

Threshold phenomena for quantum marginals

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Case Western Reserve/Paris 6

Cambridge, October 15, 2013

Collaborators: G. Aubrun, D. Ye

Comm. Pure Appl. Math. (2014), *arXiv:1106.2264v3*

Phys. Rev. A. 85(R) (2012), *arXiv:1112.4582v2*

<http://www.cwru.edu/artsci/math/szarek/>

Abstract

Consider a quantum system consisting of N identical particles and assume that it is in a random pure state (i.e., uniformly distributed over the sphere of the corresponding Hilbert space) and two subsystems \mathcal{A} and \mathcal{B} consisting of k particles each. Are \mathcal{A} and \mathcal{B} likely to share entanglement? Is the \mathcal{AB} -marginal typically PPT?

For many natural properties there is a sharp “phase transition.” E.g., there is a threshold $K \sim N/5$ such that - if $k > K$, then \mathcal{A} and \mathcal{B} typically share entanglement - if $k < K$, then \mathcal{A} and \mathcal{B} typically do not share entanglement.

The first statement was (essentially) shown in the talk by G. Aubrun. Here we present a general scheme for handling such questions and sketch the analysis specific to entanglement.

Talk summary

- setup, notation; random quantum states and ensembles of random matrices
- emergence of entanglement (the main result)
- a sketch of the proof using the tools of geometric functional analysis and random matrix theory

Sets of states

The set of (mixed) states on \mathcal{H} is denoted by $\mathcal{D} = \mathcal{D}(\mathcal{H})$

Separable (unentangled) states on $\mathcal{H} = \mathcal{H}_1 \otimes \mathcal{H}_2$: $\mathcal{S} = \mathcal{S}(\mathcal{H})$

When we talk about $\mathcal{S}(\mathbb{C}^n)$, we implicitly assume $n = d^2$ and $\mathbb{C}^n \sim \mathbb{C}^d \otimes \mathbb{C}^d$ (a fixed bipartition).

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The sets $\mathcal{D}(\mathbb{C}^n)$ and $\mathcal{S}(\mathbb{C}^n)$ are convex bodies in the hyperplane $H_1 := \{\text{tr}(\cdot) = 1\} \subset \mathcal{M}_n^{\text{sa}}$ with $\frac{1}{n}\mathbb{I}$ in the interior. $\frac{1}{n}\mathbb{I}$ is the only point invariant under the isometries of either set.

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Our approach will (in principle) work for any property in place of $d \times d$ separability, provided it has some minimal permanence properties.

Setup, notation; partial trace

The state of the entire system is described by

$$|\psi\rangle \in \mathcal{H} = \mathcal{A} \otimes \mathcal{B} \otimes \mathcal{E} = (\mathbb{C}^D)^{\otimes k} \otimes (\mathbb{C}^D)^{\otimes k} \otimes (\mathbb{C}^D)^{\otimes N-2k} \\ = \mathbb{C}^d \otimes \mathbb{C}^d \otimes \mathbb{C}^s$$

$$\text{with } d = D^k, s = D^{N-2k}$$

$$\text{If } N = 5k, \text{ then } s = D^{N-2k} = D^{3k} = d^3$$

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If the entire system is in the pure state $|\psi\rangle \in \mathcal{H}$ or $|\psi\rangle\langle\psi|$ in the density matrix formalism, then the \mathcal{AB} -marginal is given by the partial trace $\text{tr}_{\mathcal{E}} |\psi\rangle\langle\psi| = \text{tr}_{\mathbb{C}^s} |\psi\rangle\langle\psi|$.

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If $|\psi\rangle \in \mathbb{C}^d \otimes \mathbb{C}^d \otimes \mathbb{C}^s \sim \mathbb{C}^n \otimes \mathbb{C}^s$ is identified with a matrix $A \in \mathcal{M}_{n \times s}$, then $\rho = AA^\dagger$.

Random marginals and “standard” ensembles of random matrices

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Entanglement emergence for $\rho_{n,s}$, $\mathbb{C}^n \sim \mathbb{C}^d \otimes \mathbb{C}^d$

Theorem There is a sharp “entanglement emergence” threshold $s_{ent} = s_{ent}(d)$ for $\rho_{n,s}$ verifying $cd^3 \leq s_{ent} \leq Cd^3 \log^2 d$

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More precisely, if $\varepsilon > 0$, then

If $s \leq (1 - \varepsilon)s_{ent}(d)$, then $\mathbb{P}(\rho_{n,s} \text{ is separable}) \leq 2 \exp(-c(\varepsilon)s_{ent})$

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Earlier results in various directions/cases were in
Kendon–Życzkowski–Munro 2002, Hayden–Leung–Winter 2006,
Aubrun–S. 2006, Ye 2010

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$$d = D^k \text{ with } d_1 = D^{k-1} = \frac{d}{D} \leq \frac{d}{2} \quad \text{and}$$

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Problem Is $\pi_{d,s}$ an increasing function of s ?

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$$\rho \text{ is unentangled} \iff \rho - \frac{\mathbb{I}}{n} \in \mathcal{S}_0 \iff \left\| \rho - \frac{\mathbb{I}}{n} \right\|_{\mathcal{S}_0} \leq 1$$

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Need to show that, for appropriate values of n, s and for $\rho = \rho_{n,s}$, the above occur with probability close to 1 or 0.

The strategy : concentration

Step 1: Show that, for the appropriate values of n, s and $\varepsilon > 0$, $\mathbb{E}\|\rho_{n,s} - \frac{\mathbb{I}}{n}\|_{\mathcal{S}_0} \leq 1 - \varepsilon$ or $\mathbb{E}\|\rho_{n,s} - \frac{\mathbb{I}}{n}\|_{\mathcal{S}_0} \geq 1 + \varepsilon$, as needed.

Step 2: Show that $A \rightarrow f(A) := \|AA^\dagger - \frac{\mathbb{I}}{n}\|_{\mathcal{S}_0}$ is smooth enough and so it concentrates around its median (and its mean).

Recall that A varies over the Frobenius sphere in $\mathcal{M}_{n \times s}$; the relevant metric will also given by the Frobenius norm $\|\cdot\|_2$.

Step 2 is nontrivial, but relatively routine.

Step 1 is harder and requires a few new tricks.

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The latter is a Lipschitz function with constant r^{-1} , where r is the (Euclidean/Frobenius) **inradius of \mathcal{S}_0** , which is known to be the same as the **inradius of \mathcal{D}** (Gurvits-Barnum 2002), which is $1/\sqrt{n(n-1)} \sim 1/n$.

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So the Lipschitz constant is $\sim n$, in fact $< n$.

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We can now appeal to **Levy's lemma**:

Let $f : S^{m-1} \rightarrow \mathbb{R}$ be an L -Lipschitz function and let $\varepsilon > 0$. Then

$$\mathbb{P}(\{|f - M_f| > \varepsilon\}) \leq C \exp\left(-\frac{m}{2L^2} \varepsilon^2\right),$$

where M_f is the median (or mean) of f , \mathbb{P} is the normalized uniform measure on the sphere and $C > 0$ is a universal constant.

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Verification of the exponent:

$$m = 2ns, L \sim 2n, \text{ so } \frac{m}{2L^2} \sim \frac{2ns}{8n^2} = \frac{s}{4n}.$$

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Not quite the $\Omega(s)$ that we wanted, but since ultimately we are interested in the range $s = \Omega(d^3) = \Omega(n^{3/2})$, this yields some concentration for $\varepsilon = o(1)$ and so is “marginally sufficient” (even if not optimal).

A better argument

Let $\mathcal{T} := \{A : \|A\|_\infty = O(\frac{1}{\sqrt{n}})\}$. Then

- $\mathbb{P}(\mathcal{T}) \geq 1 - e^{-cs}$
- $A \rightarrow AA^\dagger$ is “locally Lipschitz on \mathcal{T} ” with constant $O(\frac{1}{\sqrt{n}})$.

The bottom line is that the local Lipschitz constant of f on \mathcal{T} is $O(\frac{1}{\sqrt{n}} \times n) = O(\sqrt{n})$.

This follows from AA^\dagger being a rescaled Wishart distribution, whose eigenvalue distribution approximates (for large n, s) the rescaled **Marchenko-Pastur** distribution. In particular, the singular values of A are typically in the interval $\frac{1}{\sqrt{n}} \times [1 - \sqrt{\frac{n}{s}}, 1 + \sqrt{\frac{n}{s}}]$, hence $\|A\|_\infty = O(\frac{1}{\sqrt{n}})$. The probability estimate stated above is a consequence of the corresponding **large deviation** bound.

Local Levy's lemma

Let $\mathcal{T} \subset S^{m-1}$ be a subset of measure larger than $3/4$. Let $f : S^{m-1} \rightarrow \mathbb{R}$ be a function such that the restriction of f to \mathcal{T} is L -Lipschitz. Then, for every $\varepsilon > 0$,

$$\mathbb{P}(\{|f - M_f| > \varepsilon\}) \leq \mathbb{P}(S^{m-1} \setminus \mathcal{T}) + C \exp\left(-\frac{m}{2L^2} \varepsilon^2\right),$$

where M_f is the median of f and $C > 0$ a universal constant.

Recalculation of the exponent:

$$m = 2ns, \quad L = O(\sqrt{n}), \quad \text{so } \frac{m}{2L^2} = \frac{ns}{L^2} = \Omega(s).$$

We had $\exp(-c(\varepsilon)s)$ in the Theorem, so this is about right.

Step 1, estimating $\mathbb{E} \|\rho_{n,S} - \frac{\mathbb{I}}{n}\|_{\mathcal{S}_0}$: *the difficulty*

The gauge $\|\cdot\|_{\mathcal{S}_0}$ is hard to work with directly (NP-hard to calculate; Gurvits 2003).

Yesterday's argument was based on working with the "dual picture" and the somewhat easier dual quantity, $\|\cdot\|_{\mathcal{S}_0^\circ}$, where for $K \subset \mathbb{R}^m$

$$K^\circ := \{x \in \mathbb{R}^m : \langle x, y \rangle \leq 1 \text{ for all } y \in K\},$$

and subsequently on estimating $w(\mathcal{S}_0) := \int_{S_{H_0}} \|u\|_{\mathcal{S}_0^\circ} du$.

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Substep (i) [Random matrices] : **if** $n, \frac{s}{n} \rightarrow \infty$, **then**

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Thus the mean (and the median) of $\mathbb{E}\|\rho_{n,s} - \frac{\mathbb{I}}{n}\|_{\mathcal{S}_0}$ are $\frac{\varepsilon}{4}$ -separated from 1 and so, by concentration, $\mathbb{P}(\|\rho_{n,s} - \frac{\mathbb{I}}{n}\|_{\mathcal{S}_0} \leq 1) \approx 1$.

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If we had the MM^* -estimate for \mathcal{S}_0 , it would follow that $w(\mathcal{S}_0^\circ) = O(d^{3/2} \log d)$, hence $s_{ent} = w(\mathcal{S}_0^\circ)^2 = O(d^3 \log^2 d)$, exactly as needed.

The workaround

One guarantee for $K \subset \mathbb{R}^m$ to be in the ℓ -position is when the isometry group of K acts irreducibly on \mathbb{R}^m .

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Specifically, the action of the isometry group (local unitaries) splits into three irreducible factors, one of dimension $(d^2 - 1)^2$ and two of dimension $d^2 - 1$. It then follows from general theory that to bring \mathcal{S}_0 to the ℓ -position we only need to apply some dilations in the two smaller factors, and since their dimensions are relatively small, this does not affect in a major way the mean width of \mathcal{S}_0 or its polar.

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Oops... need central symmetry for the MM^* -estimate and \mathcal{S}_0 is not symmetric... there is another workaround based on Santaló, inverse Santaló, Urysohn, Rogers-Shephard inequalities ...

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The trick : linearization

We approximate $\rho_{n,s} - \frac{\mathbb{I}}{n}$ by $\frac{1}{n\sqrt{s}} GUE_0$, where GUE_0 is the standard H_0 -valued random Gaussian matrix, and use the relationship between the spherical and the Gaussian mean

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The bottom line is then

$$\mathbb{E} \|\rho_{n,s} - \frac{\mathbb{I}}{n}\|_{\mathcal{S}_0} \sim \frac{1}{n\sqrt{s}} \|GUE_0\|_{\mathcal{S}_0} \sim \frac{\sqrt{n^2-1}}{n\sqrt{s}} \int_{S_{H_0}} \|u\|_{\mathcal{S}_0} du$$

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from Marchenko-Pastur, while – by Wigner – the spectrum of $\frac{1}{\sqrt{n}} GUE_0$ lives on $[-2, 2]$, so at least the scaling is right:

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The needed equivalence $\mathbb{E} \|\rho_{n,s} - \frac{\mathbb{I}}{n}\|_{\mathcal{S}_0} \sim \mathbb{E} \|\frac{1}{n\sqrt{s}} GUE_0\|_{\mathcal{S}_0}$ will now follow from the general theory. However, since the gauge $\|\cdot\|_{\mathcal{S}_0}$ is rather poorly continuous, some finesse is needed.

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Proposition Let μ be a 0-centered non-point mass measure.

Then $\forall \varepsilon > 0 \exists \eta > 0$ such that if

- $x, y \in \mathbb{R}^{n,0}$ verify $d_\infty(\nu_x, \mu) \leq \eta$ and $d_\infty(\nu_y, \mu) \leq \eta$, and
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We now apply the Proposition with $\phi = \phi_{\mathcal{S}_0}$, μ – the semicircular distribution, x, y – the spectra of the two ensembles, and take the expected values. In addition to the convergence of the ensembles to a common limit spectral distribution, the calculation uses their invariance under conjugation by unitaries.

The concept behind the Proposition: *majorization*

For $x, y \in \mathbb{R}^{n,0}$, we write $x \prec y$ if, for every $k \in \{1, \dots, n\}$,

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Note that the assertion of the Proposition is in the spirit of (a), while $\frac{1}{n} \sum_{i=1}^n |x_i - t|$ from (b) can be rewritten as $\int \Phi d\nu_x$, where $\Phi(u) := |u - t|$, and so can be related to convergence of measures.

THANK YOU