# Shannon's monotonicity problem for free and classical entropy 

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Recall that for a real valued random variable $X$ with density $p$, its entropy is defined as

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H(X)=-\int_{\mathbb{R}} p(x) \log p(x) d x
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Amongst random variables with $E(X)=0$ and $E\left(X^{2}\right)=1$, entropy is maximized by the standard Gaussian random variables $G$ with variance 1. Given a sequence $X_{1}, X_{2}, \ldots$ of independent, identically distributed random variables with $E\left(X_{n}\right)=0$ and $\operatorname{Var}\left(X_{n}\right)=1$ then the central limit states that their central limit sums

$$
z_{N}=\frac{x_{1}+\cdots+x_{N}}{\sqrt{N}}
$$

converge in law to $G$. Moreover, the entropy of this sequence is nondecreasing; a result due to Artstein, Ball, Barthe, and Naor [1].

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- Then $(A, \tau)$ is a non-commutative probability space, and any element $X \in A$ is a non-commutative random variable.
- $\tau(X)$ is the expectation, or first moment, and in general the law of $X$ refers to its moments $\left\{\tau\left(X^{n}\right): n \in \mathbb{N}\right\}$.
- Can think of the law of $X$ as a linear functional on polynomials $\mu_{X}: \mathbb{C}[t] \rightarrow \mathbb{C}$ so that $\mu_{X}(p(t))=\tau(p(X))$.


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- This example is more important to the non-commutative case than it first seems.

If $X \in A$ is self-adjoint, then there exists a measure $\mu$ supported on the spectrum of $X$ so that

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Moreover, given two such functions $f$ and $g, f(X) g(X)=(f \cdot g)(X)$. For the remainder of the talk we assume all non-commutative random variables are self-adjoint.

Given several non-commutative random variables $X_{1}, \ldots, X_{k}$, their joint law can be thought of as a linear functional on non-commutative polynomials $\mu_{X_{1}, \ldots, X_{k}}: \mathbb{C}\left\langle t_{1}, \ldots, t_{k}\right\rangle \rightarrow \mathbb{C}$ such that $\mu_{X_{1}, \ldots, X_{k}}\left(p\left(t_{1}, \ldots, t_{k}\right)\right)=\tau\left(p\left(X_{1}, \ldots, X_{k}\right)\right)$.

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Also, the law is no longer encoded by a measure.

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\left\langle e_{1} \otimes \cdots \otimes e_{k}, f_{1} \otimes \cdots \otimes f_{k}\right\rangle=\left\langle e_{1}, f_{1}\right\rangle \cdots\left\langle e_{k}, f_{k}\right\rangle .
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- The Fock space is defined as $\mathcal{F}(\mathcal{H})=\mathbb{C} \Omega \oplus \bigoplus_{k=1}^{\infty} \mathcal{H}^{\otimes k}$, where $\Omega$ is the vacuum vector (think "zero length tensor product"). Its inner product is the extension of the above where tensor products of different lengths are orthogonal.
- Fix a vector $e \in \mathcal{H}$ with $\|e\|=1$. The left creation operator is defined by

$$
\begin{aligned}
& I(e) \Omega=e \\
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- Define $c(e)=I(e)+I(e)^{*} \in \mathcal{B}(\mathcal{F}(\mathcal{H}))$.
- Define a linear functional $\tau: \mathcal{B}(\mathcal{F}(\mathcal{H})) \rightarrow \mathbb{C}$ by $\tau(X)=\langle\Omega, X \Omega\rangle$, then $\tau$ is a positive trace.
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- Consider the law of $c(e)$ with respect to $\tau$ :

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- Then

$$
\tau\left(c(e)^{n}\right)=\left\{\begin{array}{cl}
\frac{1}{k+1}\binom{2 k}{k} & \text { if } n=2 k \\
0 & \text { if } n=2 k+1
\end{array}\right.
$$

where $C_{k}=\frac{1}{k+1}\binom{2 k}{k}$ are the Catalan numbers.

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- Thus we say that $c(e)$ has the semicircle law or is semicircular.

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Consequently, any joint moment can be expressed as a product of their individual moments:

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In particular, if $W \in \mathbf{A l g}(1, X)$ and $Z \in \mathbf{A} \lg (1, Y)$ are two random variables so that $E(W)=E(Z)=0$, then $E(W Z)=0$.

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In particular, if $W \in \mathbf{A} \boldsymbol{\operatorname { l g }}(1, X)$ and $Z \in \mathbf{A l g}(1, Y)$ are two random variables so that $E(W)=E(Z)=0$, then $E(W Z)=0$.
Free independence captures this idea in the non-commutative case.

Let $F_{1}, F_{2} \subset(A, \tau)$ be two families of non-commutative random variables. Then we say that these families are freely independent if

$$
\tau\left(W_{1} W_{2} \cdots W_{n}\right)=0
$$

when $W_{j} \in \mathbf{A l g}\left(1, F_{i(j)}\right)$ are such that $\tau\left(W_{j}\right)=0$ and $i(j) \neq i(j+1)$, with $j=1, \ldots, n$ and $i(j) \in\{1,2\}$.

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More generally, $F_{1}, \ldots, F_{k} \subset(A, \tau)$ are freely independent if the above holds but now we simply take $i(j) \in\{1, \ldots, k\}$.

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For $i=1, \ldots, k$, let $X_{i}=\left(X_{i}^{(1)}, \ldots, X_{i}^{(p)}\right) \in A^{p}$ be a $p$-tuple of non-commutative random variables. Then we say these $p$-tuples are freely independent if the families $\left\{X_{1}^{(1)}, \ldots, X_{1}^{(p)}\right\}, \ldots,\left\{X_{k}^{(1)}, \ldots, X_{k}^{(p)}\right\}$ are freely independent.

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For example, to compute $\tau\left(X_{1} X_{2}\right)$. First write $X_{j}=\dot{\circ}_{j}+\tau\left(X_{j}\right) 1$, where $\dot{X}_{j}=X_{j}-\tau\left(X_{j}\right) 1$. Then $\dot{X}_{j} \in \operatorname{Alg}\left(1, X_{j}\right)$ and is centered. Thus we have

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\begin{aligned}
\tau\left(X_{1} X_{2}\right) & =\tau\left(\left(\stackrel{\circ}{X}_{1}+\tau\left(X_{1}\right) 1\right)\left(\stackrel{\circ}{X}_{2}+\tau\left(X_{2}\right) 1\right)\right) \\
& =\tau\left(\stackrel{\circ}{X}_{1} \stackrel{\circ}{X}_{2}\right)+\tau\left(\dot{X}_{1}\right) \tau\left(X_{2}\right)+\tau\left(X_{1}\right) \tau\left(\dot{X}_{2}\right)+\tau\left(X_{1}\right) \tau\left(X_{2}\right) \\
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& =\tau\left(X_{1}\right) \tau\left(X_{2}\right)
\end{aligned}
$$

where the first term vanishes because of free independence and the other two vanish because the $\check{X}_{j}$ are centered.

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\begin{aligned}
\tau\left(X_{1} X_{2}^{2} X_{1}\right) & =\tau\left(X_{1}^{2}\right) \tau\left(X_{2}^{2}\right) \\
\tau\left(X_{1} X_{2} X_{1} X_{2}\right) & =\tau\left(X_{1}^{2}\right) \tau\left(X_{2}\right)^{2}+\tau\left(X_{1}\right)^{2} \tau\left(X_{2}^{2}\right)-\tau\left(X_{1}\right)^{2} \tau\left(X_{2}\right)^{2}
\end{aligned}
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$$
\begin{aligned}
\tau\left(X_{1} X_{2}^{2} X_{1}\right) & =\tau\left(X_{1}^{2}\right) \tau\left(X_{2}^{2}\right) \\
\tau\left(X_{1} X_{2} X_{1} X_{2}\right) & =\tau\left(X_{1}^{2}\right) \tau\left(X_{2}\right)^{2}+\tau\left(X_{1}\right)^{2} \tau\left(X_{2}^{2}\right)-\tau\left(X_{1}\right)^{2} \tau\left(X_{2}\right)^{2}
\end{aligned}
$$

Given the above two examples, if $X_{1} X_{2}=X_{2} X_{1}$ then it would follow that

$$
\tau\left(\left(X_{1}-\tau\left(X_{1}\right) 1\right)^{2}\right) \tau\left(\left(X_{2}-\tau\left(X_{2}\right) 1\right)^{2}\right)=0
$$

i.e. the variance of $X_{1}$ or $X_{2}$ must vanish.

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- $\tau\left(I_{i}^{n} l_{i}^{* m}\right)=0$ unless $m=0=n$, hence a polynomial $p_{k} \in \mathbf{A l g}\left(1, l_{i}\right)$ has $\tau\left(p_{k}\right)=0$ iff it can be written as a sum of monomials (each having zero expectation) and no constant term.

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- $\tau\left(l_{i}^{n} l_{i}^{* m}\right)=0$ unless $m=0=n$, hence a polynomial $p_{k} \in \operatorname{Alg}\left(1, l_{i}\right)$ has $\tau\left(p_{k}\right)=0$ iff it can be written as a sum of monomials (each having zero expectation) and no constant term.
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- Considering $\left\langle\Omega, p_{1} \cdots p_{r} \Omega\right\rangle$, it is easy to see $m_{k} \neq 0$ for any $k$ implies this is zero, but $m_{k}=0$ for all $k$ implies $n_{k}>0$ for all $k$ and hence $p_{1} \cdots p_{r} \Omega \in \mathcal{H}^{\otimes\left(n_{1}+\cdots+n_{r}\right)} \perp \mathbb{C} \Omega$.

With the notion of free independence, we can state one of the first parallels to the classical case:

## Theorem 1 (Free central limit theorem, [6])

Let $X_{1}, X_{2}, \ldots$ be a sequence of freely independent random variables in some non-commutative probability space $(A, \tau)$. Assume $\tau\left(X_{n}\right)=0$ and $\tau\left(X_{n}^{2}\right)=1$ for all $n$, and that $\sup _{n}\left|\tau\left(X_{n}^{p}\right)\right|<\infty$ for all $p$. Then the laws of the sequence

$$
Z_{N}=\frac{1}{\sqrt{N}}\left(X_{1}+\cdots+X_{N}\right)
$$

converge in moments to the semicircle law $d \mu=\frac{1}{2 \pi} \sqrt{4-t^{2}} d t$.

## Using a standard construction in operator algebras (the

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For $X, Y \in A$ let

$$
\langle X, Y\rangle_{L^{2}(A, \tau)}=\langle X, Y\rangle_{2}:=\tau\left(X^{*} Y\right)
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This defines a sesquilinear form on the complex vector space $A$, which is complex linear in the second coordinate.

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To correct this we consider the set $N=\left\{X \in A:\langle X, X\rangle_{2}=0\right\}$. We want to mod out by $N$, but in order for $A / N$ to still be a vector space we need $N$ to be a vector subspace.

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To correct this we consider the set $N=\left\{X \in A:\langle X, X\rangle_{2}=0\right\}$. We want to mod out by $N$, but in order for $A / N$ to still be a vector space we need $N$ to be a vector subspace.
This follows from the fact that $N=\left\{X \in A:\langle Y, X\rangle_{2}=0 \forall Y \in A\right\}$. Now $A / N$ is a vector space on which $\langle\cdot, \cdot\rangle_{2}$ is an inner product. Let $L^{2}(A, \tau)$ be the Hilbert space obtained by taking the completion of $A / N$ with respect to the norm induced by $\langle\cdot, \cdot\rangle_{2}$.

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\pi(X) \hat{Y}=\widehat{X Y}
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The positivity of $\tau$ allows us to define the non-commutative $L^{p}$ spaces:

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L^{p}(A, \tau)=\left\{X \in A: \tau\left(\left(X^{*} X\right)^{p}\right)<\infty\right\}, \quad\|X\|_{p}=\tau\left(\left(X^{*} X\right)^{p}\right)^{\frac{1}{p}}
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$$

The von Neumann algebra generated by $A$ is then what plays the role of the non-commutative $L^{\infty}$ space:

$$
W^{*}(A)=\overline{\pi(A)}^{S O T}=\overline{\pi(A)}^{W O T}=\pi(A)^{\prime \prime} \cap \mathcal{B}\left(L^{2}(A, \tau)\right)
$$

More generally, given a family $F \subset A$ we let

$$
W^{*}(F)={\overline{\pi\left(\operatorname{Alg}_{*}(1, F)\right)}}^{S O T} \subset W^{*}(A)
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Let that $L^{2}\left(W^{*}(F), \tau\right)=\overline{W^{*}(F) \cdot \hat{1}^{\|\cdot\|_{2}} \subset L^{2}(A, \tau) \text {. Then we define the }}$ orthogonal projection $E_{W^{*}(F)}: L^{2}(A, \tau) \rightarrow L^{2}\left(W^{*}(F), \tau\right)$ onto this subspace.

Given non-commutative random variables $X_{1}, \ldots, X_{n} \in(A, \tau)$ we define Voiculescu's free difference quotients

$$
\partial_{X_{j}: X_{1}, \ldots, \hat{X}_{j}, \ldots, X_{n}}=\partial_{j}: \mathbf{A} \lg \left(1, X_{1}, \ldots, X_{n}\right) \rightarrow L^{2}(A, \tau) \bar{\otimes} L^{2}(A, \tau)
$$

by $\partial_{j}\left(X_{k}\right)=\delta_{j=k} 1 \otimes 1$ and the Leibniz rule:

$$
\partial_{j}(W Z)=\partial_{j}(W) \cdot Z+W \cdot \partial_{j}(Z)
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$$

For example,

$$
\begin{aligned}
\partial_{2}\left(X_{1} X_{2} X_{3}\right) & =\partial_{2}\left(X_{1}\right) \cdot\left(X_{2} X_{3}\right)+X_{1} \cdot \partial_{2}\left(X_{2} X_{3}\right) \\
& =0+X_{1} \cdot\left[\partial_{2}\left(X_{2}\right) \cdot X_{3}+X_{2} \cdot \partial_{2}\left(X_{3}\right)\right] \\
& =X_{1} \cdot(1 \otimes 1) \cdot X_{3}+0=X_{1} \otimes X_{3}
\end{aligned}
$$

If we think of $\partial_{j}$ as a map on $L^{2}\left(W^{*}\left(X_{1}, \ldots, X_{n}\right), \tau\right) \subset L^{2}(A, \tau)$, then we can consider its adjoint $\partial_{j}^{*}: L^{2}(A, \tau) \bar{\otimes} L^{2}(A, \tau) \rightarrow L^{2}(A, \tau)$.

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\xi_{j}=J\left(X_{j}: X_{1}, \ldots, \hat{X}_{j}, \ldots, X_{n}\right)=\partial_{j}^{*}(1 \otimes 1)
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That is, if $Y \in L^{2}\left(W^{*}\left(X_{1}, \ldots, X_{n}\right), \tau\right)$ then

$$
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$$

The free Fisher information of the $n$-tuple $\left(X_{1}, \ldots, X_{n}\right)$ is then defined as

$$
\Phi^{*}\left(X_{1}, \ldots, X_{n}\right):=\sum_{j=1}^{n}\left\|\xi_{j}\right\|_{L^{2}(A, \tau)}^{2}
$$

The free entropy of the $n$-tuple $\left(X_{1}, \ldots, X_{n}\right)$ is defined as

$$
\chi^{*}\left(X_{1}, \ldots, X_{n}\right)=\frac{1}{2} \int_{0}^{\infty}\left[\frac{n}{1+t}-\Phi^{*}\left(X_{1}^{t}, \ldots, X_{n}^{t}\right)\right] d t+\frac{n}{2} \log 2 \pi e,
$$

where $X_{j}^{t}=X_{j}+\sqrt{t} S_{j}$ and $S_{1}, \ldots, S_{n}$ are freely independent, identically distributed, centered, semicircular variables of variance 1 , which are also freely independent from $X_{1}, \ldots, X_{n}$.

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When $n=1$, the above definition is equivalent to

$$
\chi(X)=\iint_{\mathbb{R}^{2}} \log |s-t| d \mu_{X}(s) d \mu_{X}(t)+\frac{3}{4}+\frac{1}{2} \log 2 \pi
$$

where $\mu_{X}$ is the law of $X$.

## Theorem 2

Let $(A, \tau)$ be a non-commutative probability space. Let $X_{j}=\left(X_{j}^{(1)}, \ldots, X_{j}^{(p)}\right) \in A^{p}, j=1,2, \ldots$ be a sequence of $p$-tuples of random variables, such that $X_{1}, X_{2}, \ldots$ are freely independent, identically distributed, and have finite second moments. Define $Z_{N}=N^{-1 / 2}\left(X_{1}+\cdots+X_{N}\right)$. Then the function $N \mapsto \chi^{*}\left(Z_{N}^{(1)}, \ldots, Z_{N}^{(p)}\right)$ is monotone nondecreasing.

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To show this we first show that the free Fisher information is monotone nonincreasing.

## Lemma 3

Assume $Z \in(A, \tau)$ is freely independent from $X, Y_{1}, \ldots, Y_{n} \in(A, \tau)$. Then $J\left(X: Y_{1}, \ldots, Y_{n}\right)$ exists iff $J\left(X: Y_{1}, \ldots, Y_{n}, Z\right)$ exists, in which case we have

$$
J\left(X: Y_{1}, \ldots, Y_{n}\right)=J\left(X: Y_{1}, \ldots, Y_{n}, Z\right)
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## Proof.

Let $\partial_{0}=\partial_{X:} Y_{1}, \ldots, Y_{n}$ and $\partial_{1}=\partial_{X:} Y_{1}, \ldots, Y_{n}, Z$. Then we can think of these as maps on $B_{0}:=W^{*}\left(X, Y_{1}, \ldots, Y_{n}\right)$ and $B_{1}:=W^{*}\left(X, Y_{1}, \ldots, Y_{n}, Z\right)$.

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$$
\xi_{0}:=J\left(X: Y_{1}, \ldots, Y_{n}\right)=E_{0}\left(\xi_{1}\right)
$$

## Proof of Lemma 3 (cont.)

Thus it suffices to show that if $\xi_{0}$ exists, then $\xi_{1}$ exists and $\xi_{0}=\xi_{1}$.

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Thus it suffices to show that if $\xi_{0}$ exists, then $\xi_{1}$ exists and $\xi_{0}=\xi_{1}$. Consider an arbitrary element of $\operatorname{Alg}\left(1, X, Y_{1}, \ldots, Y_{n}, Z\right)$ :

$$
R=Q_{0} P_{1} Q_{1} \cdots P_{r} Q_{r}
$$

with $P_{k} \in \mathbf{A} \lg \left(1, X, Y_{1}, \ldots, Y_{n}\right), Q_{k} \in \mathbf{A l g}(1, Z)$.

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We must show $\left\langle\xi_{0}, R\right\rangle_{2}=\left\langle 1 \otimes 1, \partial_{1}(R)\right\rangle_{L^{2}(A, \tau) \bar{\otimes} L^{2}(A, \tau)}$. We proceed by induction on $r$.

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For $r=0$, note that

$$
\tau\left(\xi_{0}\right)=\left\langle\xi_{0}, 1\right\rangle_{2}=\left\langle 1 \otimes 1, \partial_{0}(1)\right\rangle_{L^{2}(A, \tau) \bar{\otimes} L^{2}(A, \tau)}=0
$$

and so by free independence...

## Proof of Lemma 3 (cont.)

$$
\begin{aligned}
\left\langle\xi_{0}, Q_{0}\right\rangle_{2} & =\tau\left(\xi_{0} Q_{0}\right) \\
& =\tau\left(\xi_{0}\left(Q_{0}-\tau\left(Q_{0}\right)\right)\right)+\tau\left(\xi_{0}\right) \tau\left(Q_{0}\right)=0
\end{aligned}
$$

On the other hand, $\partial_{1}\left(Q_{0}\right)=0$, so the base case holds.

## Proof of Lemma 3 (cont.)

$$
\begin{aligned}
\left\langle\xi_{0}, Q_{0}\right\rangle_{2} & =\tau\left(\xi_{0} Q_{0}\right) \\
& =\tau\left(\xi_{0}\left(Q_{0}-\tau\left(Q_{0}\right)\right)\right)+\tau\left(\xi_{0}\right) \tau\left(Q_{0}\right)=0
\end{aligned}
$$

On the other hand, $\partial_{1}\left(Q_{0}\right)=0$, so the base case holds. For $r>0$, we note that we can assume $P_{1}, Q_{1}, \ldots, P_{r-1}, Q_{r-1}, P_{r}$ are centered.

## Proof of Lemma 3 (cont.)

$$
\begin{aligned}
\left\langle\xi_{0}, Q_{0}\right\rangle_{2} & =\tau\left(\xi_{0} Q_{0}\right) \\
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On the other hand, $\partial_{1}\left(Q_{0}\right)=0$, so the base case holds. For $r>0$, we note that we can assume $P_{1}, Q_{1}, \ldots, P_{r-1}, Q_{r-1}, P_{r}$ are centered.
Indeed, by expanding each term in $P_{1} Q_{1} \cdots Q_{r-1} P_{r}$ into its centered part and scalar part, the resulting products have either all centered terms or are covered by the induction hypothesis.

## Proof of Lemma 3 (cont.)

Thus, for $r=1$ we have

$$
\begin{aligned}
\tau\left(\xi_{0} Q_{0} P_{1} Q_{1}\right)= & \tau\left(\xi_{0} \stackrel{\circ}{0}_{0} P_{1} \grave{Q}_{1}\right)+\tau\left(Q_{0}\right) \tau\left(\xi_{0} P_{1} \grave{Q}_{1}\right) \\
& +\tau\left(Q_{1}\right) \tau\left(\xi_{0} \grave{Q}_{0} P_{1}\right)+\tau\left(Q_{0}\right) \tau\left(Q_{1}\right) \tau\left(\xi_{0} P_{1}\right) \\
= & \tau\left(Q_{0}\right) \tau\left(Q_{1}\right) \tau\left(\xi_{0} P_{1}\right) \\
= & \tau\left(Q_{0}\right) \tau\left(Q_{1}\right)\left\langle 1 \otimes 1, \partial_{0} P_{1}\right\rangle_{L^{2}(A, \tau) \bar{\otimes} L^{2}(A, \tau)},
\end{aligned}
$$

## Proof of Lemma 3 (cont.)

Thus, for $r=1$ we have

$$
\begin{aligned}
\tau\left(\xi_{0} Q_{0} P_{1} Q_{1}\right)= & \tau\left(\xi_{0} \check{Q}_{0} P_{1} \grave{Q}_{1}\right)+\tau\left(Q_{0}\right) \tau\left(\xi_{0} P_{1} \grave{Q}_{1}\right) \\
& +\tau\left(Q_{1}\right) \tau\left(\xi_{0} \grave{Q}_{0} P_{1}\right)+\tau\left(Q_{0}\right) \tau\left(Q_{1}\right) \tau\left(\xi_{0} P_{1}\right) \\
= & \tau\left(Q_{0}\right) \tau\left(Q_{1}\right) \tau\left(\xi_{0} P_{1}\right) \\
= & \tau\left(Q_{0}\right) \tau\left(Q_{1}\right)\left\langle 1 \otimes 1, \partial_{0} P_{1}\right\rangle_{L^{2}(A, \tau) \bar{\otimes} L^{2}(A, \tau)}
\end{aligned}
$$

while on the other hand

$$
\begin{aligned}
\left\langle\xi_{1}, Q_{0} P_{1} Q_{1}\right\rangle_{2} & =\tau \otimes \tau\left(\partial_{1}\left(Q_{0} P_{1} Q_{1}\right)\right) \\
& =\tau \otimes \tau\left(Q_{0} \cdot \partial_{1}\left(P_{1}\right) \cdot Q_{1}\right) \\
& =\tau \otimes \tau\left(Q_{0} \cdot \partial_{0}\left(P_{1}\right) \cdot Q_{1}\right) \\
& =\tau\left(Q_{0}\right) \tau\left(Q_{1}\right) \tau \otimes \tau\left(\partial_{0} P_{1}\right)
\end{aligned}
$$

where the last equality follows from free independence.

## Proof of Lemma 3 (cont.)

For $r \geq 2$ we have

$$
\begin{aligned}
& \tau\left(\xi_{0} Q_{0} P_{1} \ldots P_{r} Q_{r}\right) \\
& =\tau\left(\xi_{0} \stackrel{\circ}{Q}_{0} P_{1} \cdots P_{r} \stackrel{\circ}{Q}_{r}\right)+\tau\left(Q_{0}\right) \tau\left(\xi_{0} P_{1} \ldots P_{r} \stackrel{\circ}{Q}_{r}\right) \\
& \quad+\tau\left(Q_{r}\right) \tau\left(\xi_{0} \stackrel{\circ}{0}_{0} P_{1} \cdots P_{r}\right)+\tau\left(Q_{0}\right) \tau\left(Q_{r}\right) \tau\left(\xi_{0} P_{1} \ldots P_{r}\right)=0 .
\end{aligned}
$$

And

$$
\begin{aligned}
& \tau \otimes \tau\left(\partial_{1}\left(Q_{0} P_{1} \cdots P_{r} Q_{r}\right)\right) \\
& \quad=\sum_{l=1}^{r} \tau \otimes \tau\left(\left[Q_{0} P_{1} \cdots Q_{l-1}\right] \cdot \partial_{1}\left(P_{l}\right) \cdot\left[Q_{l} \cdots P_{r} Q_{r}\right]\right)=0 .
\end{aligned}
$$

## Lemma 4

Let $\left\{X_{j}^{(k)}\right\} \subset(A, \tau), k=1, \ldots, p, j=1,2, \ldots$ be non-commutative random variables. Fix $N \in \mathbb{N}, j=1, \ldots, N+1$, and $k=1, \ldots, p$. Then one has

$$
\begin{aligned}
& J\left(\sum_{i=1}^{N+1} X_{i}^{(k)}:\left\{\sum_{i=1}^{N+1} X_{i}^{(r)}\right\}_{r \neq k}\right) \\
& =E_{W^{*}\left(\left\{\sum_{i=1}^{N+1} X_{i}^{(r)}\right\}_{r=1}^{p}\right)^{J}\left(\sum_{i \neq j} X_{i}^{(k)}:\left\{\sum_{i \neq j} X_{i}^{(r)}\right\}_{r \neq k},\left\{X_{j}^{(r)}\right\}_{r=1}^{p}\right)}
\end{aligned}
$$

assuming the conjugate variables on the right-hand side exists.

## Lemma 4

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\end{aligned}
$$

assuming the conjugate variables on the right-hand side exists.
The basic idea, is that if $y=\tilde{y}+x$ then $\partial_{y}(p(y))=\partial_{\tilde{y}}(p(\tilde{y}+x))$.

## Proof.

Let $Y_{k}=\sum_{i=1}^{N+1} X_{i}^{(k)}$ and $Y_{k}^{\prime}=\sum_{i \neq j} X_{i}^{(k)}$, so that $Y_{k}=Y_{k}^{\prime}+X_{j}^{(k)}$. Then a polynomial $P$ in $Y_{1}, \ldots, Y_{p}$ can be viewed as a polynomial in $Y_{1}^{\prime}, \ldots, Y_{p}^{\prime}, X_{j}^{(1)}, \ldots, X_{j}^{(p)}$.

## Proof.

Let $Y_{k}=\sum_{i=1}^{N+1} X_{i}^{(k)}$ and $Y_{k}^{\prime}=\sum_{i \neq j} X_{i}^{(k)}$, so that $Y_{k}=Y_{k}^{\prime}+X_{j}^{(k)}$. Then a polynomial $P$ in $Y_{1}, \ldots, Y_{p}$ can be viewed as a polynomial in $Y_{1}^{\prime}, \ldots, Y_{p}^{\prime}, X_{j}^{(1)}, \ldots, X_{j}^{(p)}$.
In particular,

$$
\partial_{Y_{k}^{\prime}:\left\{Y_{r}^{\prime}: r \neq k\right\},\left\{X_{j}^{(r)}: r=1, \ldots, p\right\}} P=\partial_{Y_{k}:\left\{Y_{r}: r \neq k\right\}} P,
$$

since the derivation is determined by the Leibniz rule and the values

$$
\partial_{Y_{k}^{\prime}:\left\{Y_{t}^{\prime}: r \neq k\right\},\left\{X_{j}^{(r)}: r=1, \ldots, p\right\}}\left(Y_{q}^{\prime}+X_{j}^{(q)}\right)=\partial_{Y_{k}:\left\{Y_{r}: r \neq k\right\}}\left(Y_{q}\right)=\delta_{k=q} 1 \otimes 1 .
$$

## Proof of Lemma 4 (cont.)

Hence

$$
\begin{aligned}
\left\langle P, J\left(Y_{k}^{\prime}:\left\{Y_{r}^{\prime}: r\right.\right.\right. & \left.\left.\neq k\},\left\{X_{j}^{(r)}: r=1, \ldots, p\right\}\right)\right\rangle_{2} \\
& =\left\langle P, J\left(Y_{k}:\left\{Y_{r}: r \neq k\right\}\right)\right\rangle_{2},
\end{aligned}
$$

which concludes the proof as $P \in W^{*}\left(Y_{1}, \ldots, Y_{p}\right)$.

## Theorem 5

Let $(A, \tau)$ be a non-commutative probability space. Let $X_{j}=\left(X_{j}^{(1)}, \ldots, X_{j}^{(p)}\right) \in A^{p}, j=1,2, \ldots$ be a sequence of $p$-tuples of random variables, such that $X_{1}, X_{2}, \ldots$ are freely independent, identically distributed, and have finite second moments. Define $Z_{N}=N^{-1 / 2}\left(X_{1}+\cdots+X_{N}\right)$. Then the function $N \mapsto \Phi^{*}\left(Z_{N}^{(1)}, \ldots, Z_{N}^{(p)}\right)$ is monotone nonincreasing.

## Theorem 5

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## Proof.

First note that

$$
J\left(c X: b Y_{1}, \ldots, b Y_{n}\right)=c^{-1} J\left(X: Y_{1}, \ldots, Y_{n}\right)
$$

this follows from the observation $\partial_{c t}(t)=c^{-1} \partial_{c t}(c t)=c^{-1} \partial_{t}(t)$.

## Proof of Theorem 5 (cont.)

Fix $j$ and let $Y_{k}$ and $Y_{k}^{\prime}$ be as in the previous proof. Then $Z_{N+1}^{(r)}=(N+1)^{-1 / 2} Y_{r}$. Let $B=W^{*}\left(Z_{N+1}^{(1)}, \ldots, Z_{N+1}^{(p)}\right)$.

## Proof of Theorem 5 (cont.)

Fix $j$ and let $Y_{k}$ and $Y_{k}^{\prime}$ be as in the previous proof. Then
$Z_{N+1}^{(r)}=(N+1)^{-1 / 2} Y_{r}$. Let $B=W^{*}\left(Z_{N+1}^{(1)}, \ldots, Z_{N+1}^{(p)}\right)$.
Then

$$
\begin{aligned}
J\left(Z_{N+1}^{(k)}:\left\{Z_{N+1}^{(r)}: r \neq k\right\}\right) & =(N+1)^{\frac{1}{2}} J\left(Y_{k}:\left\{Y_{r}: r \neq k\right\}\right) \\
& =(N+1)^{\frac{1}{2}} E_{B} J\left(Y_{k}^{\prime}:\left\{Y_{r}^{\prime}: r \neq k\right\},\left\{X_{j}^{(r)}\right\}_{r=1}^{n}\right) \\
& =(N+1)^{\frac{1}{2}} E_{B} J\left(Y_{k}^{\prime}:\left\{Y_{r}^{\prime}: r \neq k\right\}\right),
\end{aligned}
$$

where we have used (in order) our initial observation, Lemma 4, and Lemma 3. (Recall that the $X_{j}^{(r)}$ are freely independent from the $Y_{r}^{\prime}=\sum_{i \neq j} X_{i}^{(r)}$.)

## Proof of Theorem 5 (cont.)

Thus we have

$$
\begin{aligned}
(N+1)^{\frac{1}{2}} J & \left(Z_{N+1}^{(k)}:\left\{Z_{N+1}^{(r)}: r \neq k\right\}\right) \\
& =E_{B}(N+1) J\left(\sum_{i \neq j} X_{i}^{(k)}:\left\{\sum_{i \neq j} X_{i}^{(r)}\right\}_{r \neq k}\right) \\
& =E_{B} \sum_{j=1}^{N+1} J\left(\sum_{i \neq j} X_{i}^{(k)}:\left\{\sum_{i \neq j} X_{i}^{(r)}\right\}_{r \neq k}\right)
\end{aligned}
$$

where we have used the fact that our initial choice of $j$ was arbitrary.

## Proof of Theorem 5 (cont.)

Since $E_{B}$ is a contraction on $L^{2}(A, \tau)$ we then have

$$
\begin{equation*}
\left\|J\left(Z_{N+1}^{(k)}:\left\{Z_{N+1}^{(r)}: r \neq k\right\}\right)\right\|_{2}^{2} \leq(N+1)^{-1}\left\|\sum_{j=1}^{N+1} \zeta_{j}\right\|_{2}^{2} \tag{1}
\end{equation*}
$$

where

$$
\zeta_{j}=J\left(\sum_{i \neq j} X_{i}^{(k)}:\left\{\sum_{i \neq j} X_{i}^{(r)}\right\}_{r \neq k}\right)
$$

To proceed, we need to appeal to a lemma from the classical proof:

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## Lemma 6 ([1])

Let $E_{1}, \ldots, E_{N+1}$ be commuting orthogonal projections in a Hilbert space. Assume that we have $N+1$ vectors $\zeta_{1}, \ldots, \zeta_{N+1}$ such that for every $j$, $E_{1} \cdots E_{N+1} \zeta_{j}=0$. Then

$$
\left\|\sum_{j=1}^{N+1} E_{j} \zeta_{j}\right\|^{2} \leq N \sum_{j=1}^{N+1}\left\|\zeta_{j}\right\|^{2} .
$$

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$$
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$$

## Proof of Theorem 5 (cont.)

We let $M_{j}=W^{*}\left(X_{j}\right), M=W^{*}\left(X_{1}, \ldots, X_{N+1}\right), Q_{j}=W^{*}\left(X_{i}: i \neq j\right)$, and $E_{j}=E_{Q_{j}}$.

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## Proof of Theorem 5 (cont.)

We let $M_{j}=W^{*}\left(X_{j}\right), M=W^{*}\left(X_{1}, \ldots, X_{N+1}\right), Q_{j}=W^{*}\left(X_{i}: i \neq j\right)$, and $E_{j}=E_{Q_{j}}$.
We claim the $E_{j}$ commute, $E_{j} \zeta_{j}=\zeta_{j}$, and $E_{1} \cdots E_{N+1} \zeta_{j}=\tau\left(\zeta_{j}\right)=0$.

## Proof of Theorem 5 (cont.)

That $E_{j} \zeta_{j}=\zeta_{j}$ follows from the definition of $\zeta_{j}$, and $\tau\left(\zeta_{j}\right)=\left\langle 1, \zeta_{j}\right\rangle_{2}=\left\langle\partial_{j}(1), 1 \otimes 1\right\rangle=0$.

## Proof of Theorem 5 (cont.)

That $E_{j} \zeta_{j}=\zeta_{j}$ follows from the definition of $\zeta_{j}$, and $\tau\left(\zeta_{j}\right)=\left\langle 1, \zeta_{j}\right\rangle_{2}=\left\langle\partial_{j}(1), 1 \otimes 1\right\rangle=0$.
Let $\stackrel{\circ}{M}_{j}=M_{j} \ominus \mathbb{C} 1$ (i.e. the centered elements in $M_{j}$ ). Then $M$ has the following orthogonal decomposition:

$$
L^{2}(M, \tau)=\mathbb{C} 1 \oplus \bigoplus_{n=1}^{\infty}\left[\bigoplus_{j_{1} \neq \cdots \neq j_{n}} \check{M}_{j_{1}} \check{M}_{j_{2}} \cdots \check{M}_{j_{n}}\right] .
$$

It is orthogonal precisely because of the free independence.

## Proof of Theorem 5 (cont.)

That $E_{j} \zeta_{j}=\zeta_{j}$ follows from the definition of $\zeta_{j}$, and $\tau\left(\zeta_{j}\right)=\left\langle 1, \zeta_{j}\right\rangle_{2}=\left\langle\partial_{j}(1), 1 \otimes 1\right\rangle=0$.
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$$

It is orthogonal precisely because of the free independence.
$Q_{j}$ has the same decomposition except that $j_{k}$ is never allowed to be $j$, and $E_{j}$ is determined by $E_{j} 1=1$ and

$$
\left.E_{j}\right|_{\dot{M}_{j_{1}} \dot{M}_{j_{2}} \cdots \dot{M}_{j_{n}}}=\left\{\begin{array}{ll}
\text { id } & \text { if } j \notin\left\{j_{1}, \ldots, j_{n}\right\} \\
0 & \text { otherwise }
\end{array} .\right.
$$

## Proof of Theorem 5 (cont.)

From this characterization is clear that the $E_{j}$ commute with one another and

$$
\left.E_{j} E_{i}\right|_{\dot{M}_{j_{1}} \dot{M}_{j_{2}} \cdots \dot{M}_{j_{n}}}=\left\{\begin{array}{cl}
\text { id } & \text { if }\{i, j\} \cap\left\{j_{1}, \ldots, j_{n}\right\}=\emptyset \\
0 & \text { otherwise }
\end{array} .\right.
$$

## Proof of Theorem 5 (cont.)

From this characterization is clear that the $E_{j}$ commute with one another and

$$
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0 & \text { otherwise }
\end{array} .\right.
$$

Moreover, $E_{1} \cdots E_{N+1}$ is the orthogonal projection onto the scalars $\mathbb{C} 1$. Hence we can determine $E_{1} \cdots E_{N+1} \zeta_{j}$ by considering the inner products of the $\zeta_{j}$ against scalars:

$$
\left\langle 1, \zeta_{j}\right\rangle_{2}=\tau\left(\zeta_{j}\right),
$$

so $E_{1} \cdots E_{N+1} \zeta_{j}=\tau\left(\zeta_{j}\right) 1$, as claimed.

## Proof of Theorem 5 (cont.)

Applying Lemma 6 to (1) yields

$$
\left\|J\left(Z_{N+1}^{(k)}:\left\{Z_{N+1}^{(r)}: r \neq k\right\}\right)\right\|_{2}^{2} \leq(N+1)^{-1} N \sum_{j=1}^{N+1}\left\|\zeta_{j}\right\|_{2}^{2} .
$$

## Proof of Theorem 5 (cont.)

Applying Lemma 6 to (1) yields

$$
\left\|J\left(Z_{N+1}^{(k)}:\left\{Z_{N+1}^{(r)}: r \neq k\right\}\right)\right\|_{2}^{2} \leq(N+1)^{-1} N \sum_{j=1}^{N+1}\left\|\zeta_{j}\right\|_{2}^{2} .
$$

However, since the $p$-tuples are identically distributed we have

$$
\begin{aligned}
\sum_{j=1}^{N+1}\left\|\zeta_{j}\right\|_{2}^{2} & =(N+1)\left\|\zeta_{N+1}\right\|_{2}^{2} \\
& =(N+1)\left\|J\left(\sum_{i=1}^{N} X_{i}^{(k)}:\left\{\sum_{i=1}^{N} X_{i}^{(r)}: r \neq k\right\}\right)\right\|_{2}^{2} \\
& =\frac{N+1}{N} \| J\left(Z_{N}^{(k)}:\left\{Z_{N}^{(r)}: r \neq k\right\} \|_{2}^{2}\right.
\end{aligned}
$$

## Proof of Theorem 5 (cont.)

Combining the two previous equations then yields

$$
\left\|J\left(Z_{N+1}^{(k)}:\left\{Z_{N+1}^{(r)}: r \neq k\right\}\right)\right\|_{2}^{2} \leq \| J\left(Z_{N}^{(k)}:\left\{Z_{N}^{(r)}: r \neq k\right\} \|_{2}^{2}\right.
$$

## Proof of Theorem 5 (cont.)

Combining the two previous equations then yields

$$
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$$

Finally, summing over $k$ yields

$$
\Phi^{*}\left(Z_{N+1}^{(1)}, \ldots, Z_{N+1}^{(p)}\right) \leq \Phi^{*}\left(Z_{N}^{(1)}, \ldots, Z_{N}^{(p)}\right)
$$

## Proof of Theorem 2.

We wish to show

$$
\chi^{*}\left(Z_{N}^{(1)}, \ldots, Z_{N}^{(p)}\right) \leq \chi^{*}\left(Z_{N+1}^{(1)}, \ldots, Z_{N+1}^{(p)}\right)
$$

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$$

where for an arbitrary p-tuple $\left(W^{(1)}, \ldots, W^{(p)}\right)$

$$
\begin{aligned}
& \chi^{*}\left(W^{(1)}, \ldots, W^{(p)}\right) \\
& \quad=\frac{1}{2} \int_{0}^{\infty}\left[\frac{p}{1+t}-\Phi^{*}\left(W^{(1, t)}, \ldots, W^{(p, t)}\right)\right] d t+\frac{p}{2} \log 2 \pi e
\end{aligned}
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\end{aligned}
$$

where $W^{(k, t)}=W^{(k)}+\sqrt{t} S^{(k)}$ with $S^{(1)}, \ldots, S^{(p)}$ a freely iid centered semicircular random variables of variance 1 , freely independent from $W^{(1)}, \ldots, W^{(p)}$.

## Proof of Theorem 2 (cont.)

Let $\left\{S_{j}^{(k)}: j=1, \ldots, N+1, k=1, \ldots, p\right\}$ be freely iid centered semicircular variables of variance 1 , which are freely independent from $\left\{X_{j}^{(k)}\right\}_{j, k}$.

## Proof of Theorem 2 (cont.)

Let $\left\{S_{j}^{(k)}: j=1, \ldots, N+1, k=1, \ldots, p\right\}$ be freely iid centered semicircular variables of variance 1 , which are freely independent from $\left\{X_{j}^{(k)}\right\}_{j, k}$.
Define $X_{j}^{(k, t)}=X_{j}^{(k)}+\sqrt{t} S_{j}^{(k)}$ and $Z_{N}^{(k, t)}=N^{-1 / 2}\left(X_{1}^{(k, t)}+\cdots+X_{N}^{(k, t)}\right)$.

## Proof of Theorem 2 (cont.)

Let $\left\{S_{j}^{(k)}: j=1, \ldots, N+1, k=1, \ldots, p\right\}$ be freely iid centered semicircular variables of variance 1 , which are freely independent from $\left\{X_{j}^{(k)}\right\}_{j, k}$.
Define $X_{j}^{(k, t)}=X_{j}^{(k)}+\sqrt{t} S_{j}^{(k)}$ and $Z_{N}^{(k, t)}=N^{-1 / 2}\left(X_{1}^{(k, t)}+\cdots+X_{N}^{(k, t)}\right)$.
The $p$-tuples $X_{j}^{t}=\left(X_{j}^{(1, t)}, \ldots, X_{j}^{(p, t)}\right)$ are freely iid with finite second moments (via the Cauchy-Schwarz inequality).

## Proof of Theorem 2 (cont.)

Let $\left\{S_{j}^{(k)}: j=1, \ldots, N+1, k=1, \ldots, p\right\}$ be freely iid centered semicircular variables of variance 1 , which are freely independent from $\left\{X_{j}^{(k)}\right\}_{j, k}$.
Define $X_{j}^{(k, t)}=X_{j}^{(k)}+\sqrt{t} S_{j}^{(k)}$ and $Z_{N}^{(k, t)}=N^{-1 / 2}\left(X_{1}^{(k, t)}+\cdots+X_{N}^{(k, t)}\right)$.
The $p$-tuples $X_{j}^{t}=\left(X_{j}^{(1, t)}, \ldots, X_{j}^{(p, t)}\right)$ are freely iid with finite second moments (via the Cauchy-Schwarz inequality).
Hence we may apply Theorem 5 to obtain

$$
\Phi^{*}\left(Z_{N}^{(1, t)}, \ldots, Z_{N}^{(p, t)}\right) \geq \Phi^{*}\left(Z_{N+1}^{(1, t)}, \ldots, Z_{N+1}^{(p, t)}\right)
$$

## Proof of Theorem 2 (cont.)

Note that $Z_{N}^{(k, t)}=Z_{N}^{(k)}+\sqrt{t} S^{(N, k)}$ where for each fixed $N$, $S^{(N, k)}=N^{-1 / 2}\left(S_{1}^{(k)}+\cdots+S_{N}^{(k)}\right), k=1, \ldots, p$ is a family of centered freely iid semicircular variables freely independent from $\left\{Z_{N}^{(k)}\right\}_{k=1}^{p}$ and having variance 1.

## Proof of Theorem 2 (cont.)

Note that $Z_{N}^{(k, t)}=Z_{N}^{(k)}+\sqrt{t} S^{(N, k)}$ where for each fixed $N$, $S^{(N, k)}=N^{-1 / 2}\left(S_{1}^{(k)}+\cdots+S_{N}^{(k)}\right), k=1, \ldots, p$ is a family of centered freely iid semicircular variables freely independent from $\left\{Z_{N}^{(k)}\right\}_{k=1}^{p}$ and having variance 1.
The definition of $\chi^{*}$ and the free Fisher information inequality gives

$$
\begin{aligned}
\chi^{*} & \left(Z_{N}^{(1)}, \ldots, Z_{N}^{(p)}\right) \\
& =\frac{1}{2} \int_{0}^{\infty}\left[\frac{p}{1+t}-\Phi^{*}\left(Z_{N}^{(1, t)}, \ldots, Z_{N}^{(p, t)}\right)\right] d t+\frac{p}{2} \log 2 \pi e \\
& \leq\left[\frac{p}{1+t}-\Phi^{*}\left(Z_{N+1}^{(1, t)}, \ldots, Z_{N+1}^{(p, t)}\right)\right] d t+\frac{p}{2} \log 2 \pi e \\
& =\chi^{*}\left(Z_{N+1}^{(1)}, \ldots, Z_{N+1}^{(p)}\right)
\end{aligned}
$$

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