta \notin \mathbb{Q}$ it holds

$$
\varphi_{1}(x)=\lim _{n \rightarrow \infty} \frac{1}{2 n+1} \sum_{j=-n}^{n} u^{j} x u^{-j}, \quad \varphi_{2}(x)=\lim _{n \rightarrow \infty} \frac{1}{2 n+1} \sum_{j=-n}^{n} v^{j} x v^{-j}
$$

This statement is false for $\vartheta \in \mathbb{Q}$ !
Proof: (i) $\rho_{\xi, \mu}$ exists by the universal property and because $\rho_{\bar{\xi}, \bar{\mu}} \circ \rho_{\xi, \mu}=\mathrm{id}$, $\rho_{\xi, \mu}$ is an automorphism.
(ii) Elements in $A_{\vartheta}$ may be approximated by $\sum_{k, l \in \mathbb{Z}} a_{k, l} u^{k} v^{l}$ with $a_{k, l} \neq 0$ finitely often. For $x \in A_{\vartheta}$, the map

$$
f_{x}: \mathbb{T}^{2} \longrightarrow A_{\vartheta}, \quad(\xi, \mu) \longmapsto \rho_{\xi, \mu}(x)
$$

is norm continuous, since for $x=\sum_{k, l=-n}^{n} a_{k, l} u^{k} v^{l}$ it holds

$$
\left\|f_{x}\left(\xi_{1}, \mu_{1}\right)-f_{x}\left(\xi_{2}, \mu_{2}\right)\right\| \leq \sum_{k, l=-n}^{n}\left|\xi_{1}^{k} \mu_{1}^{l}-\xi_{2}^{k} \mu_{2}^{l}\right|\left|a_{k, l}\right| \rightarrow 0 \text { for }\left(\xi_{1}, \mu_{1}\right) \rightarrow\left(\xi_{2}, \mu_{2}\right)
$$

thus

$$
g_{x}:[0,1] \longrightarrow A_{\vartheta}, \quad t \longmapsto f_{x}\left(1, e^{2 \pi \mathrm{i} t}\right)
$$

is also norm-continuous, hence $\frac{1}{n} \sum_{j=1}^{n} g_{x}\left(t_{j}\right) \rightarrow \int_{0}^{1} g_{x}(t) d t=: \varphi_{1}(x)$ exists as a limit of Riemann sums (for $0 \leq t_{0} \leq \cdots \leq t_{n}=1$ ).

To see that $\varphi_{1}$ is contractive, we notice that

$$
\left\|\varphi_{1}(x)\right\| \leq \frac{1}{n}\left\|\sum_{j} \rho_{1, e^{2 \pi \mathrm{i} t_{j}}}(x)\right\| \leq \frac{1}{n} \sum_{j}\|x\|=\|x\|
$$

$\varphi_{1}$ is positive because for $x \geq 0, \rho_{1, e^{2 \pi \mathrm{i} t}}(x) \geq 0$, hence $\varphi_{1}(x)$ is positive as limit of positive elements. One checks, that $\varphi_{1}$ is linear, unital and faithful.

Now, by (iii), $\varphi_{1}\left(\varphi_{1}(x)\right)=\left.\mathrm{id}\right|_{C^{*}(u)}\left(\varphi_{1}(x)\right)=\varphi_{1}(x)$, as $\varphi_{1}(x) \in C^{*}(u)$.
(iii) It holds

$$
\varphi_{1}\left(v^{l}\right)=\int_{0}^{1} \rho_{1, e^{2 \pi \mathrm{i} t}}\left(v^{l}\right) d t=\int_{0}^{1} e^{2 \pi \mathrm{i} l t} v^{l} d t=\left(\int_{0}^{1} e^{2 \pi \mathrm{i} l t} d t\right) v^{l}=\delta_{0, l}
$$

hence $\varphi_{1}\left(u^{k} v^{l}\right)=\int_{0}^{1} \rho_{1, e^{2 \pi \mathrm{it}}}\left(u^{k} v^{l}\right) d t=u^{k} \varphi_{1}\left(v^{l}\right)=\delta_{0, l} u^{k}$, so $\varphi_{1}\left(A_{\vartheta}\right) \subseteq C^{*}(u)$ and $\left.\varphi_{1}\right|_{C^{*}(u)}=\operatorname{id}_{C^{*}(u)}$.
(iv) This was also shown in (iii).
(v) We compute

$$
\frac{1}{2 n+1} \sum_{j=-n}^{n} u^{j}\left(u^{k} v^{l}\right) u^{-j}=\frac{1}{2 n+1}\left(\sum_{j=-n}^{n} \lambda^{j l}\right) u^{k} v^{l} \rightarrow \delta_{0, l} u^{k} v^{l}=\varphi_{1}\left(u^{k} v^{l}\right)
$$

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Corollary 17.6: Let $\vartheta \notin \mathbb{Q}$. The map

$$
\tau:=\varphi_{1} \circ \varphi_{2}=\varphi_{2} \circ \varphi_{1}: A_{\vartheta} \longrightarrow \mathbb{C}
$$

is the unique unital faithful trace on $A_{\vartheta}$. We have $\tau\left(\sum_{k, l \in \mathbb{Z}} a_{k, l} u^{k} v^{l}\right)=a_{0,0}$, this is "the $(0,0)$-th Fourier-coefficient".

Proof: It holds for $k, l \in \mathbb{Z}$ :

$$
\varphi_{1} \varphi_{2}\left(u^{k} v^{l}\right)=\delta_{k, 0} \varphi_{1}\left(v^{l}\right)=\delta_{k, 0} \delta_{l, 0}=\varphi_{2} \varphi_{1}\left(u^{k} v^{l}\right)
$$

hence $\tau$ is well-defined. Furthermore $\tau$ is linear, positive, unital and faithful as a composition of two maps with these properties, also

$$
\begin{aligned}
\tau\left(\left(u^{k} v^{l}\right)\left(u^{m} v^{n}\right)\right)=\lambda^{-l m} \tau\left(u^{m+k} v^{l+n}\right) & =\delta_{m+k, 0} \delta_{l+n, 0} \lambda^{-l m} \\
& =\lambda^{-} n k \tau\left(u^{m+k} v^{l+n}\right)=\tau\left(\left(u^{m} u^{n}\right)\left(u^{k} v^{l}\right)\right)
\end{aligned}
$$

One can check, that $\tau$ is unique.
Theorem 17.7: For $\vartheta \notin \mathbb{Q}, A_{\vartheta}$ is simple.
Theorem 17.7 is false for $\vartheta \in \mathbb{Q}$, see Exercise sheet 13 .
Proof: Let $0 \neq I \triangleleft A_{\vartheta}$, hence we find $0 \neq x \in I$, i.e., $0 \neq x^{*} x \in I$. Then

$$
0 \neq \varphi_{1}\left(x^{*} x\right)=\lim _{n \rightarrow \infty} \frac{1}{2 n+1} \sum_{j=-n}^{n} u^{j} x^{*} x u^{-j} \in I
$$

but $0 \neq \tau\left(x^{*} x\right)=\varphi_{2} \varphi_{1}\left(x^{*} x\right) \in I$, thus $1 \in I$, i. e., $I=A_{\vartheta}$ (because $\tau\left(x^{*} x\right) \in \mathbb{C}$ and $\left.\varphi_{1}\left(x^{*} x\right) \in I\right)$.

Remark 17.8: It holds $\vartheta^{\prime}= \pm \vartheta(\bmod \mathbb{Z})$ if and only if $A_{\vartheta} \cong A_{\vartheta^{\prime}}$ (" $\Rightarrow$ " is trivial, $" \Leftarrow$ " requires certain Powers-Rieffel projections).

## 18 Cuntz-Algebra

Definition 18.1: The Cuntz-Algebra $\mathcal{O}_{n}$ for $2 \leq n<\infty$ is defined as

$$
\mathcal{O}_{n}:=C^{*}\left(S_{1}, \ldots, S_{n} \text { isometries } \mid \sum_{i=1}^{n} S_{i} S_{i}^{*}=1\right)
$$

for $n=\infty$ put $\mathcal{O}_{\infty}:=C^{*}\left(S_{1}, S_{n}, \ldots\right.$ projections $\left.\mid S_{i}^{*} S_{j}=\delta_{i, j}\right)$.
Remark 18.2: (i) $\mathcal{O}_{n}$ is the structure of decomposing a space into $n$ copies.

(ii) J. Cuntz introduced these $C^{*}$-algebras in 1977. They are important (counter-) examples for certain questions. But they are also building blocks of the theory of $C^{*}$-algebras via theorems like:

- Let $A$ be a separable $C^{*}$-Algebra. $A$ is exact if and only if there is an embedding $A \hookrightarrow \mathcal{O}_{2}$.
- Let $A$ be a $C^{*}$-algebra. $A$ is a unital, simple, separable, nuclear $C^{*}$-algebra if and only if $A \cong A \otimes \mathcal{O}_{2}$.
- Let $A$ be a $C^{*}$-algebra. $A$ is purely infinte if and only if $A \cong A \otimes \mathcal{O}_{\infty}$.

The above stated theorems are called the Kirchberg-Phillips theorems.
(iii) Also, it has nice properties, like being purely infinite (a strengthening of being simple). We will prove this property in this chapter.

Definition 18.3: A word in $\mathcal{O}_{n}$ is $S_{\mu}:=S_{i_{1}} \cdots S_{i_{k}}$, where $\mu=\left(i_{1}, \ldots, i_{k}\right) \in$ $\{1, \ldots, n\}^{k}$ is a multi-index. We call $|\mu|=k$ the length of the multi-index $\mu$.

Lemma 18.4: (i) It holds $S_{i}^{*} S_{j}=\delta_{i, j}$,
(ii) Let $\mu, \nu$ be multi-indices with $|\mu|=|\nu|$. Then $S_{\mu}^{*} S_{\nu}=\delta_{\mu, \nu}$, where $\delta_{\mu, \nu}$ is defined standing to reason.
(iii) Let $\mu, \nu$ be multi-indices. Then we have the following:
(1) If $|\mu|<|\nu|$, then

$$
S_{\mu}^{*} S_{\nu}= \begin{cases}S_{n u^{\prime}}, & \text { if } \nu=\mu \nu^{\prime} \\ 0, & \text { otherwise }\end{cases}
$$

(2) If $|\mu|>|\nu|$, then

$$
S_{\mu}^{*} S_{\nu}= \begin{cases}S_{\mu^{\prime}}^{*}, & \text { if } \mu=\nu \mu^{\prime} \\ 0, & \text { otherwise }\end{cases}
$$

Proof: (i) The case $i=j$ is clear: $S_{i}^{*} S_{i}=1$, because $S_{i}$ are isometries.
For $i \neq j$, we have $S_{i} S_{i}^{*}+S_{j} S_{j}^{*} \leq \sum_{k} S_{k} S_{k}^{*}=1$, thus $S_{i}^{*} S_{j} S_{j}^{*} S_{i} \leq 0$. Hence

$$
S_{i}^{*}\left(S_{i} S_{i}^{*}+S_{j} S_{j}^{*}\right) S_{i} \leq S_{i}^{*} S_{i}=S_{i}^{*}\left(1+S_{i}^{*} S_{j} S_{j}^{*} S_{i}\right) S_{i}=1
$$

but as $S_{i}^{*} S_{j} S_{j}^{*} S_{i}=\left(S_{j}^{*} S_{i}\right)^{*}\left(S_{j}^{*} S_{i}\right) \geq 0$, we get that $\left(S_{j}^{*} S_{i}\right)^{*}\left(S_{j}^{*} S_{i}\right)=0$, so finally $S_{j}^{*} S_{i}=0$.
(ii) Exercise on Sheet 13,
(iii) Exercise on Sheet 13.

Proposition 18.5: (i) For $k \in \mathbb{N}$ define

$$
\mathcal{F}_{k}^{n}:=\operatorname{span}\left\{S_{\mu} S_{\nu}^{*}| | \mu|=|\nu|=k\} \subseteq \mathcal{O}_{n} .\right.
$$

Note that $\mathcal{F}_{k}^{n} \cong M_{n^{k}}(\mathbb{C})$, here $S_{1}^{k} S_{1}^{k *} \leftrightarrow e_{1,1}$.
(ii) The set $\left\{S_{\mu} S_{\nu}^{*} \mid \mu, \nu\right.$ arbitary $\} \subseteq \mathcal{O}_{n}$ is dense.

Proof: (i) Put $e_{\mu \nu}:=S_{\mu} S_{\nu}^{*} \in \mathcal{F}_{k}^{n}$. Then we have the relations $e_{\mu \nu}^{*}=e_{\nu \mu}$ as well as $e_{\mu \nu} e_{\rho \sigma}=S_{\mu} S_{\nu}^{*} S_{\rho} S_{\sigma}^{*}=\delta_{\nu \rho} e_{\mu \sigma}$. Due to $\left|\left\{\mu \in\{1, \ldots, n\}^{k}\right\}\right|=n^{k}$, we obtain that

$$
\varphi: M_{n^{k}}(\mathbb{C}) \longrightarrow \mathcal{F}_{n}^{k}, \quad e_{\mu \nu} \longmapsto S_{\mu} S_{\nu}^{*}
$$

is a *-isomorphism.
(ii) Monomials in $\mathcal{O}_{n}$ are of the form $S_{\mu} S_{\nu}^{*}$ by (Lemma 18.4).

Lemma 18.6: (i) Let $k \in \mathbb{N}$. Then

$$
\sum_{\substack{\text { multi-index, } \\|\delta|=k}} S_{\delta} S_{\delta}^{*}=1
$$

(ii) For $l \leq k$, we have $\mathcal{F}_{l}^{n} \subseteq \mathcal{F}_{k}^{n}$.

Proof: (i) This is exercise 2 b ) on sheet 13.
(ii) Let $S_{\mu} S_{\nu}^{*} \in \mathcal{F}_{l}^{n}$ where $l \leq k$ and $|\mu|=|\nu|=l$. Then by (i) it holds

$$
\sum_{\substack{\delta \text { multi-index, } \\|\delta|=k-l}} S_{\nu} S_{\delta} S_{\delta}^{*} S_{\nu}^{*} \in \mathcal{F}_{k}^{n}
$$

Lemma 18.7: Let $\mu, \nu$ be multi-indices with $|\mu|,|\nu| \leq k,|\mu| \neq|\nu|$. Then, for $S_{\gamma}:=S_{1}^{k} S_{2}$,

$$
S_{\gamma}^{*}\left(S_{\mu} S_{\nu}^{*}\right) S_{\gamma}=0
$$

Proof: Exercise 2 c) on Sheet 13.
Lemma 18.8: Let $k \in \mathbb{N}$. Then there is an isometry $w \in \mathcal{O}_{n}$ such that
(i) $w x=x w$ for all $x \in \mathcal{F}_{k}^{n}$,
(ii) For all multi-indices $\mu, \nu$ with $|\mu|,|\nu| \leq k$, it holds

$$
w^{*}\left(S_{\mu} S_{\nu}^{*}\right) w= \begin{cases}S_{\mu} S_{\nu}^{*}, & \text { if }|\mu|=|\nu| \\ 0 & \text { otherwise }\end{cases}
$$

Proposition 18.9: There is a faithful expectation

$$
\varphi: \mathcal{O}_{n} \longrightarrow \mathcal{F}^{n}
$$

where $\mathcal{F}^{n}=\operatorname{cl}\left(\operatorname{span}\left\{F_{k}^{n} \mid k \in \mathbb{N}\right\}\right)=\operatorname{span}\left\{S_{\mu} S_{\nu}^{*}| | \mu|=|\nu|\}\right.$, with

$$
\varphi\left(S_{\mu} S_{\nu}^{*}\right)= \begin{cases}S_{\mu} S_{\nu}^{*}, & \text { if }|\mu|=|\nu|, \\ 0, & \text { otherwise }\end{cases}
$$

and $\varphi(x)=w^{*} x w \in \mathcal{F}_{k}^{n}$ for $x \in \operatorname{span}\left\{S_{\mu} S_{\nu}^{*}| | \mu|,|\nu| \leq k\}\right.$ and $w$ from (Proposition 18.8).

Proof: For $\zeta \in \mathbb{S}^{1}$, the map

$$
\rho_{\zeta}: \mathcal{O}_{n} \longrightarrow \mathcal{O}_{n}, \quad S_{i} \longmapsto \zeta S_{i}
$$

is an isomorphism with $\rho_{\bar{\zeta}} \circ \rho_{\zeta}=\mathrm{id}$. For $x \in \mathcal{O}_{n}$,

$$
f_{x}: \mathbb{S}^{1} \longrightarrow \mathcal{O}_{n}, \quad \zeta \longmapsto \rho_{\zeta}(x)
$$

is norm-continuous (like in the $A_{\vartheta}$-proof). Then $\varphi(x):=\int_{0}^{1} f_{x}\left(e^{2 \pi \mathrm{i} t}\right) d t$ is positive, linear, unital and faithful, where again

$$
\frac{1}{n} \sum_{j=1}^{n} f_{x}\left(e^{2 \pi \mathrm{i} t_{j}}\right) \longrightarrow \int_{0}^{1} f_{x}\left(e^{2 \pi \mathrm{i} t}\right) d t
$$

It holds

$$
\varphi\left(S_{\mu} S_{\nu}^{*}\right)=\int_{0}^{1} \rho_{e^{2 \pi \mathrm{it}}}\left(S_{\mu} S_{\nu}^{*}\right) d t=\left(\int_{0}^{1} e^{2 \pi \mathrm{i} t(|\mu|-|\nu|)} d t\right) S_{\mu} S_{\nu}^{*}=\delta_{|\mu|,|\nu|} S_{\mu} S_{\nu}^{*}
$$

thus $\varphi^{2}=\varphi$ and $\varphi\left(S_{\mu} S_{\nu}^{*}\right)=w^{*} S_{\mu} S_{\nu}^{*} w$ and $\varphi\left(S_{\mu} S_{\nu}^{*}\right) \in \mathcal{F}_{k}^{n}$ by (Lemma 18.6).

Definition 18.10: Let $A$ be a unital $C^{*}$-algebra. $A$ is purely infinite if and only if for all $0 \neq x \in A$ there are $a, b \in A$ such that $a x b=1$.

Remark 18.11: (i) If $A$ is purely infinite, then $A$ is simple. To see this, let $0 \neq I \triangleleft A$ be an ideal. Because $A$ is purely infinite, for $0 \neq x \in I$, there are $a, b \in A$ such that $1=a x b \in I$, i. e. $I=A$.
(ii) The notion "purely infinite" comes from von Neumann algebras: If $M$ is a von Neumann algebra of type III, then $M$ has no finite projections. Let now $A$ be a $C^{*}$-algebra. $A$ is purely infinite if and only if it holds: "For all hereditary $C^{*}$-subalgebras $B \subseteq A$ (i.e., if $0 \leq a \leq b$ and $b \in B$, then $a \in B$ ), $B$ has a finite projection".

Theorem 18.12: $\mathcal{O}_{n}$ is purely infinite (and thus in particular simple).
Proof: (i) Let $0 \neq x \in \mathcal{O}_{n}$, then $x^{*} x \nexists 0$, also $\varphi\left(x^{*} x\right) \nRightarrow 0$. Without loss of generality, $\left\|\varphi\left(x^{*} x\right)\right\|=1$.
(ii) Find $y=y^{*} \in \operatorname{span}\left\{S_{\mu} S_{\nu}^{*} \mid \mu, \nu\right.$ arbitrary $\}$ close to $x^{*}, x$, i. e., $\left\|x^{*} x-y\right\|<\frac{1}{4}$. This is possible: Because $\operatorname{cl}\left(\operatorname{span}\left\{S_{\mu} S_{\nu}^{*} \mid \mu, \nu\right.\right.$ arbitrary $\left.\}\right)=\mathcal{O}_{n}$, we find $y_{0} \in$ $\operatorname{span}\left\{S_{\mu} S_{\nu}^{*} \mid \mu, \nu\right.$ arbitrary $\}$, such that $\left\|x^{*} x-y_{0}\right\|<\frac{1}{4}$; then put $y:=2^{-1}\left(y_{0}+y_{0}^{*}\right)$.
(iii) Find $z \in \mathcal{O}_{n}$ such that $z y z^{*}=1$.

Proof: We have $\|\varphi(y)\|>\frac{3}{4}$ since

$$
1=\left\|\varphi\left(x^{*} x\right)\right\| \leq\left\|\varphi\left(x^{*} x-y\right)\right\|+\|\varphi(y)\|<\frac{1}{4}+\|\varphi(y)\| .
$$

Let now $y=\sum_{i=1}^{n} a_{i} S_{\mu_{i}} S_{\nu_{i}}^{*}$ and $k \geq \max \left\{\left|\mu_{i}\right|,\left|\nu_{i}\right|\right\}$. Thus there is an isometry $w \in \mathcal{O}_{n}$ with $\varphi(y)=w^{*} y w$. Since $\varphi(y) \in \mathcal{F}_{k}^{n} \cong M_{n^{k}}(\mathbb{C})$, we may view $\varphi(y)$ as a matrix which we can diagonalise, thus there exists a one-dimensional projection $e \in \mathcal{F}_{k}^{n}$ such that $e \varphi(y)=\varphi(y) e= \pm\|\varphi(y)\| e$ and a unitary $u \in \mathcal{F}_{k}^{n}$ such that $u e u^{*}=S_{1}^{k} S_{1}^{* k}$ (corresponding to $e_{1,1}$ ). Now put

$$
z:=\|\varphi(y)\|^{-\frac{1}{2}} S_{1}^{* k} u e w^{*}
$$

It holds $z \in \mathcal{O}_{n}$ and $\|z\|<\frac{2}{\sqrt{3}}$, because

$$
\|z\| \leq\|\varphi(y)\|^{-\frac{1}{2}}\left\|S_{1}^{* k} u e w^{*}\right\| \leq\|\varphi(y)\|^{-\frac{1}{2}}<\frac{2}{\sqrt{3}}
$$

and $z y z^{*}=1$, because

$$
z y z^{*}=\|\varphi(y)\|^{-1} S_{1}^{* k} u e w^{*} y w e u^{*} S_{1}^{k}=\cdots=1
$$

(iv) $z x^{*} x z^{*}$ is invertible, since

$$
\left\|1-z x^{*} x z^{*}\right\|=\left\|z\left(y-x^{*} x\right) z^{*}\right\| \leq\|z\|^{2}\left\|y-x^{*} x\right\| \leq \frac{4}{3} \cdot \frac{1}{4}=\frac{1}{3}<1
$$

Put $a:=b^{*} x^{*}, b:=z^{*}\left(z x^{*} x z^{*}\right)^{-\frac{1}{2}}$, then

$$
a x b=b^{*} x^{*} x b=\left(z x^{*} x z^{*}\right)^{-\frac{1}{2}} z x^{*} x z^{*}\left(z x^{*} x z^{*}\right)^{-\frac{1}{2}}=1 .
$$

Remark 18.13: (i) There is a generalisation of $\mathcal{O}_{n}$ to graph $C^{*}$-algebras.
(ii) $\mathcal{O}_{n} \not \not \mathcal{O}_{m}$ for $n \neq m$.

## 19 Group $C^{*}$-algebras

Motivation: Let $G$ be a locally compact group, furthermore let $G$ be abelian. Then

$$
\widehat{G}:=\{\varphi: G \rightarrow \mathbb{T} \subseteq \mathbb{C} \text { group homomorphism }\}
$$

is the dual group of $G$. $\widehat{G}$ is again an abelian locally compact group. Pontrjagin duality states that

$$
G \longrightarrow \hat{\hat{G}}, \quad x \longmapsto \mathrm{ev}_{x}
$$

is an isomorphism of topological groups. Hence $G$ and $\widehat{G}$ are dual to each other. For example we have the following dual groups:

- $\widehat{\mathbb{Z}}=\mathbb{T}$ (think of the Fourier transformation),
- $\widehat{\mathbb{Z} / n \mathbb{Z}}=\mathbb{Z} / n \mathbb{Z}$,
- $\widehat{\mathbb{R}}=\mathbb{R}$.

What about non-abelian groups $G$ ? Then $\widehat{G}$ has too little information and is no group in general. Instead, consider $\{G \rightarrow B(H)$ unitary representation $\}$ - note that if $G$ is abelian and $\operatorname{dim} H=1$, then $\{G \rightarrow B(H)$ unitary representation $\}=\widehat{G}$. By some Schur-Weyl / Tannaka-Krein / Peter-Weyl duality, we may reconstruct $G$ from its representation theory.

Consider $C^{*}(G)$ or $C_{\mathrm{rad}}^{*}(G)$ whose representation theory is intuitively linked with the representation theory of $G$, i.e., group $C^{*}$-algebras arise from the idea to study the representation theory of groups by means from $C^{*}$-algebra theory. In fact, this was one of the reasons for introducing $C^{*}$-algebras as a concept (see J. Rosenberg, $C^{*}$-algebras and Mackey's theory of group representations, 1994 in Doran, R (ed.): $C^{*}$-algebras 1943-1993: A fifty year celebration, 1994).

By the way, one could also solve the problem of finding Pontrjagin duality in the non-abelian case by the following:

> Sketch missing
which leads to topological quantum groups.
In conclusion: $C^{*}(G), C_{\mathrm{rad}}^{*}(G)$ help to understand the representation theory of groups, but also they provide a huge class of examples of $C^{*}$-algebras.

Definition 19.1: Let $G$ be a locally compact group. A unitary representation of $G$ is a group homomorphism

$$
\pi: G \longrightarrow \mathfrak{U}(H):=\{u \in B(H) \text { unitary }\}
$$

such that

$$
G \longrightarrow H, \quad g \longmapsto \pi(g) \xi
$$

is continuous for all $\xi \in H$, i.e., $\pi$ is continuous in the strong operator topology.

Remark 19.2: On $\mathfrak{U}(H)$, the strong- and the weak operator topology coincide. To see this, consider a net $\left(u_{\lambda}\right)_{\lambda \in \Lambda}$ in $\mathfrak{U}(H)$ such that $u_{\lambda} \rightarrow$ wot $u$. Then

$$
\begin{aligned}
\left\|u_{\lambda} \xi-u \xi\right\|^{2} & =\left\|u_{\lambda} \xi\right\|^{2}-\left\langle u_{\lambda} \xi, u \xi\right\rangle-\left\langle u \xi, u_{\lambda} \xi\right\rangle+\|u \xi\|^{2} \\
& \rightarrow\|\xi\|-\|u \xi\|^{2}-\|u \xi\|^{2}+\|\xi\|^{2}=0 .
\end{aligned}
$$

Theorem 19.3: Let $G$ be a locally compact group. Then there exists a left-invariant Radon measure (dt or $\mu_{G}$ ) on $G$, unique up to a multiple by a constant. Hence for all $s \in G$ and $f \in L^{1}(G, d t)$ it holds

$$
\int_{G} f(t) d t=\int_{G} f(s t) d t
$$

$d t$ respectively $\mu_{G}$ is called the Haar measure.
Proof (idea): Construct on $C_{c}(G):=\{f: G \rightarrow \mathbb{C}$ continuous, $\operatorname{supp}(f)$ compact $\}$ a positive linear functional $\Lambda: C_{c}(G) \rightarrow \mathbb{C}$ such that $\Lambda(f)=\lambda\left({ }_{s} f\right)$, where we denote ${ }_{s} f(t):=f(s t)$. By the Riesz representation theorem, there is a measure $\mu_{G}$ such that $\int_{G} f d \mu_{G}=\Lambda(f)$.

Remark 19.4: (i) If $G$ is discrete, the Haar measure on $G$ is the counting measure (up to normalisation $\mu_{G}(\{e\})=1$ ).
(ii) If $G$ is compact, one usually normalises $\mu_{G}(G)=1$.
(iii) In case $G=\mathbb{R}^{n}$, the Haar measure coincides with the Lebesgue measure.

Proposition 19.5: The space $L^{1}(G):=L^{1}(G, d t)$ is $a^{*}$-Banach algebra via

$$
(f \star g)(s):=\int_{G} f(t) g\left(t^{-1} s\right) d t, \quad f^{*}(s):=\Delta(s)^{-1} \overline{f\left(s^{-1}\right)}, \quad\|f\|_{1}:=\int_{G}|f(t)| d t
$$

for $f, g \in L^{1}(G)$ and $s \in G$. Here, $\Delta: G \rightarrow \mathbb{R}_{+}$is the "modular function" with $\mu_{G}\left(E_{s}\right)=\Delta(s) \mu_{G}(E)$, where $E_{s}:=\{t s \mid t \in E\}$ (note that $\mu_{G}$ is left-invariant, i.e., $\mu_{G}\left({ }_{s} E\right)=\mu_{G}(E)$, but not right-invariant). $G$ is called unimodular, if $\Delta(s) \equiv 1$. $f \star g$ is called the convulution of $f$ and $g$, which shall be the product on $L^{1}(G)$.

Proof (Sketch): For $f, g \in L^{1}(G)$, Fubini's theorem yields

$$
\begin{aligned}
\|f \star g\|_{1} & =\int_{G}|f \star g(s)| d s \\
& \leq \int_{G} \int_{G}\left|f(t)\left\|g\left(t^{-1} s\right)\left|d s d t=\int_{G}\right| f(t) d t \int_{G}|g(s)| d s=\right\| f\left\|_{1}\right\| g \|_{1}\right.
\end{aligned}
$$

thus $\|\cdot\|_{1}$ is submultiplicative with resect to the convolution and $f \star g \in L^{1}(G)$.
$\Delta$ is a group homomorphism, thus

$$
f^{* *}(s)=\Delta(s)^{-1} \overline{f^{*}\left(s^{-1}\right)}=\Delta\left(s^{-1}\right) \Delta\left(s^{-1}\right)^{-1} f(s)=f(s)
$$

Remark 19.6: (i) If $G$ is abelian, discrete or compact, then $G$ is unimodular.
(ii) If $G$ is discrete, then $\mu_{G}$ is the counting measure and for characteristic functions $\delta_{t}, t \in G$, we have

$$
\left(\delta_{t_{1}} \star \delta_{t_{2}}\right)(s)=\sum_{g \in G} \delta_{t_{1}}(g) \delta_{t_{2}}\left(g^{-1} s\right)=\delta_{t_{1}, t_{2}}(s)
$$

More generally,

$$
\left(\sum_{i} \alpha_{i} \delta_{t_{i}}\right)\left(\sum_{j} \beta_{j} \delta_{j}\right)=\sum_{i, j} \alpha_{i} \beta_{j} \delta_{t_{i}, s_{j}},
$$

thus convolution corresponds to multiplication of indices; furthermore $\delta_{s}^{*}=\delta_{s^{-1}}$.
(iii) We always have that $C_{c}(G) \subseteq L^{1}(G)$ is dense (the Haar measure is finite on compact sets).
(iv) Let $G$ be discrete. Then

$$
\mathbb{C} G=\left\{\sum_{g \in G} \alpha_{g} g: \alpha_{g} \in \mathbb{C}, \alpha_{g} \neq 0 \text { for finitely many } g\right\}
$$

is called the group algebra, that becomes an algebra with the multiplication

$$
\left(\sum_{g} \alpha_{g} g\right)\left(\sum_{h} \beta_{h} h\right)=\sum_{g, h} \alpha_{g} \beta_{h} g h .
$$

It holds $\mathbb{C} G=C_{c}(G)$, hence $L^{1}(G)$ is the completion of $\mathbb{C} G$ to a Banach algebra. There is an involution on $\mathbb{C} G$ via $\delta_{g}^{*}=\delta_{g^{-1}}$ (extend linearly).

Proposition 19.7: Let $G$ be locally compact.
(i) $G$ is discrete if and only if $L^{1}(G)$ is unital; $L^{1}(G)$ is unital if and only if $\mu_{G}(\{e\}) \neq 0$ ( $\delta_{e}$ then is the unit),
(ii) $G$ is commutative if and only if $L^{1}(G)$ is commutative.

Remark 19.8: $\left(L^{1}(G),\|\cdot\|_{1}\right)$ is no $C^{*}$-algebra (consider for instance the function $f=-\delta_{t^{-1}}+\delta_{e}+\delta_{t}$, then $\|f\|^{2}=9$, but $\left\|f^{*} \star f\right\| \leq 5$ ). Still, for any ${ }^{*}$-Banach algebra, there are two canonical way to turn it into a $C^{*}$-algebra (see Dixmier, 2.7.1).

Definition 19.9: Let $G$ be a locally compact group. Put for $f \in L^{1}(G)$ :

$$
\|f\|:=\sup \left\{\|\pi(f)\| \mid \pi \text { representation of } L^{1}(G)\right\} \leq\|f\|_{1} .
$$

The completion of $L^{1}(G)$ with respect to this norm is the full (or maximal) group $C^{*}$-algebra $C^{*}(G)$ or $C_{\max }^{*}(G)$ or $C_{f}^{*}(G)$.

Proposition 19.10: Let $G$ be a locally compact group.
(i) Every unitary representation $\pi: G \rightarrow B(H)$ induces a representation $\tilde{\pi}$ of $L^{1}(G)$.
(ii) The left regular representation

$$
\lambda: G \longrightarrow B\left(L^{2}(G)\right), \quad(\lambda(s) f)(t):=f\left(s^{-1} t\right)
$$

has an extension $\tilde{\lambda}$ to a faithful representation of $L^{1}(G)$.
(iii) $\|\cdot\|$ from Definition 19.9 is indeed a norm.

Proof (Sketch): (i) For $f \in L^{1}(G)$ put $\tilde{\pi}(f):=\int f(t) \pi(t) d t$, i. e.,

$$
\langle\tilde{\pi} \xi, \eta\rangle:=\int_{G} f(t)\langle\pi(t) \xi, \eta\rangle d t
$$

This $\tilde{\pi}$ is indeed a representation.
(ii) Use (i) and check for faithfulness.
(iii) Let $0 \neq x \in L^{1}(G)$. Then $\tilde{\lambda}(x) \neq 0$, but as $\tilde{\lambda}(x)$ appears in the supremum, $\|x\| \neq 0$. The rest is clear, because the supremum is taken over $C^{*}$-seminorms.

Definition 19.11: The reduced group $C^{*}$-algebra $C_{\mathrm{rad}}^{*}(G)$ or $C_{\mathrm{min}}^{*}(G)$ or $C_{r}^{*}(G)$ is given by

$$
C_{\mathrm{red}}^{*}(G):=\operatorname{cl}\left(\tilde{\lambda}\left(L^{1}(G)\right)\right) \subseteq B\left(L^{2}(G)\right) .
$$

We always have a homomorphism $C^{*}(G) \rightarrow C_{\text {red }}^{*}(G)$ which is an isomorphism if and only if $G$ is amenable and $L(G)=C_{\mathrm{red}}^{*}(G)^{\prime \prime} \subseteq B\left(L^{2}(G)\right)$.

Remark 19.12: Let $G$ be discrete.
(i) Then $G$ is locally compact and the Haar measure is the counting measure. Also, $G$ is unimodular.
(ii) We can define $C_{\max }^{*}(G):=\mathrm{cl}_{\|\cdot\|}(\mathbb{C} G)$ with the norm

$$
\|x\|:=\sup \left\{\|\pi(x)\| \mid \pi \text { is }^{*} \text {-representation of } \mathbb{C} G\right\}<\infty
$$

and $C_{\max }^{*}=\mathrm{cl}_{\|\cdot\|}\left(L^{1}(G)\right)$ from Definition 19.9 coincides with this definition, because $\mathbb{C} G=C_{c}(G) \subseteq L^{1}(G)$ is dense.
(iii) We have $C_{\max }^{*}(G) \cong C^{*}\left(u_{g}, g \in G \mid u_{g}\right.$ unitary, $\left.u_{g} u_{h}=u_{g h}, u_{g}^{*}=u_{g^{-1}}\right)$, where $u_{e}=1$ in the universal $C^{*}$-algebra.

Proof: By the universal property, we have a map

$$
\varphi: C^{*}\left(u_{g}, g \in G \mid u_{g} \text { unitary, } u_{g} u_{h}=u_{g h}, u_{g}^{*}=u_{g^{-1}}\right) \longrightarrow C_{\max }^{*}(G), \quad u_{g} \longmapsto \delta_{g}
$$

Note that $u_{e} u_{g}=u_{g} u_{e}=u_{g}, u_{e}^{2}=u_{e}$, thus $u_{e}$ is the unit. $\varphi$ is surjective, since $\mathbb{C} G \subseteq \operatorname{im}(\varphi)$ and $\mathbb{C} G$ is dense. For the injectivity of $\varphi$, let

$$
\pi: C *\left(u_{g}, g \in G \mid u_{g} \text { unitary, } u_{g} u_{h}=u_{g h}, u_{g}^{*}=u_{g^{-1}}\right) \longrightarrow B(H)
$$

be a faithful representation. Then

$$
\alpha^{\circ}: G \longrightarrow B(H), \quad g \longmapsto \pi\left(u_{g}\right)
$$

is a unitary representation of $G$, thus there is an extension $\alpha: C_{\max }^{*}(G) \rightarrow B(H)$, hence

commutes and if $\varphi(x)=0$, then $\pi(x)=\alpha \circ \varphi(x)=0$ and thus, as $\pi$ is injective, $x=0$.
(iv) If $\varphi: G \rightarrow B(H)$ is a unitary representation of the group $G$, then

$$
\tilde{\varphi}: C_{\max }^{*}(G) \longrightarrow B(H), \quad \delta_{g} \longmapsto \varphi(g),
$$

is a representation the full group- $C^{*}$ algebra. Conversely, if a representation $\tilde{\varphi}: C_{\max }^{*}(G) \rightarrow B(H)$ is given, then we get a representation of the group $G$ via

$$
\varphi: G \longrightarrow B(H), \quad g \longmapsto \tilde{\varphi}\left(\delta_{g}\right) .
$$

(v) Let $\lambda: G \rightarrow B\left(\ell^{2}(G)\right)$ be a left regular representation. Then $\lambda(g) \delta_{h}=\delta_{g h}$ is just left multiplication. Thus there is a faithful representation of $\mathbb{C} G$ via

$$
\tilde{\lambda}: \mathbb{C} G \longrightarrow B\left(\ell^{2}(G)\right), \quad \sum \alpha_{g} \delta_{g} \longmapsto \sum \alpha_{g} \lambda(g) .
$$

To see that $\tilde{\lambda}$ is faithful, let $x=\sum \alpha_{g} \delta_{g} \neq 0$, i. e., there is $g_{0}$ such that $\alpha_{g_{0}} \neq 0$. Then

$$
\left\langle\tilde{\lambda}(x) \delta_{e}, \delta_{g_{0}}\right\rangle=\left\langle\sum \alpha_{g} \delta_{g}, \delta_{g_{0}}\right\rangle=\alpha_{g_{0}} \neq 0
$$

Because $\tilde{\lambda}$ is faithful, we have $\mathbb{C} G \cong \tilde{\lambda}(\mathbb{C} G) \subseteq B\left(\ell^{2}(G)\right)$ and thus can define $C_{\text {red }}^{*}(G):=\operatorname{cl}(\tilde{\lambda}(\mathbb{C} G)) \subseteq B\left(\ell^{2}(G)\right)$.
(vi) By (iv), the left regular representation extends to $\tilde{\lambda}: C_{\max }^{*}(G) \rightarrow B\left(\ell^{2}(G)\right)$, so in fact $\tilde{\lambda}: C_{\max }^{*}(G) \rightarrow C_{\text {red }}^{*}$ is a surjective *-homomorphism, $\lambda$ is no isomorphism. One can show there are surjective ${ }^{*}$-homomorphisms

$$
C_{\max }^{*}(G) \longrightarrow \mathrm{cl}_{\|\cdot\|}(\mathbb{C} G) \longrightarrow C_{\mathrm{red}}^{*}
$$

identity on $\mathbb{C} G$, hence: $C_{\max }^{*}(G)$ is the maximal- and $C_{\min }^{*}$ is the minimal $C^{*}$ completion of $\mathbb{C} G$. This map from $C_{\max }^{*}(G)$ to $C_{\text {red }}^{*}$ exists for general locally compact groups and we have that $\tilde{\lambda}$ is an isomorphism if and only if $G$ is amenable. $G$ is called amenable, if there is a state $m: L^{\infty}(G) \rightarrow \mathbb{C}$ such that $m\left(f_{s}\right)=m(f)$ for all $f \in L^{\infty}(G), s \in G$ and $f_{s}(t):=f\left(s^{-1} t\right) .{ }^{1}$ For example: If $G$ is compact, finite or abelian, then $G$ is amenable.
(vii) $L(G):=\tilde{\lambda}(\mathbb{C} G)^{\prime \prime}=C_{\text {red }}^{*}(G)^{\prime \prime} \subseteq B\left(\ell^{2}(G)\right)$ is the group von Neumann algebra for a group $G$.
(viii) If $G$ is abelian and locally compact, then $C_{\max }^{*}(G) \cong C_{0}(\widehat{G})$.
(ix) There also are right regular representations, but those yield the same $C^{*}$ algebras.
(x) We have a trace $\tau(x)=\left\langle x \delta_{e}, \delta_{e}\right\rangle$ on $C_{\mathrm{red}}^{*}(G)$ by (Theorem 11.8).

Example 19.13: Consider $G=\mathbb{Z}$. Then

$$
C^{*}(\mathbb{Z})=C^{*}\left(u_{n}, n \in \mathbb{Z} \mid u_{n} \text { unitary, } u_{n} u_{m}=u_{n+m}, u_{-m}=u_{m}^{*}\right)
$$

hence $u_{0}=1, u_{n}=\left(u_{1}\right)^{n}$ for $n>0$ and $u_{n}=\left(u_{1}^{*}\right)^{|n|}$ for $n<0$. Thus we have the equalities $C^{*}(\mathbb{Z})=C^{*}\left(u_{1}\right.$ unitary $)=C\left(\mathbb{S}^{1}\right)=C(\widehat{\mathbb{Z}})$. Because $\mathbb{Z}$ is abelian, $\mathbb{Z}$ is amenable, hence $C_{\text {max }}^{*}(\mathbb{Z})=C_{\text {red }}^{*}(\mathbb{Z})$.

Example 19.14: Consider $G=\mathbf{F}_{2}$, the free group on two generators $x$ and $y$. This $G$ is not amenable, thus $C_{\max }^{*}\left(\mathbf{F}_{2}\right) \neq C_{\text {red }}^{*}\left(\mathbf{F}_{2}\right)$. Also, we can write $C_{\max }^{*}\left(\mathbf{F}_{2}\right)$ as universal $C^{*}$ algebra: $C_{\max }^{*}\left(\mathbf{F}_{2}\right)=C^{*}(u, v$ unitaries).

These group $C^{*}$-algebras and the von Neumann algebras $L\left(\mathbf{F}_{n}\right)$ on $n$ generators are famous. Some properties of these are:

- $C^{*}\left(\mathbf{F}_{n}\right)$ and $C_{\text {red }}^{*}\left(\mathbf{F}_{n}\right)$ have faithful traces,
- $C^{*}\left(\mathbf{F}_{n}\right)$ and $C_{\text {red }}^{*}\left(\mathbf{F}_{n}\right)$ have no projections (a hot topic was to find such examples until $1981^{2}$ ).
- $C_{\mathrm{red}}^{*}\left(\mathbf{F}_{n}\right)$ is simple.

The problem for von Nemann algebras is the following: It is well known that $\mathbf{F}_{n} \neq \mathbf{F}_{m}$, if $n \neq m$. Furthermore it can be shown that for $n \neq m$, also $\mathbb{C F}_{n} \neq \mathbb{C F}_{m}$ and $C_{\text {red }}^{*}\left(\mathbf{F}_{n}\right) \not \not C_{\text {red }}^{*}\left(\mathbf{F}_{m}\right)$. But for the group von Neumann algebras, it is still unknown whether $L\left(\mathbf{F}_{n}\right) \nsubseteq L\left(\mathbf{F}_{m}\right)$ or $L\left(\mathbf{F}_{n}\right) \cong L\left(\mathbf{F}_{m}\right)$. This problem is called the free group factor problem and it has been open for more than 80 years now.

[^10]
## 20 Products of $C^{*}$-algebras

Given two $C^{*}$-algebras $A, B$. How can we form a "product" $A$ ? $B$ ?
Definition 20.1: Let $A, B$ be unital $C^{*}$-algebras and $C$ be a $C^{*}$-algebra, $j_{1}: C \hookrightarrow A$, $j_{2}: C \hookrightarrow B$ two embeddings.
(i) The $C^{*}$-algebra

$$
A *_{\mathbb{C}} B:=C^{*}\left(a \in B, b \in B \mid \text { relations of } A, \text { relations of } B, 1_{A}=1_{B}\right)
$$

is called the (unital) free product of $A$ and $B$,
(ii) The $C^{*}$-algebra
$A *_{C} B:=C^{*}\left(a \in A, b \in B \mid\right.$ relations of $A$, relations of $B, j_{1}(x)=j_{2}(x)$ for all $\left.x \in C\right\}$ is called the amalgamated free product of $A$ and $B$.

Proposition 20.2: The free product has the following universal property:


Example 20.3: We have

$$
C_{\max }^{*}\left(\mathbf{F}_{2}\right)=C\left(\mathbb{S}^{1}\right) * C\left(\mathbb{S}^{1}\right)=C^{*}(\mathbb{Z}) * C^{*}(\mathbb{Z})=C^{*}(\mathbb{Z} * \mathbb{Z})
$$

Remark 20.4: There are reduced free products.
Remark 20.5: Let $A, B$ be $C^{*}$-algebras. Then

$$
\begin{aligned}
& A \odot B:=\operatorname{span}\left\{\sum_{i=1}^{n} a_{i} \otimes b_{i}: a_{i} \in A, b_{i} \in B, n \in \mathbb{N}\right\} / \\
& \\
& \\
& \\
& \left(\left(x_{1}+x_{2}\right) \otimes y=x_{1} \otimes y+x_{2} \otimes y,\right. \\
& \\
& \\
& \\
& \\
& \left.\left.\lambda \otimes\left(y_{1}+y_{2}\right)=x \otimes y\right)=(\lambda x) \otimes y=x \otimes(\lambda y)\right)
\end{aligned}
$$

is the algebraic tensor product of $A$ and $B$.

Our concern shall now be to find a $C^{*}$-norm on $A \odot B$, of which there are in fact many.

Definition 20.6: Let $A, B$ be $C^{*}$-algebras. For $\sum_{i=1}^{n} a_{i} \otimes b_{i} \in A \odot B$ put

$$
\left\|\sum_{i=1}^{n} a_{i} \otimes b_{i}\right\|_{\max }:=\sup \left\{\|\pi(x)\| \mid \pi: A \odot B \longrightarrow B(H)^{*} \text {-homomorphism }\right\}
$$

then $A \otimes_{\max } B:=\mathrm{cl}_{\|\cdot\|_{\max }}(A \odot B)$ is called the maximal tensor product of $A$ and B;

$$
\left\|\sum_{i=1}^{n} a_{i} \otimes b_{i}\right\|_{\min }:=\left\|\sum_{i=1}^{n} \pi\left(a_{i}\right) \otimes \sigma\left(b_{i}\right)\right\|,
$$

where $\pi: A \hookrightarrow B(H), \sigma: B \hookrightarrow B(K)$ are faithful representations and $\pi \otimes \sigma: A \odot$ $B \rightarrow B(H \otimes K)$. Then $A \otimes_{\min } B:=\mathrm{cl}_{\|\cdot\|_{\text {min }}}(A \odot B)$ is called the minimal- or spatial tensor product.

Proposition 20.7: (i) We have a universal property:

(ii) For all $C^{*}$-norms $\|\cdot\|_{\alpha}$ on $A \odot B$ we have the diagram


Proposition 20.8: Let $A, B$ be unital $C^{*}$-algebras. Then
$A \otimes_{\max } B=C^{*}\left(a \in A\right.$, relations of $A, b \in B$ relations of $\left.B \mid a b=b a, 1_{A}=1_{B}\right)$.
Example 20.9: We have

$$
C^{*}\left(\mathbb{Z}^{2}\right)=C^{*}(\mathbb{Z} \times \mathbb{Z})=C^{*}(\mathbb{Z}) \otimes C^{*}(\mathbb{Z})=C\left(\mathbb{S}^{1}\right) \otimes C\left(\mathbb{S}^{1}\right)
$$

Remark 20.10: A $C^{*}$-algebra $A$ is called nuclear, if for all $C^{*}$-algebras $B$ it holds $A \otimes_{\min } B=A \otimes_{\max } B$.

If $A$ is commutative, then $A$ is commutative (for example, $\mathcal{T}, A_{\vartheta}$ and $\mathcal{O}_{n}$ are nuclear).
$G$ is amenable if and only if $C^{*}(G)$ is nuclear.

Remark 20.11: Recall: For $A, B C^{*}$-algebras,

$$
A \oplus B:=\{(a, b) \mid a \in A, b \in B\}, \quad\|(a, b)\|:=\max \{\|a\|,\|b\|\}
$$

is called the direct sum of the $C^{*}$-algebras $A$ and $B$. We can express $A \oplus B$ as universal $C^{*}$-algebra via

$$
A \oplus B=C^{*}(a \in A, \text { relations of } A, b \in B, \text { relations of } B \mid a b=0)
$$

Definition 20.12: Let $G$ be a locally compact group, $A$ a $C^{*}$-algebra and let $\alpha: G \rightarrow \operatorname{Aut}(A), g \mapsto \alpha_{g}$ be a continuous group homomorphism (i. e., $g \mapsto \alpha_{g}(x)$ is continuous for all $x \in A$ ). Then $\alpha$ is called an action of $G$ on $A, A$ is called a $G$ - $C^{*}$-algebra, the triple $(A, G, \alpha)$ is called a covariant system or $C^{*}$-dynamical system.

Definition 20.13: A covariant representation of a $C^{*}$-dynamical system ( $A, G, \alpha$ ) is a non-degenerate representation $\pi: A \rightarrow B(H)$ together with a unitary representation $G \rightarrow \mathfrak{U}(H), g \mapsto u_{g}$ (on the same Hilbert space!) such that for all $a \in A$, $g \in G$, it holds

$$
\pi\left(\alpha_{g}(x)\right)=u_{g} \pi(x) u_{g}^{-1}
$$

Remark 20.14: In the above situation, the unitaries $\left(u_{g}\right)_{g \in G}$ implement the automorphisms $\left(\alpha_{g}\right)_{g \in G}$, i. e., they make them inner.

Indeed, consider a $C^{*}$-algebra $A \subseteq B(H)$ and a unitary $u \in A$. Then

$$
A \longrightarrow A, \quad u \longmapsto u x u^{*}
$$

is a *-homomorphism $\left(x \mapsto u^{*} x u\right.$ is the inverse of $\left.x \mapsto u x u^{*}\right)$. It is an inner automorphism of $A$.
Remark 20.15: Let $G$ be discrete. Consider

$$
A G:=\left\{\sum_{g \in G}^{\mathrm{fin}} a_{g} \delta_{g}: a_{g} \in A\right\}
$$

similar to the group algebra but with coefficients from $A$ rather than from $\mathbb{C}$, with twisted multiplication

$$
\left(\delta_{t_{1}} * a \delta_{t_{2}}\right):=\alpha_{t_{1}}(a) \delta_{t_{1} t_{2}}, \quad \delta_{t}^{*}:=\delta_{t^{-1}}
$$

or rather

$$
\left(\sum a_{t} \delta_{t}\right)\left(\sum b_{s} \delta_{s}\right):=\sum_{s, t} a_{t} \alpha_{t}\left(b_{s}\right) \delta_{t s}
$$

How does this make $\left(\alpha_{t}\right)_{t \in G}$ inner? We have

$$
\left(\sum_{t} a_{t} \delta_{t}\right)\left(\sum_{s} b_{s} \delta_{s}\right)=\sum_{s, t} a_{t} \delta_{t} b_{s} \delta_{t}^{*} \delta_{t} \delta_{s}=\sum_{s, t} a_{t} \alpha_{t}\left(b_{s}\right) \delta_{s}
$$

This construction works analogously for locally compact groups using $L^{1}(G)$ instead of $\mathbb{C} G$.

Definition 20.16: Let $G$ be discrete. Put
$\|x\|:=\sup \{\|\pi(x)\| \mid \pi: A G \rightarrow B(H)$ non-degenerate
*-homomorphism induced from a covariant representation of $(A, G, \alpha)\}$
and call $A \rtimes_{\alpha} G:=\operatorname{cl}_{\|\cdot\|}(A G)$ the cross-product of $A$ with $G$ (given $\alpha$ ).
Remark 20.17: One can define left regular representations for $C^{*}$-dynamical systems and one does obtain a reduced version $A \rtimes_{r} G$.
Proposition 20.18: Let $G$ be discrete, $A$ be a unital $C^{*}$-algebra and $\alpha: G \rightarrow \operatorname{Aut}(A)$ be an action, i. e., let $(A, G, \alpha)$ be a $C^{*}$-dynamical system. Then

$$
\begin{aligned}
& A \rtimes_{\alpha} G=C^{*}\left(a \in A, \text { relations of } A, u_{g} \text { unitaries, } g \in G \mid\right. \\
& \left.\quad u_{g} u_{h}=u_{g h}, u_{g}^{*}=u_{g^{-1}}, \alpha_{g}(a)=u_{g} \text { aug } u_{g}^{*} \text { for all } g, h \in G, a \in A, 1_{A}=u_{e}\right)
\end{aligned}
$$

Remark 20.19: (i) $A \rtimes_{\alpha} G$ contains a copy of $A$, a copy of $G$ and it makes the automorphism $\alpha_{g}$ inner. It is the smallest $C^{*}$-algebra with this property. One can also say that $A$ is obtained by adjoining unitaries making the $\alpha_{g}$ inner.
Example 20.20: (i) Let $X$ be a locally compact space, $G$ be a locally compact group acting on $X$, i.e.,

$$
\alpha^{\circ}: X \times G \longrightarrow X, \quad(x, g) \longmapsto g x
$$

Then $\left(C_{0}(X), G, \alpha\right)$ is a $C^{*}$-dynamical system with

$$
\alpha: G \longrightarrow \operatorname{Aut}\left(C_{0}(X)\right), \quad \alpha_{g}(f)(x):=f\left(\alpha^{\circ}(x, g)\right)=f(g x)
$$

for $x \in X, g \in G$. This classical understanding of a dynamical system matches the notion of a $C^{*}$-dynamical system.
(ii) If $A=\mathbb{C}$, then $\alpha: G \rightarrow \operatorname{Aut}(G)$ acts trivially via $\alpha_{g}(\lambda)=\lambda$. Then for the cross-product $A \rtimes_{\alpha} G$ it holds: $A \rtimes_{\alpha} G=\mathbb{C} \rtimes_{\alpha} G=C^{*}(G)$.

More generally: If $G$ acts trivially of some $C^{*}$-algebra $A$ (i.e., $\alpha_{g} \equiv \mathrm{id}$ ), then $A \rtimes_{\alpha} G \cong A \otimes_{\max } C^{*}(G)$. We then also have $A \rtimes_{\alpha}^{r} G \cong A \otimes_{\min } C_{\text {red }}^{*}(G)$.
(iii) We have $C^{*}\left(N \rtimes_{\alpha} H\right)=C^{*}(N) \rtimes_{\alpha} H$, where on the left side we mean the semidirect product of groups and on the right side the cross-product of $C^{*}$-algebras.
(iv) For $G=\mathbb{Z}$, the action $\alpha: \mathbb{Z} \rightarrow \operatorname{Aut}(A)$ consists in one single automorphism $\alpha_{1}$, since $\alpha_{n}=\alpha_{1} \circ \cdots \circ \alpha_{1}$ and every automorphism $\alpha: A \rightarrow A$ induces an action

$$
\mathbb{Z} \longrightarrow \operatorname{Aut}(A), \quad u \longmapsto \alpha \circ \cdots \circ \alpha .
$$

Then

$$
A \rtimes_{\alpha} \mathbb{Z}=C^{*}\left(a \in A, \text { relations of } A, u \text { unitary } \mid u a u^{*}=\alpha(a)\right) .
$$

If for instance $A=C\left(\mathbb{S}^{1}\right)$ and

$$
\alpha: C\left(\mathbb{S}^{1}\right) \longrightarrow C\left(\mathbb{S}^{1}\right)=C^{*}(u \text { unitary }), \quad u \longmapsto e^{2 \pi \mathrm{i} \vartheta} u
$$

then $C\left(\mathbb{S}^{1}\right) \rtimes_{\alpha_{\vartheta}} \mathbb{Z}=A_{\vartheta}$ is the rotation algebra.


[^0]:    ${ }^{1}$ This is the so called theorem of Parseval, refer to Theorem 5.28 from the Functional Analysis I lecture notes.

[^1]:    ${ }^{2}$ For the notation, refer to the proof of Proposition 7.4 in the lecture notes of Functional Analysis I.

[^2]:    ${ }^{3}$ Check Functional Analysis I, Sheet 10, Exercise 3.

[^3]:    ${ }^{4}$ Check Functional Analysis I, Sheet 11, Exercise 3.

[^4]:    ${ }^{1}$ Note that the spectral radius is defined as $r(x):=\max \{|\lambda| \mid \lambda \in \operatorname{Sp}(x)\} \leq\|x\|$ and the spectrum $\operatorname{Sp}(x):=\{\lambda \in \mathbb{C} \mid \lambda 1-x$ not invertible in $B\} \subseteq \mathbb{C}$ is compact. Refer to Proposition 8.8 from Functional Analysis I.

[^5]:    ${ }^{2} \mathrm{We}$ call it the first fundamental theorem of $C^{*}$-algebras. This is non-standard naming.

[^6]:    ${ }^{1}$ Refer to the Functional Analysis I lecture notes, Chapter 2.

[^7]:    ${ }^{1}$ In the following we will denote with $(X, \mathfrak{T})^{*}$ the set

    $$
    \{f: X \rightarrow \mathbb{C} \mid f \text { continuous with respect to } \mathfrak{T}\} .
    $$

[^8]:    ${ }^{1}$ Here, $\lambda^{1}$ denotes the one-dimensional Lebesgue measure on $\mathfrak{B}([0,1])$.

[^9]:    ${ }^{1}$ We follow the convention $0 \cdot \infty:=0$.

[^10]:    ${ }^{1}$ This is only one of very many possible definitions for "amenable".
    ${ }^{2}$ Kaplansky started the interest in this question in 1958, when he first asked for an example of a simple $C^{*}$-algebra without projections.

