# Mathematical Studies on Quantum Systems and Locally Quantum Systems 

Yuuya Yoshida

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#### Abstract

General probabilistic theories (GPTs) are theoretical models having states, measurements, and probabilities of obtaining measurement outcomes. For example, classical probability theory and quantum theory are GPTs with probability vectors and density matrices as states, respectively. In the study of GPTs, one often describes a finite number of GPTs $A_{1}, \ldots, A_{n}$ together, as one GPT $A=A_{1} \cdots A_{n}$. In this case, we say that $A$ is a whole system of subsystems $A_{1}, \ldots, A_{n}$. If $A_{1}, \ldots, A_{n}$ are quantum systems, then $A$ is called a locally quantum system. A quantum system is locally quantum, but there are many other locally quantum systems. In this thesis, we study three topics associated with quantum systems and locally quantum systems.

The first one is capacity. For a GPT, the maximum number of simultaneously and perfectly distinguishable states is called the capacity. It is known that the capacity does not really depend on the GPT. In this thesis, we propose a statement S to determine the capacities of special locally quantum systems, and prove a weaker statement WS. Statement WS is a statement associated with tensor product spaces over $\mathbb{C}$. Our proof works even if replacing the scalar field $\mathbb{C}$ with an arbitrary infinite field $\mathcal{F}$. We also show a partial result for finite fields.

The second one is differential privacy (DP). DP was born in the study to utilize private data while protecting the data (privacy-preserving data mining, PPDM). For this reason, DP has been studied by using classical probability theory, and has almost never been studied in quantum information theory. In this thesis, we define a quantum version of DP (called classical-quantum DP) and investigate its mathematical aspects. We define the set $\mathrm{CQ}_{n}$ of all classical-quantum differentially private $n$-tuples of quantum states and the subset $\mathrm{EC}_{n}$ of $\mathrm{CQ}_{n}$ that is essentially classical. We prove that $\mathrm{EC}_{n}=\mathrm{CQ}_{n}$ if $n=2$ and $\mathrm{EC}_{n} \neq \mathrm{CQ}_{n}$ if $n \geq 3$.

The third one is perfect discrimination of two states. In quantum theory, orthogonality is a necessary and sufficient condition to distinguish two states perfectly (denote this equivalence by E). However, equivalence E does not hold for a general GPT. In this thesis, we construct continuous one-parameter families of locally quantum systems in a certain natural manner, and investigate perfect discrimination of two states in the framework of those locally quantum systems. These one-parameter families contain a locally quantum system sufficiently close to a quantum system. Nevertheless, we show that they violate equivalence E.


## Publication list

This thesis is organized with special emphasis on the author's single-authored articles, and is based on the following papers and preprint.

1. Y. Yoshida, Maximum dimension of subspaces with no product basis, Linear Algebra Appl. 620, 228-241 (2021).
2. Y. Yoshida, Classical-quantum differentially private mechanisms beyond classical ones, arXiv:2011.09960 (2020).
3. Y. Yoshida and M. Hayashi, Classical mechanism is optimal in classical-quantum differentially private mechanisms, Proc. of 2020 IEEE International Symposium on Information Theory (ISIT), pp. 1991-1995 (2020). ${ }^{1}$
4. Y. Yoshida, H. Arai, and M. Hayashi, Perfect discrimination in approximate quantum theory of general probabilistic theories, Phys. Rev. Lett. 125, 150402 (2020) ${ }^{2}$

The contents of this thesis have the following correspondence. My contribution to paper 4 is its mathematical parts: Definition 2, Theorems 1, 2, and their proofs. Definition 1 in paper 4 is a definition discussed together with Hayato Arai.

| Section 1.1 | Section 1 of paper 1 |
| :---: | :---: |
| Section 1.2 | Section 1 of preprint 2 |
| Section 1.3 | None |
| Chapter 2 | Section 4 of paper 1 |
| Chapter 3 | Sections 2 and 3 of paper 1 |
| Chapter $\overline{4}$ excluding Section 4.2 | Sections 2-4 and Appendix of preprint 2 |
| Section 4.2 | Paper 3 |
| Chapter 5 | Paper 4 |
| Section 6.1 | None |
| Section $\overline{6.2}$ | Section 5 of preprint 2 |
| Section 6.3 | None |

I also wrote the following papers, which are not included in this thesis.
5. K. Saito and Y. Yoshida, Distributions of finite sequences represented by polynomials in Piatetski-Shapiro sequences, J. Number Theory 222, 115-156(2021).

[^0]6. Y. Yoshida and M. Hayashi, Asymptotic properties for Markovian dynamics in quantum theory and general probabilistic theories, J. Phys. A 53, 215303 (2020).
7. H. Arai, Y. Yoshida, and M. Hayashi, Perfect discrimination of non-orthogonal separable pure states on bipartite system in general probabilistic theory, $J$. Phys. A 52, 465304 (2019).

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## Contents

1 Introduction ..... 1
1.1 Statement determining the capacities of GPTs, and weaker statement ..... 3
1.2 Differential privacy ..... 6
1.3 Perfect discrimination of two states ..... 10
2 GPTs and their capacities ..... 13
2.1 Framework of GPTs ..... 13
2.2 Capacities of locally quantum systems ..... 15
3 Weaker statement: Maximum dimension of subspaces with no prod- uct basis ..... 17
3.1 More general proposition and proof ..... 17
3.2 Case of finite fields ..... 22
4 Classical-quantum differential privacy ..... 25
4.1 Linear mappings in quantum information theory ..... 25
4.2 Case $n=2$ ..... 26
4.3 Case $n \geq 3$ : Main idea and result ..... 28
4.4 Case $n \geq 3$ : Proof ..... 29
4.5 Concrete CQ $\varepsilon$-DP $n$-tuples that do not lie in $\mathrm{EC}_{n}(\varepsilon)$ ..... 35
5 Perfect discrimination of two states in approximate quantum theory ..... 37
5.1 One-parameter families of proper cones ..... 37
5.2 Proofs ..... 39
6 Conclusion ..... 45
6.1 Statement determining the capacities of GPTs, and weaker statement ..... 45
6.2 Differential privacy ..... 46
6.3 Perfect discrimination of two states ..... 46

## Chapter 1

## Introduction

General probabilistic theories.-A general probabilistic theory (GPT) is a theoretical model $[36,40,41,43,48,54]$, which is defined as a triplet $(\mathcal{V}, \mathcal{K}, u)$ of a real Hilbert space space $\mathcal{V}$, a convex cone $\mathcal{K} \subset \mathcal{V}$ and an element $u \in \mathcal{V}$ (more precisely, see Section 2.1). For example, the GPTs $\left(\mathbb{R}^{d},[0, \infty)^{d}, \mathbf{1}_{d}\right)$ and $\left(\operatorname{Herm}(d), \operatorname{PSD}(d), I_{d}\right)$ are classical probability theory and quantum theory, respectively, where $\mathbf{1}_{d}$ denotes the vector $[1, \ldots, 1]^{\top} \in \mathbb{R}^{d} ; \operatorname{Herm}(d)$ denotes the set of all Hermitian matrices on $\mathbb{C}^{d} ; \operatorname{PSD}(d)$ denotes the set of all positive semi-definite matrices on $\mathbb{C}^{d} ; I_{d}$ denotes the identity matrix on $\mathbb{C}^{d}$. The Hilbert spaces $\mathbb{R}^{d}$ and $\operatorname{Herm}(d)$ are equipped with the inner products $\sum_{i=1}^{d} x_{i} y_{i}$ and $\operatorname{Tr} X Y$, respectively, where $\operatorname{Tr} X$ denotes the (nonnormalized) trace of a square matrix $X$. See also Table 1.1 which summarizes the definitions of classical probability theory and quantum theory as GPTs. Each GPT $(\mathcal{V}, \mathcal{K}, u)$ has the state class $\mathcal{S}(\mathcal{K}, u)$ and measurement class $\mathcal{M}\left(\mathcal{K}^{*}, u\right)$ defined as

$$
\begin{align*}
\mathcal{S}(\mathcal{K}, u) & =\{x \in \mathcal{K}:\langle x, u\rangle=1\} \\
\mathcal{M}\left(\mathcal{K}^{*}, u\right) & =\left\{\left(y_{i}\right)_{i=1}^{m} m \text {-tuple of elements in } \mathcal{K}^{*}: m \in \mathbb{N}, \sum_{i=1}^{m} y_{i}=u\right\} \tag{1.1}
\end{align*}
$$

where $\mathcal{K}^{*}$ denotes the dual cone of $\mathcal{K}$ :

$$
\mathcal{K}^{*}=\{y \in \mathcal{V}: \forall x \in \mathcal{K},\langle y, x\rangle \geq 0\} .
$$

An element in the state class (resp. measurement class) is called a state (resp. measurement). For example, a state in classical probability theory (resp. quantum theory) is a probability vector (resp. density matrix). Although a measurement in classical probability theory has no particular name, a measurement in quantum theory is a positive-operator-valued measure (POVM).

Table 1.1: Classical probability theory and quantum theory as GPTs.

| GPT | Hilbert space $\mathcal{V}$ | Inner product $\langle\cdot, \cdot\rangle$ | Convex cone $\mathcal{K}$ |
| :---: | :---: | :---: | :---: |
| Classical | $\mathbb{R}^{d}$ | $\sum_{i=1}^{d} x(i) y(i)$ | $[0, \infty)^{d}$ |
| Quantum | Herm $(d)$ | $\operatorname{Tr} X Y$ | $\operatorname{PSD}(d)$ |
| GPT | Element $u$ | State $x$ | Measurement $\left(y_{i}\right)_{i=1}^{m}$ |
| Classical | $\mathbf{1}_{d}$ | Probability vector $p$ | No name |
| Quantum | $I_{d}$ | Density matrix $\rho$ | POVM $\left(M_{i}\right)_{i=1}^{m}$ |

Given a state $x \in \mathcal{S}(\mathcal{K}, u)$ and a measurement $\left(y_{i}\right)_{i=1}^{m} \in \mathcal{M}\left(\mathcal{K}^{*}, u\right)$, the probability of obtaining each outcome $i=1, \ldots, m$ is $\left\langle x, y_{i}\right\rangle$. Note that $\left(\left\langle x, y_{i}\right\rangle\right)_{i=1}^{m}$ is a probability vector. For instance, a dice roll is described by classical probability theory, and its state is $x=(1 / 6) \mathbf{1}_{6}$. If we want to know whether the number on a cast dice is even or odd, then we should use the measurement consisting of $y_{\text {even }}=e_{2}+e_{4}+e_{6}$ and $y_{\text {odd }}=e_{1}+e_{3}+e_{5}$, where $\left(e_{i}\right)_{i=1}^{6}$ denotes the standard basis of $\mathbb{R}^{6}$. The probability for the number to be even (resp. odd) is $\left\langle x, y_{\text {even }}\right\rangle$ (resp. $\left\langle x, y_{\text {odd }}\right\rangle$ ).

As stated above, classical probability theory and quantum theory are typical GPTs, but there are many other GPTs. A simple example is the PR box [35] (named after Popescu and Rohrlich $[47])$. The convex cone $\mathcal{K}$ of the PR box $\left(\mathbb{R}^{3}, \mathcal{K}, e_{3}\right)$ is a square pyramid, and is different from $[0, \infty)^{d}$ and $\operatorname{PSD}(d)$ which appear in classical probability theory and quantum theory, respectively. ${ }^{1}$

Whole systems and subsystems.-Consider classical probability theory. In classical probability theory, one often considers the joint distribution $\mathbb{P}_{X_{1} \cdots X_{n}}$ of random variables $X_{1}, \ldots, X_{n}$ and its marginal distributions $\mathbb{P}_{X_{1}}, \ldots, \mathbb{P}_{X_{n}}$. The marginal distributions come from the joint distribution, but in general, one cannot recover the joint distribution from the marginal distributions.

The study of GPTs also has an analog of the above situation. We say that a $\operatorname{GPT}(\mathcal{V}, \mathcal{K}, u)$ is a whole system of subsystems $\left(\mathcal{V}^{[i]}, \mathcal{K}^{[i]}, u^{[i]}\right), i=1, \ldots, n$, if the $n$ $\operatorname{GPTs}\left(\mathcal{V}^{[i]}, \mathcal{K}^{[i]}, u^{[i]}\right)$ satisfy (i) $\mathcal{V}=\mathcal{V}^{[1]} \otimes \cdots \otimes \mathcal{V}^{[n]}$, (ii) $\mathcal{K}_{\text {min }} \subset \mathcal{K} \subset \mathcal{K}_{\text {max }}$, and (iii) $u=u^{[1]} \otimes \cdots \otimes u^{[n]}$, where $\mathcal{K}_{\text {min }}$ and $\mathcal{K}_{\text {max }}$ are defined as

$$
\begin{align*}
& \mathcal{K}_{\text {min }}=\operatorname{conv}\left\{x^{[1]} \otimes \cdots \otimes x^{[n]}: x^{[1]} \in \mathcal{K}^{[1]}, \ldots, x^{[n]} \in \mathcal{K}^{[n]}\right\}, \\
& \mathcal{K}_{\text {max }}=\left\{x \in \mathcal{V}: \forall y^{[1]} \in\left(\mathcal{K}^{[1]}\right)^{*}, \ldots, \forall y^{[n]} \in\left(\mathcal{K}^{[n]}\right)^{*},\left\langle x, y^{[1]} \otimes \cdots \otimes y^{[n]}\right\rangle \geq 0\right\} . \tag{1.2}
\end{align*}
$$

That is, an element in $\mathcal{K}_{\text {min }}$ is a convex combination of tensor products of $x^{[1]} \in$ $\mathcal{K}^{[1]}, \ldots, x^{[n]} \in \mathcal{K}^{[n]}$, and an element in $\mathcal{K}_{\text {max }}^{*}$ is a convex combination of tensor prod-

[^1]ucts of $y^{[1]} \in\left(\mathcal{K}^{[1]}\right)^{*}, \ldots, y^{[n]} \in\left(\mathcal{K}^{[n]}\right)^{*} \cdot 2$ It is the most important here that the whole system (namely, its convex cone $\mathcal{K}$ ) is not uniquely determined even if its subsystems are given, since $\mathcal{K}_{\text {min }}$ and $\mathcal{K}_{\text {max }}$ are not equal to each other in general. However, it follows that $\mathcal{K}_{\text {min }}=\mathcal{K}_{\text {max }}$ if all subsystems are classical (i.e., described by classical probability theory). In this case, the whole system (namely, convex cone $\mathcal{K}$ ) is uniquely determined and classical.

However, the quantum case is different from the classical case. We say that a $\operatorname{GPT}(\mathcal{V}, \mathcal{K}, u)$ is a locally quantum system if $(\mathcal{V}, \mathcal{K}, u)$ is a whole system of quantum subsystems (i.e., subsystems are described by quantum theory). When ( $\mathcal{V}, \mathcal{K}, u)$ is a locally quantum system, we simply write $\mathcal{S}(\mathcal{K}, u)$ and $\mathcal{M}\left(\mathcal{K}^{*}, u\right)$ as $\mathcal{S}(\mathcal{K})$ and $\mathcal{M}\left(\mathcal{K}^{*}\right)$, respectively, since $u$ is the identity matrix. A quantum system is locally quantum, but there are many other locally quantum systems. We are interested in common or different properties for locally quantum systems, which are studied, e.g., in 42]. However, many of such properties are not made clear and not studied well.

Structure of this section.-In this thesis, we study three topics associated with quantum systems and locally quantum systems: (i) capacity, (ii) differential privacy, and (iii) perfect discrimination of two states. The first one is a property that almost never depends on the locally quantum system, but the third one is a property that strongly depends on the locally quantum system. Although the first and third ones are discussed for locally quantum systems, the second one is only discussed in quantum theory. This is because the quantum case is still not sufficiently compared with the classical case in the viewpoint of (ii). Note that all the vector spaces we are dealing with in this thesis are finite-dimensional. In Sections 1.1, 1.2 and 1.3, we give a description of each of these three contributions.

### 1.1 Statement determining the capacities of GPTs, and weaker statement

Any problem for locally quantum systems can be regarded as a problem in matrix theory, since a locally quantum system is described by Hermitian matrices on $\mathbb{C}^{d_{1}} \otimes \cdots \otimes \mathbb{C}^{d_{n}}$. In particular, the so-called separability problem is one of the most important problems, and attracts attention in quantum theory [4, 25, 28, 31, 32]. The separability problem is the problem to investigate whether a given positive semidefinite matrix is separable or not (we state the definition of separability in the

[^2]next paragraph). There are several necessary or sufficient conditions for a positive semi-definite matrix to be separable [25, 31,32]. While motivated by a problem from locally quantum systems, in this thesis, we discuss a sufficient condition for separability (statement S below) as a linear algebraic one.

Let $\mathcal{H}\left(d_{1}, \ldots, d_{n}\right)$ be an $n$-partite complex Hilbert space $\mathbb{C}^{d_{1}} \otimes \cdots \otimes \mathbb{C}^{d_{n}}$, and $\tilde{d}$ be the dimension $d_{1} d_{2} \cdots d_{n}$ of $\mathcal{H}\left(d_{1}, \ldots, d_{n}\right)$. For each GPT, the maximum number of simultaneously and perfectly distinguishable states is called the capacity (for the definition of simultaneously and perfectly distinguishable states, see Section 1.3). The capacities of GPTs have been studied, e.g., in $[41-43]$. To determine the capacities of special GPTs, the following statement plays an important role (see Section 2.2):
(S) for every unit vector $u \in \mathcal{H}\left(d_{1}, \ldots, d_{n}\right)$, the two matrices $I_{\tilde{d}} \pm|u\rangle\langle u|$ lie in the set $\operatorname{Sep}\left(d_{1}, \ldots, d_{n}\right):=\operatorname{conv}\left\{X^{[1]} \otimes \cdots \otimes X^{[n]}: X^{[1]}, \ldots, X^{[n]}\right.$ positive semi-definite $\}$,
where $\operatorname{conv}(\mathcal{S})$ denotes the convex hull of a subset $\mathcal{S}$. Throughout this thesis, we use the bra-ket notation (see Section 2.2), and use the superscript [j] (resp. [k:l]) to express the $j$ th site (resp. the sites from $k$ th to $l$ th). A matrix in $\operatorname{Sep}\left(d_{1}, \ldots, d_{n}\right)$ is called separable. Every separable matrix is positive semi-definite, but the converse does not necessarily hold. If $n=2$, then statement S is true for all integers $d_{1}, d_{2} \geq 2$ [25]. However, statement S is still open for $n \geq 3$ to the best of our knowledge.

Main result.-Our main interest is whether statement $S$ holds for all integers $n \geq 2$ and $d_{1}, \ldots, d_{n} \geq 2$. In this thesis, we show the following statement weaker than statement S : for all integers $n \geq 2$ and $d_{1}, \ldots, d_{n} \geq 2$,
(WS) every $(\tilde{d}-1)$-dimensional subspace $\mathcal{L}$ of $\mathcal{H}\left(d_{1}, \ldots, d_{n}\right)$ has a product basis, where a vector $u \in \mathcal{H}\left(d_{1}, \ldots, d_{n}\right)$ is called a product vector if $u=u^{[1]} \otimes \cdots \otimes u^{[n]}$ for some $u^{[1]} \in \mathbb{C}^{d_{1}}, \ldots, u^{[n]} \in \mathbb{C}^{d_{n}} ;$ a basis composed of product vectors is called a product basis. In Section 3.1, we give a $(\tilde{d}-2)$-dimensional subspace with no product basis. These results yield the following theorem.

Theorem 1.1. For all integers $n \geq 2$ and $d_{1}, \ldots, d_{n} \geq 2$,

$$
\max \left\{\operatorname{dim} \mathcal{L}: \begin{array}{l}
\mathcal{L} \text { is a subspace of } \mathcal{H}\left(d_{1}, \ldots, d_{n}\right) \text { and } \\
\text { has no product basis }
\end{array}\right\}=\tilde{d}-2
$$

Since Theorem 1.1 can be regarded simply as a mathematical one, we generalize the scalar field $\mathbb{C}$ to an arbitrary field. Theorem 1.1 is still true even if the scalar field $\mathbb{C}$ is replaced with an arbitrary infinite field (Remark 3.5). We address the case of finite fields in Section 3.2. For every finite field $\mathcal{F}$, we show that Theorem 1.1 is
also true if either (i) $n=2$ or (ii) $n \geq 3$ and $\# \mathcal{F}>\max \left\{d_{i}: i \neq n_{1}, n_{2}\right\}$ for some $n_{1}$ and $n_{2}$ (Theorem 3.9), where $\# \mathcal{F}$ denotes the order of $\mathcal{F}$.

Related work.-Existing studies often consider an orthogonal product basis [2,7, 10, 13, 21, 32, which is defined as an orthonormal basis composed of product vectors. An orthogonal product basis of a subspace $\mathcal{L}$ is called unextendible if the orthogonal complement of $\mathcal{L}$ contains no non-zero product vector. Unextendible orthogonal product bases (UPBs; "orthogonal" is usually omitted) are used to construct bound entangled states $7,73,32$. In particular, quantum information theory motivates us to find UPBs of the minimum possible number. Alon and Lovász [2] proved that the minimum dimension of subspaces of $\mathcal{H}\left(d_{1}, \ldots, d_{n}\right)$ with UPBs is equal to $d_{1}+\cdots+d_{n}-n+1$ unless either (i) $n=2$ and $2 \in\left\{d_{1}, d_{2}\right\}$ or (ii) $d_{1}+\cdots+d_{n}-n+1$ is odd and at least one $d_{i}$ is even. Moreover, the minimum dimension is strictly greater than $d_{1}+\cdots+d_{n}-n+1$ in cases (i) and (ii). After the appearance of their work, cases (i) and (ii) have been studied in more detail [10, 21].

There are two statements similar to Theorem 1.1. A subspace of $\mathcal{H}\left(d_{1}, \ldots, d_{n}\right)$ containing no non-zero product vector is called completely entangled. Wallach [51] and Parthasarathy [45] proved that for all integers $n \geq 2$ and $d_{1}, \ldots, d_{n} \geq 2$,

$$
\begin{align*}
& \max \left\{\operatorname{dim} \mathcal{L}: \mathcal{L} \text { is a completely entangled subspace of } \mathcal{H}\left(d_{1}, \ldots, d_{n}\right)\right\} \\
= & \tilde{d}-\left(d_{1}+\cdots+d_{n}\right)+n-1 \tag{1.3}
\end{align*}
$$

Cubitt et al. [12] proved that for all integers $d_{1}, d_{2} \geq 2$ and $r \in\left[0, \min \left\{d_{1}, d_{2}\right\}-1\right]$,

$$
\begin{align*}
& \max \left\{\operatorname{dim} \mathcal{L}: \begin{array}{l}
\mathcal{L} \text { is a subspace of } \mathcal{H}\left(d_{1}, d_{2}\right) \text { satisfying that } \\
\text { s-rank } u \geq r+1 \text { for all non-zero } u \in \mathcal{L}
\end{array}\right\}  \tag{1.4}\\
& \leq\left(d_{1}-r\right)\left(d_{2}-r\right),
\end{align*}
$$

where s-rank $u$ denotes the Schmidt rank of $u \in \mathcal{H}\left(d_{1}, d_{2}\right)$, i.e., a unique number $k$ such that $u$ is expressed as $u=\sum_{i=1}^{k} \alpha_{i} u_{i}^{[1]} \otimes u_{i}^{[2]}$ with positive numbers $\alpha_{1}, \ldots, \alpha_{k}$ and two orthonormal systems $\left(u_{i}^{[1]}\right)_{i=1}^{k}$ and $\left(u_{i}^{[2]}\right)_{i=1}^{k}$. If $r=0,1$, then (1.4) has equality due to (1.3). Recently, Bag et al. [5] constructed subspaces that achieve equality in (1.4) for all $d_{1}, d_{2} \geq 4$ and $r=1,2,3$.

Actually, if $n=2$, then for all integers $d_{1}, d_{2} \geq 2$, the following statement holds (25):
(SS) the matrix $I+X$ lies in $\operatorname{Sep}\left(d_{1}, \ldots, d_{n}\right)$ for every Hermitian matrix $X$ with $\|X\|_{2}:=\left(\operatorname{Tr} X^{*} X\right)^{1 / 2} \leq 1$.

Statement SS is stronger than statement S and does not hold in general [4, 28.

Finally, we briefly describe entangled vectors which are related to product vectors. If $u \in \mathcal{H}\left(d_{1}, \ldots, d_{n}\right)$ is not any product vectors, we say that $u$ is entangled. It is widely known [32] that for every entangled vector $u \in \mathcal{H}\left(d_{1}, \ldots, d_{n}\right)$, the matrix $|u\rangle\langle u|$ does not lie in $\operatorname{Sep}\left(d_{1}, \ldots, d_{n}\right)$ (see also Proposition 3.6). Normalized entangled vectors are significant resources in quantum information processing.

### 1.2 Differential privacy

In data analysis, data analysts need to know only some statistical information about private data while protecting the private data. They hope to maximally utilize private data under some privacy protection. In general, protection and utilization of private data have a trade-off relation, and researchers optimize the trade-off relation [22, 24, 30, 38, 55].

As a way protecting private data, Warner [52] proposed randomized response in 1965, in which private data $X$ are converted to other data $Y$ subject to a conditional probability distribution $\mathbb{P}_{Y \mid X}$, and the data $Y$ is released instead of $X$. Since a data analyst collects only randomized data $Y$, private data $X$ are protected.

However, private data are not always protected in the above way. For instance, if $Y$ is always equal to $X$, then it is clear that private data are not protected. To further understand such issues, consider another case: suppose that $X$ is binary, $Y$ is ternary,

$$
\mathbb{P}_{Y \mid X}(\cdot \mid 0)=\begin{array}{r}
Y=0 \quad Y=1 \quad Y=2 \\
{[1 / 3,1 / 3,1 / 3]}
\end{array} \quad \text { and } \quad \mathbb{P}_{Y \mid X}(\cdot \mid 1)=\left[\begin{array}{c}
Y=0 \quad Y=1 \quad Y=2 \\
0,1 / 2,1 / 2
\end{array}\right] .
$$

Then the data analyst finds that $X=0$ if $Y=0$, since the conditional probability distribution satisfies that $\mathbb{P}_{Y \mid X}(0 \mid 0)>0$ and $\mathbb{P}_{Y \mid X}(0 \mid 1)=0$.

Differential privacy.-To enforce protection of data, we impose the following condition on the conditional probability distribution $\mathbb{P}_{Y \mid X}$ :

$$
\begin{equation*}
\forall x, x^{\prime}, \mathbb{P}_{Y \mid X}\left(\cdot \mid x^{\prime}\right) \leq e^{\varepsilon} \mathbb{P}_{Y \mid X}(\cdot \mid x) \tag{1.5}
\end{equation*}
$$

where $\varepsilon>0$ is a constant. This condition is called $\varepsilon$-differential privacy ( $\varepsilon-D P$ ) [17-19]. Differential privacy (DP) was introduced by Dwork et al. [19] and Dwork 18] in the global privacy context (the case when a company or government releases users' data partially for machine learning). After that, DP was also introduced by Duchi et al. [17] in the local privacy context (the case when data providers do not trust a data analyst). Definition (1.5) is that in the local privacy context.

DP has been studied intensively by using classical probability theory. However, there are only a few studies of quantum versions of DP $[1,15,16,56]$ to the best of our knowledge. In this thesis, we define a quantum version of DP and investigate its mathematical aspects when $n$-ary data $X$ are converted to quantum states $\rho$ depending on $X$ (i.e., classical-quantum setting) in the local privacy context.

To define a quantum version of DP, we consider an $n$-tuple of quantum states $\left(\rho_{x}\right)_{x=1}^{n}$, where $x$ and $\rho_{x}$ correspond to an input classical state and its output quantum state, respectively. Now, a data analyst needs to measure a quantum state by a measurement $\left(M_{y}\right)_{y=1}^{m}$ in order to obtain some information about the quantum state. Hence, following the classical definition of DP, we define the classical-quantum $\varepsilon$ differential privacy ( $C Q \varepsilon-D P$ ) as

$$
\forall\left(M_{y}\right)_{y=1}^{m} \text { POVM, the c.p.d. } \mathbb{P}(y \mid x)=\operatorname{Tr} \rho_{x} M_{y} \text { satisfies } \varepsilon \text {-DP }
$$

where "c.p.d." is an abbreviation of "conditional probability distribution". This condition is equivalent to the following one:

$$
\forall x, x^{\prime}, \rho_{x^{\prime}} \leq e^{\varepsilon} \rho_{x}
$$

where for Hermitian matrices $H$ and $H^{\prime}$ the inequality $H \leq H^{\prime}$ means for $H^{\prime}-H$ to be positive semi-definite. The definition of $\mathrm{CQ} \varepsilon$-DP is a simple extension of the classical one, because (1.5) can be written as

$$
\forall x, x^{\prime}, p_{x^{\prime}} \leq e^{\varepsilon} p_{x}
$$

if replacing the probability distributions $\mathbb{P}_{Y \mid X}(\cdot \mid x)$ with probability vectors $p_{x}$, where for probability vectors $p$ and $p^{\prime}$ the inequality $p \leq p^{\prime}$ means for $p^{\prime}-p$ to be nonnegative. From now on, we use an $n$-tuple $\left(p_{i}\right)_{i=1}^{n}$ of probability vectors instead of $\left(\mathbb{P}_{Y \mid X}(\cdot \mid x)\right)_{x=1}^{n}$.

We summarize the above definitions.
Definition 1.2 (Classical $\varepsilon$-DP 17 and classical-quantum $\varepsilon$-DP). Let $\varepsilon>0$ be a real number and $n \geq 2$ be an integer. An n-tuple $\left(p_{i}\right)_{i=1}^{n}$ of probability vectors is called $\varepsilon$-differentially private $(\varepsilon-D P)$ if $p_{i} \leq e^{\varepsilon} p_{j}$ for all $i, j=1, \ldots, n$. An n-tuple $\left(\rho_{i}\right)_{i=1}^{n}$ of density matrices is called classical-quantum $\varepsilon$-differentially private ( $C Q \varepsilon$-DP) if $\rho_{i} \leq e^{\varepsilon} \rho_{j}$ for all $i, j=1, \ldots, n$. Also, define the sets $\mathrm{C}_{n}^{(d)}(\varepsilon), \mathrm{CQ}_{n}^{(d)}(\varepsilon), \mathrm{C}_{n}(\varepsilon)$ and $\mathrm{CQ}_{n}(\varepsilon)$ as

$$
\begin{gathered}
\mathrm{C}_{n}^{(d)}(\varepsilon)=\left\{\varepsilon-D P\left(p_{i}\right)_{i=1}^{n}: \text { all } p_{i} \text { are probability vectors in } \mathbb{R}^{d}\right\} \quad(d \geq 2), \\
\mathrm{CQ}_{n}^{(d)}(\varepsilon)=\left\{C Q \varepsilon-D P\left(\rho_{i}\right)_{i=1}^{n}: \text { all } \rho_{i} \text { are density matrices on } \mathbb{C}^{d}\right\} \quad(d \geq 2), \\
\mathrm{C}_{n}(\varepsilon)=\bigcup_{d \geq 2} \mathrm{C}_{n}^{(d)}(\varepsilon), \quad \mathrm{CQ}_{n}(\varepsilon)=\bigcup_{d \geq 2} \mathrm{CQ}_{n}^{(d)}(\varepsilon) .
\end{gathered}
$$

If $\left(\rho_{i}\right)_{i=1}^{n}$ is $\mathrm{CQ} \varepsilon$ - DP, then all $\rho_{i}$ have the same support, i.e., all the ranges of $\rho_{i}$ are equal to one another. Hence, we often implicitly assume that all $\rho_{i}$ have full rank if $\left(\rho_{i}\right)_{i=1}^{n}$ is CQ $\varepsilon$-DP.

Embedding classical states into quantum ones.-Next, let us consider a subset of $\mathrm{CQ}_{n}(\varepsilon)$ that corresponds to $\mathrm{C}_{n}(\varepsilon)$. For a probability vector $p=(p(i))_{i=1}^{d} \in$ $\mathbb{R}^{d}$, define $\operatorname{diag}(p)$ as the diagonal matrix with diagonal entries $p(1), \ldots, p(d)$, which is a density matrix on $\mathbb{C}^{d}$. Since a quantum (resp. classical) state is a density matrix (resp. probability vector), the mapping $\operatorname{diag}(\cdot)$ is an embedding from the set of classical states into the set of quantum ones. Using the mapping $\operatorname{diag}(\cdot)$, we obtain the set

$$
\operatorname{diag}\left(\mathrm{C}_{n}(\varepsilon)\right):=\left\{\left(\operatorname{diag}\left(p_{i}\right)\right)_{i=1}^{n}:\left(p_{i}\right)_{i=1}^{n} \in \mathrm{C}_{n}(\varepsilon)\right\}
$$

that corresponds to $\mathrm{C}_{n}(\varepsilon)$.
Essentially classical elements.-The set $\operatorname{diag}\left(\mathrm{C}_{n}(\varepsilon)\right)$ is much smaller than $\mathrm{CQ}_{n}(\varepsilon)$, but actually, there is a set larger than $\operatorname{diag}\left(\mathrm{C}_{n}(\varepsilon)\right)$ that is "essentially classical". To describe such a set, we consider two optimization problems: one is the classical case
which is often considered in information-theoretic studies of DP [22, 24, 30, 38]; the other is the quantum case

$$
S_{n}^{\mathrm{CQ}}(\varepsilon ; \Phi)=\underbrace{\left(\rho_{i} i_{i=1}^{n} \in \operatorname{CQ}_{n}(\varepsilon)\right.}_{\text {Privacy protection }} \underbrace{\Phi\left(\rho_{1}, \ldots, \rho_{n}\right)}_{\text {Utility }} .
$$

The above $\Phi$ is a real-valued function of $n$ density matrices that represents the utility of private data, and the conditions $\left(p_{i}\right)_{i=1}^{n} \in \mathrm{C}_{n}(\varepsilon)$ and $\left(\rho_{i}\right)_{i=1}^{n} \in \mathrm{CQ}_{n}(\varepsilon)$ represent the privacy protection. Since the data analyst's purpose is to maximally utilize private data under the privacy protection, we arrive at the above optimization problems.

Now, we want to define a subset of $\mathrm{CQ}_{n}(\varepsilon)$ that is "essentially classical". For this purpose, assume that the objective function $\Phi$ must satisfy monotonicity for completely positive and trace-preserving linear maps (CPTP maps).

Definition 1.3 (Monotonicity for CPTP maps). A real-valued function $\Phi$ of $n$ density matrices is called monotone for CPTP maps if

$$
\Phi\left(\Lambda\left(\rho_{1}\right), \ldots, \Lambda\left(\rho_{n}\right)\right) \leq \Phi\left(\rho_{1}, \ldots, \rho_{n}\right)
$$

for all density matrices $\rho_{1}, \ldots, \rho_{n}$ and CPTP maps $\Lambda$. This inequality is called the information processing inequality (or data processing inequality).

Since a CPTP map is regarded as a quantum operation in quantum information theory, information-theoretic quantities usually satisfy monotonicity for CPTP maps. For example, quantum relative entropy, symmetric logarithmic derivative (SLD) Fisher information, Kubo-Mori-Bogoljubov (KMB) Fisher information, right logarithmic derivative (RLD) Fisher information, and trace distance satisfy monotonicity for CPTP maps [26, Theorems 5.7 and 6.2], [27, Theorem 6.7 and Lemma 6.9].

By monotonicity for CPTP maps, it follows that

$$
\sup _{\substack{\left(p_{i}\right)_{i=1}^{n} \in \mathrm{C}_{n}(\varepsilon) \\ \Lambda \mathrm{CPTP} \text { map }}} \Phi\left(\Lambda\left(\operatorname{diag}\left(p_{1}\right)\right), \ldots, \Lambda\left(\operatorname{diag}\left(p_{n}\right)\right)\right) \leq S_{n}^{\mathrm{C}}(\varepsilon ; \Phi) .
$$

Moreover, the opposite inequality also holds, since the identity mapping on Herm ( $d$ ) is a CPTP map. This fact leads us to the following definition.

Definition 1.4 (Essentially classical element). Let $\varepsilon>0$ be a real number and $n \geq 2$ be an integer. We say that $\left(\rho_{i}\right)_{i=1}^{n} \in \mathrm{CQ}_{n}(\varepsilon)$ is essentially classical if there exist an $\varepsilon$-DP n-tuple $\left(p_{i}\right)_{i=1}^{n} \in \mathrm{C}_{n}(\varepsilon)$ and a CPTP map $\Lambda$ such that $\Lambda\left(\operatorname{diag}\left(p_{i}\right)\right)=\rho_{i}$ for all $i=1, \ldots, n$. We denote by $\mathrm{EC}_{n}(\varepsilon)$ the set of all essentially classical elements in $\mathrm{CQ}_{n}(\varepsilon)$.

Although an element in $\mathrm{EC}_{n}(\varepsilon)$ consists of quantum states, the equality

$$
S_{n}^{\mathrm{EC}}(\varepsilon ; \Phi)=S_{n}^{\mathrm{C}}(\varepsilon ; \Phi)
$$

holds, where $S_{n}^{\mathrm{EC}}(\varepsilon ; \Phi)$ is defined in the same way as $S_{n}^{\mathrm{CQ}}(\varepsilon ; \Phi)$. Hence, the comparison of $S_{n}^{\mathrm{C}}(\varepsilon ; \Phi)$ and $S_{n}^{\mathrm{CQ}}(\varepsilon ; \Phi)$ is the same as that of $S_{n}^{\mathrm{EC}}(\varepsilon ; \Phi)$ and $S_{n}^{\mathrm{CQ}}(\varepsilon ; \Phi)$.

Comparison of $\mathrm{EC}_{n}(\varepsilon)$ and $\mathrm{CQ}_{n}(\varepsilon)$.-Although the set $\mathrm{EC}_{n}(\varepsilon)$ is a subset of $\mathrm{CQ}_{n}(\varepsilon)$, we are interested in whether they are equal to each other or not. If $\mathrm{EC}_{n}(\varepsilon)$ is equal to $\mathrm{CQ}_{n}(\varepsilon)$, then $S_{n}^{\mathrm{C}}(\varepsilon ; \Phi)=S_{n}^{\mathrm{EC}}(\varepsilon ; \Phi)=S_{n}^{\mathrm{CQ}}(\varepsilon ; \Phi)$, i.e., $\mathrm{CQ} \varepsilon$-DP mechanisms have no quantum advantage in optimization. In this perspective, it is important to compare $\mathrm{EC}_{n}(\varepsilon)$ with $\mathrm{CQ}_{n}(\varepsilon)$.

Main results.-In this thesis, we show the following theorems.
Theorem 1.5. For all $\varepsilon>0, \mathrm{EC}_{2}(\varepsilon)=\mathrm{CQ}_{2}(\varepsilon)$.
Theorem 1.6. For all $\varepsilon>0$ and $n \geq 3, \mathrm{EC}_{n}(\varepsilon) \neq \mathrm{CQ}_{n}(\varepsilon)$.

Table 1.2: Quantum versions of DP.

|  | Input | Output | Context |
| :---: | :---: | :---: | :---: |
| This thesis | Classical | Quantum | Local privacy |
| Ref. 16 | Quantum | Quantum | Local privacy |
| Refs. 15 | Q6 | Quantum | Quantum |
| Global privacy |  |  |  |

By Theorem 1.5, it follows that $S_{2}^{\mathrm{C}}(\varepsilon ; \Phi)=S_{2}^{\mathrm{EC}}(\varepsilon ; \Phi)=S_{2}^{\mathrm{CQ}}(\varepsilon ; \Phi)$. We will actually prove Theorem 4.2 that is stronger than Theorem 1.5, and prove Theorem 1.6 by giving a concrete objective function $\Phi$ such that $S_{n}^{\mathrm{C}}(\varepsilon ; \Phi)=S_{n}^{\mathrm{EC}}(\varepsilon ; \Phi)<S_{n}^{\mathrm{CQ}}(\varepsilon ; \Phi)$ for every $n \geq 3$ (Theorem 4.5). Theorem 4.5 implies a sufficient condition for a CQ $\varepsilon$-DP $n$-tuple not to lie in $\mathrm{EC}_{n}(\varepsilon)$ (Corollary 4.6). Using Corollary 4.6, we construct $\mathrm{CQ} \varepsilon$-DP $n$-tuples that do not lie in $\mathrm{EC}_{n}(\varepsilon)$ (Section 4.5 ).

We mention a relation among this thesis and existing studies briefly. As seen from Table 1.2, Refs. [1, 15, 16, 56] consider the case when input and output states are quantum. The definition of $\mathrm{CQ} \varepsilon$-DP can be regarded as a special case of quantum DP [16], but [16] does not include our results.

Supplement on the set $\mathrm{EC}_{n}(\varepsilon)$.-Actually, the set $\mathrm{EC}_{n}(\varepsilon)$ can be written without CPTP maps.
Proposition 1.7. For all $\varepsilon>0$ and $n \geq 2$,

$$
\begin{equation*}
\mathrm{EC}_{n}(\varepsilon)=\left\{\left(\sum_{k} p_{i}(k) \sigma_{k}\right)_{i=1}^{n}:\left(p_{i}\right)_{i=1}^{n} \in \mathrm{C}_{n}(\varepsilon), \text { density matrices } \sigma_{k}\right\} \tag{1.7}
\end{equation*}
$$

where the above sum is taken all over $k=1, \ldots, d$ if $d$ is the dimension of the vector space that $p_{1}, \ldots, p_{n}$ inhabit.

Proposition 1.7 can easily be checked; see Section 4.1. Although we have defined the set $\mathrm{EC}_{n}(\varepsilon)$ with CPTP maps, the same set is obtained even if replacing CPTP maps with positive and trace-preserving linear maps (PTP maps). That is, complete positivity is unnecessary, and positivity suffices in Definition 1.4. However, we have used CPTP maps in Definition 1.4 because CPTP maps are more natural in quantum information theory than PTP maps, and monotonicity for CPTP maps is used in Section 4.3.

### 1.3 Perfect discrimination of two states

Let $(\mathcal{V}, \mathcal{K}, u)$ be a GPT. We say that $m$ states $x_{1}, \ldots, x_{m} \in \mathcal{S}(\mathcal{K}, u)$ are simultaneously and perfectly distinguishable if there exists a measurement $\left(y_{j}\right)_{j=1}^{m} \in \mathcal{M}\left(\mathcal{K}^{*}, u\right)$ such
that $\left\langle x_{i}, y_{j}\right\rangle=\delta_{i, j}$ for all $i, j=1, \ldots, m$, where $\delta_{i, j}$ denotes the Kronecker delta. In the case $m=2$, two states $x_{1}$ and $x_{2}$ are called perfectly distinguishable if they are simultaneously and perfectly distinguishable. It is a fundamental problem whether two given states are perfectly distinguishable or not.

In quantum theory, the following equivalence is widely known (see a textbook in quantum information theory, e.g., [26, p. 138], 27, Proposition 5.13]):
(E) two states $\rho_{1}$ and $\rho_{2}$ are perfectly distinguishable if and only if $\operatorname{Tr} \rho_{1} \rho_{2}=0$.

However, equivalence E does not hold for a general GPT even if the GPT is locally quantum [3]. Some existing studies other than [3] show a difference between quantum theory and an alternative model, but they consider only GPTs that are too far from quantum theory. For example, quantum theory has an upper bound called Tsirelson's bound in the CHSH inequality (named after Clauser, Horne, Shimony, and Holt [11]). The GPT called PR box violates it [46, 47], but the PR box is never close to quantum theory, since the set of states in the PR box is a square. Such GPTs are unlikely to exist in reality. Hence, it is more insightful to consider GPTs that are sufficiently close to quantum theory in order to further understand quantum theory.

Main results.-In this thesis, we investigate perfect discrimination of two states in a bipartite locally quantum system $\left(\operatorname{Herm}\left(d_{1} d_{2}\right), \mathcal{K}, I_{d_{1} d_{2}}\right)$, i.e., a whole system of two quantum subsystems $\left(\operatorname{Herm}\left(d_{1}\right), \operatorname{PSD}\left(d_{1}\right), I_{d_{1}}\right)$ and $\left(\operatorname{Herm}\left(d_{2}\right), \operatorname{PSD}\left(d_{2}\right), I_{d_{2}}\right)$. Since our purpose is to make $\mathcal{K}$ sufficiently close to $\operatorname{PSD}\left(d_{1} d_{2}\right)$, we construct two convex cones $\mathcal{K}_{s}^{\text {neg }}$ and $\mathcal{K}_{s}^{\text {sc }}$ for $s \geq 0$ in a certain natural manner (for details, see Section 5.1), and consider the locally quantum systems

$$
\begin{equation*}
\left(\operatorname{Herm}\left(d_{1} d_{2}\right),\left(\mathcal{K}_{s}^{\mathrm{neg}}\right)^{*}, I_{d_{1} d_{2}}\right) \quad \text { and } \quad\left(\operatorname{Herm}\left(d_{1} d_{2}\right),\left(\mathcal{K}_{s}^{\mathrm{sc}}\right)^{*}, I_{d_{1} d_{2}}\right) \tag{1.8}
\end{equation*}
$$

The resulting convex cones $\mathcal{K}_{s}^{\text {neg }}, \mathcal{K}_{s}^{\text {sc }},\left(\mathcal{K}_{s}^{\text {neg }}\right)^{*}$ and $\left(\mathcal{K}_{s}^{\text {sc }}\right)^{*}$ are sufficiently close to $\operatorname{PSD}\left(d_{1} d_{2}\right)$ if $s$ is small enough. In this sense, the locally quantum systems (1.8) are sufficiently close to quantum theory if $s$ is small enough. Nevertheless, the locally quantum systems (1.8) violate equivalence E unless $s=0$. More precisely, we show the following theorems in Section 5.2.

Theorem 1.8 (Perfect discrimination with $\mathcal{M}\left(\mathcal{K}_{s}^{\text {neg }}\right)$ ). Let $s$ be a real number in the interval $[0,1 / 4]$, and $\rho_{1}^{[1]}, \rho_{1}^{[2]}, \rho_{2}^{[1]}, \rho_{2}^{[2]}$ be rank-one density matrices. If the point $\left(\operatorname{Tr} \rho_{1}^{[1]} \rho_{2}^{[1]}, \operatorname{Tr} \rho_{1}^{[2]} \rho_{2}^{[2]}\right)$ belongs to the set

$$
\begin{equation*}
\left\{(x, y) \in[0,1]^{2}: x y \leq 16 s^{2}(1-x)(1-y)\right\} \tag{1.9}
\end{equation*}
$$

then the two states $\rho_{1}=\rho_{1}^{[1]} \otimes \rho_{1}^{[2]}$ and $\rho_{2}=\rho_{2}^{[1]} \otimes \rho_{2}^{[2]}$ are perfectly distinguishable by some measurement in $\mathcal{M}\left(\mathcal{K}_{s}^{\text {neg }}\right)$.




Figure 1.1: The above figures illustrate the set $(1.9)$ for $s=0,1 / 8,1 / 4$, from left to right. They also correspond to the set 1.10 for $(s, t)=(0,0),(2 / 5,1 / 4),(1 / 2,1)$, from left to right. ${ }^{1}$

Theorem 1.9 (Perfect discrimination with $\mathcal{M}\left(\mathcal{K}_{s}^{\mathrm{sc}}\right)$ ). Let $t$ be a real number in the interval $[0,1]$, and $\rho_{1}^{[1]}, \rho_{1}^{[2]}, \rho_{2}^{[1]}, \rho_{2}^{[2]}$ be rank-one density matrices. If the point $\left(\operatorname{Tr} \rho_{1}^{[1]} \rho_{2}^{[1]}, \operatorname{Tr} \rho_{1}^{[2]} \rho_{2}^{[2]}\right)$ belongs to the set

$$
\begin{equation*}
\left\{(x, y) \in[0,1]^{2}: x y \leq t(1-x)(1-y)\right\} \tag{1.10}
\end{equation*}
$$

then the two states $\rho_{1}=\rho_{1}^{[1]} \otimes \rho_{1}^{[2]}$ and $\rho_{2}=\rho_{2}^{[1]} \otimes \rho_{2}^{[2]}$ are perfectly distinguishable by some measurement in $\mathcal{M}\left(\mathcal{K}_{s}^{\mathrm{sc}}\right)$ with $s=\sqrt{t} /(1+t)$.

The sets (1.9) and 1.10 are illustrated as Figure 1.1. Once the parameter $s$ decreases, the sets (1.9) and (1.10) become smaller. Since the orthogonality $\operatorname{Tr} \rho_{1} \rho_{2}=$ 0 can be rewritten as $\left(\operatorname{Tr} \rho_{1}^{[1]} \rho_{2}^{[1]}\right)\left(\operatorname{Tr} \rho_{1}^{[2]} \rho_{2}^{[2]}\right)=0$, we find that the locally quantum systems (1.8) violate equivalence E unless $s=0$. Also, a rank-one density matrix is called pure, and often plays an important role in quantum information processing. Hence, Theorems 1.8 and 1.9 are meaningful results.

[^3]
## Chapter 2

## GPTs and their capacities

In this chapter, we describe the framework of GPTs, define the capacities of GPTs, and prove a few statements on capacities briefly. The contents of this chapter except for Proposition 2.3 rely on [54] (framework of GPTs) and [42] (capacities) essentially.

### 2.1 Framework of GPTs

A GPT is defined as the triplet of (i) a real Hilbert space $\mathcal{V}$ equipped with an inner product $\langle\cdot, \cdot\rangle$, (ii) a proper cone $\mathcal{K}$ of $\mathcal{V}$, and (iii) an element $u$ in the interior of the dual cone $\mathcal{K}^{*}$, where a subset $\mathcal{K}$ of $\mathcal{V}$ is called a convex cone if $\alpha u+\beta v \in \mathcal{K}$ for all $u, v \in \mathcal{K}$ and $\alpha, \beta \geq 0$; a convex cone $\mathcal{K}$ is called proper if $\mathcal{K}$ is closed, has an interior point, and satisfies $\mathcal{K} \cap(-\mathcal{K})=\{0\}$. Throughout this thesis, we assume that $\mathcal{V}$ is finite-dimensional. It is widely known that (see a textbook in convex analysis, e.g., [8, p. 53])

- if $\mathcal{K}$ is a non-empty closed convex cone, then $\mathcal{K}^{* *}=\mathcal{K}$;
- if $\mathcal{K}$ is a proper cone, then so is $\mathcal{K}^{*}$.

The element $u$ is called a unit effect and fixed for each GPT. Also, for a GPT $(\mathcal{V}, \mathcal{K}, u)$, the state class $\mathcal{S}(\mathcal{K}, u)$ and measurement class $\mathcal{M}\left(\mathcal{K}^{*}, u\right)$ are defined as (1.1). As a convex cone to define a measurement class, one can use another proper cone $\mathcal{K}^{\prime} \subset \mathcal{K}^{*}$ instead of $\mathcal{K}^{*}$, but the condition $\mathcal{K}^{\prime}=\mathcal{K}^{*}$ is imposed in usual, which is assumed in this thesis.

Given a state $x \in \mathcal{S}(\mathcal{K}, u)$ and a measurement $\left(y_{i}\right)_{i=1}^{m} \in \mathcal{M}\left(\mathcal{K}^{*}, u\right)$, the probability of obtaining each outcome $i=1, \ldots, m$ is $\left\langle x, y_{i}\right\rangle$. Note that $\left(\left\langle x, y_{i}\right\rangle\right)_{i=1}^{m}$ is a probability vector due to definition (1.1).

For $n \operatorname{GPTs}\left(\mathcal{V}^{[i]}, \mathcal{K}^{[i]}, u^{[i]}\right)$ (which are called subsystems), a whole system $(\mathcal{V}, \mathcal{K}, u)$ of these subsystems is defined as stated in Chapter 1. It is important that $\mathcal{K}$ is not uniquely determined even if its subsystems are given, since $\mathcal{K}_{\text {min }}$ and $\mathcal{K}_{\text {max }}$ defined in (1.2) are not equal to each other in general.

We recall the three types of GPTs below which have been defined in Chapter 1. For classical and quantum systems, see also Table 1.1.

Classical system.-The triplet $\left(\mathbb{R}^{d},[0, \infty)^{d}, \mathbf{1}_{d}\right)$ is a GPT called d-level classical system. The proper cone $[0, \infty)^{d}$ is self-dual, i.e., $\left([0, \infty)^{d}\right)^{*}=[0, \infty)^{d}$. If $\mathcal{K}^{[i]}=[0, \infty)^{d_{i}}$ for all $i=1, \ldots, n$, then it can easily be checked that $\mathcal{K}_{\text {min }}=\mathcal{K}_{\text {max }}$. Therefore, a whole system of classical subsystems is classical.

Quantum system.-The triplet $\left(\operatorname{Herm}(d), \operatorname{PSD}(d), I_{d}\right)$ is a GPT called d-level quantum system. Note that $\operatorname{PSD}(d)$ is self-dual, i.e., $\operatorname{PSD}(d)^{*}=\operatorname{PSD}(d)$. The state class $\mathcal{S}\left(\operatorname{PSD}(d), I_{d}\right)$ is the set of all density matrices, and the measurement class $\mathcal{M}\left(\operatorname{PSD}(d), I_{d}\right)$ is the set of all POVMs.

Locally quantum system.-For integers $d_{1}, \ldots, d_{n} \geq 2$, denote by $\tilde{d}$ the product $d_{1} d_{2} \cdots d_{n}$. For $d_{i}$-level quantum systems, $i=1, \ldots, n$, a whole system $(\mathcal{V}, \mathcal{K}, u)$ of these subsystems is called a $\left(d_{1}, \ldots, d_{n}\right)$-level locally quantum system. Let $(\mathcal{V}, \mathcal{K}, u)$ be a $\left(d_{1}, \ldots, d_{n}\right)$-level locally quantum system. Then

$$
\mathcal{V}=\operatorname{Herm}\left(d_{1}\right) \otimes \cdots \otimes \operatorname{Herm}\left(d_{n}\right) \quad \text { and } \quad u=I_{d_{1}} \otimes \cdots \otimes I_{d_{n}}=I_{\tilde{d}}
$$

The tensor product $\mathcal{V}$ can be regarded as a subspace of $\operatorname{Herm}(\tilde{d})$ by using the Kronecker product below. The tensor product of two matrices $A^{[1]}=\left(a_{i, j}^{[1]}\right)_{1 \leq i, j \leq d_{1}}$ and $A^{[2]}=\left(a_{i, j}^{[2]}\right)_{1 \leq i, j \leq d_{2}}$ can be expressed as a block matrix:

$$
A^{[1]} \otimes A^{[2]}=\left[\begin{array}{ccc}
a_{1,1}^{[1]} A^{[2]} & \cdots & a_{1, d_{1}}^{[1]} A^{[2]} \\
\vdots & \ddots & \vdots \\
a_{d_{1}, 1}^{[1]} A^{[2]} & \cdots & a_{d_{1}, d_{1}}^{[1]} A^{[2]}
\end{array}\right]
$$

This is called the Kronecker product. For $n \geq 3$, the tensor product of $n$ matrices can also be expressed in the same way (but we do not use it in this thesis). Under this identification, $\mathcal{V}$ is equal to $\operatorname{Herm}(\tilde{d})$, since (i) $\mathcal{V}$ is a subspace of $\operatorname{Herm}(\tilde{d})$, and (ii) the dimensions of $\mathcal{V}$ and $\operatorname{Herm}(\tilde{d})$ are equal to each other.

Let us return to the $\left(d_{1}, \ldots, d_{n}\right)$-level locally quantum system $(\mathcal{V}, \mathcal{K}, u)$. Using the set $\operatorname{Sep}\left(d_{1}, \ldots, d_{n}\right)$ defined in Section 1.1, we have

$$
\mathcal{K}_{\min }=\operatorname{Sep}\left(d_{1}, \ldots, d_{n}\right) \quad \text { and } \quad \mathcal{K}_{\max }=\operatorname{Sep}\left(d_{1}, \ldots, d_{n}\right)^{*}
$$

Thus, $\mathcal{K}$ satisfies the inclusion relation $\operatorname{Sep}\left(d_{1}, \ldots, d_{n}\right) \subset \mathcal{K} \subset \operatorname{Sep}\left(d_{1}, \ldots, d_{n}\right)^{*}$. Moreover, for all integers $n \geq 2$ and $d_{1}, \ldots, d_{n} \geq 2$, the two proper cones $\operatorname{Sep}\left(d_{1}, \ldots, d_{n}\right)$ and $\operatorname{Sep}\left(d_{1}, \ldots, d_{n}\right)^{*}$ are not equal to each other. Also, the $\tilde{d}$-level quantum system is a $\left(d_{1}, \ldots, d_{n}\right)$-level locally quantum system, since the inclusion relation $\operatorname{Sep}\left(d_{1}, \ldots, d_{n}\right) \subset$ $\operatorname{PSD}(\tilde{d}) \subset \operatorname{Sep}\left(d_{1}, \ldots, d_{n}\right)^{*}$ holds.

### 2.2 Capacities of locally quantum systems

For a $\operatorname{GPT}(\mathcal{V}, \mathcal{K}, u)$, simultaneously and perfectly distinguishable states are defined in Sections 1.3. We define the capacity of a GPT.

Definition 2.1 (Capacity). For a $\operatorname{GPT}(\mathcal{V}, \mathcal{K}, u)$, the maximum number of simultaneously and perfectly distinguishable states is called the capacity. We denote by $\operatorname{capa}(\mathcal{V}, \mathcal{K}, u)$ the capacity of a $\operatorname{GPT}(\mathcal{V}, \mathcal{K}, u)$.

For example, it is known that the capacity of $d$-level quantum system is equal to $d$. As proved below, the capacity of each $\left(d_{1}, d_{2}\right)$-level locally quantum system is equal to $d_{1} d_{2}$. This fact is found in (42] (without proof).

Proposition 2.2. Let $d_{1}, d_{2} \geq 2$ be integers. For every $\left(d_{1}, d_{2}\right)$-level locally quantum system, the capacity is equal to $d_{1} d_{2}$.

Proposition 2.2 asserts that the capacity of a locally quantum system only depends on its dimension. However, another property discussed in Chapter 5 changes depending on the convex cone of the locally quantum system.

Before proving Proposition 2.2, we describe notational conventions. A vector $u \in \mathbb{C}^{d}$ is expressed as a column vector. Also, we use the bra-ket notation: for $u \in \mathbb{C}^{d},|u\rangle$ and $\langle u|$ denote the column vector $u$ and its conjugate transpose, respectively. Hence, $\langle\cdot \mid \cdot\rangle$ gives the standard Hermitian inner product on $\mathbb{C}^{d}$, and $|u\rangle\langle u|$ is a rank-one orthogonal projection for every unit vector $u \in \mathbb{C}^{d}$. When $(\mathcal{V}, \mathcal{K}, u)$ is a locally quantum system, we simply write $\mathcal{S}(\mathcal{K}, u)$ and $\mathcal{M}\left(\mathcal{K}^{*}, u\right)$ as $\mathcal{S}(\mathcal{K})$ and $\mathcal{M}\left(\mathcal{K}^{*}\right)$, respectively.

Proof of Proposition 2.2. We use the fact that statement SS holds in the case $n=2$ [25]. First, let us show that every $Y \in \operatorname{Sep}\left(d_{1}, d_{2}\right)^{*}$ satisfies $\|Y\|_{2} \leq \operatorname{Tr} Y$. Since the case $Y=0$ is trivial, assume that $Y \in \operatorname{Sep}\left(d_{1}, d_{2}\right)^{*}$ is non-zero. Set $X=-Y /\|Y\|_{2}$. Then $I_{d_{1} d_{2}}+X$ lies in $\operatorname{Sep}\left(d_{1}, d_{2}\right)$. Thus, $\operatorname{Tr} Y-\|Y\|_{2}=\operatorname{Tr}\left(I_{d_{1} d_{2}}+X\right) Y \geq 0$.

Let $\left(\operatorname{Herm}\left(d_{1} d_{2}\right), \mathcal{K}, I_{d_{1} d_{2}}\right)$ be a $\left(d_{1}, d_{2}\right)$-level locally quantum system. Next, we show that the capacity is equal to $d_{1} d_{2}$. Suppose that $m$ states $\rho_{1}, \ldots, \rho_{m} \in \mathcal{S}(\mathcal{K})$
are simultaneously and perfectly distinguishable by a measurement $\left(M_{j}\right)_{j=1}^{m}$. Then

$$
\begin{align*}
& m \stackrel{(a)}{=} \sum_{i=1}^{m} \operatorname{Tr} \rho_{i} M_{i} \leq \sum_{i=1}^{m}\left\|\rho_{i}\right\|_{2}\left\|M_{i}\right\|_{2} \stackrel{(b)}{\leq} \sum_{i=1}^{m}\left(\operatorname{Tr} \rho_{i}\right)\left(\operatorname{Tr} M_{i}\right) \\
& \stackrel{(c)}{=} \sum_{i=1}^{m} \operatorname{Tr} M_{i} \stackrel{(d)}{=} \operatorname{Tr} I_{d_{1} d_{2}}=d_{1} d_{2}, \tag{2.1}
\end{align*}
$$

where $(a),(b),(c)$ and $(d)$ follow from the facts $\operatorname{Tr} \rho_{i} M_{i}=1, \rho_{i}, M_{i} \in \operatorname{Sep}\left(d_{1}, d_{2}\right)^{*}$, $\operatorname{Tr} \rho_{i}=\operatorname{Tr} \rho_{i} I_{d_{1} d_{2}}=1$ and $\sum_{i=1}^{m} M_{i}=I_{d_{1} d_{2}}$, respectively. Therefore, the capacity is less than or equal to $d_{1} d_{2}$. Since the $d_{1} d_{2}$ states

$$
\left|e_{i}^{[1]}\right\rangle\left\langle e_{i}^{[1]}\right| \otimes\left|e_{j}^{[2]}\right\rangle\left\langle e_{j}^{[2]}\right| \in \mathcal{S}(\mathcal{K}) \quad\left(1 \leq i \leq d_{1}, 1 \leq j \leq d_{2}\right)
$$

are simultaneously and perfectly distinguishable by the measurement $\left(\left|e_{i}^{[1]}\right\rangle\left\langle e_{i}^{[1]}\right| \otimes\right.$ $\left.\left|e_{j}^{[2]}\right\rangle\left\langle e_{j}^{[2]}\right|\right)_{i, j} \in \mathcal{M}\left(\mathcal{K}^{*}\right)$, we find that the capacity is equal to $d_{1} d_{2}$.

We have used statement SS with $n=2$ in the above proof, but statement SS is false in general [4, 28]. Instead of statement SS , let us focus on statement S . As already stated in Section 1.1, statement $S$ is still open for $n \geq 3$ to the best of our knowledge. Finally, assuming statement S , we show the following proposition.

Proposition 2.3. Assume that statement $S$ is true for all integers $n \geq 2$ and $d_{1}, \ldots, d_{n} \geq 2$. Let $\left(\operatorname{Herm}(\tilde{d}), \mathcal{K}, I_{\tilde{d}}\right)$ be a $\left(d_{1}, \ldots, d_{n}\right)$-level locally quantum system satisfying either $\mathcal{K} \subset \operatorname{PSD}(\tilde{d})$ or $\mathcal{K} \supset \operatorname{PSD}(\tilde{d})$. Then capa $\left(\operatorname{Herm}(\tilde{d}), \mathcal{K}, I_{\tilde{d}}\right)$ is equal to $\tilde{d}$.

Proof. First, let us show that every $Y \in \operatorname{Sep}\left(d_{1}, \ldots, d_{n}\right)^{*}$ satisfies $\|Y\| \leq \operatorname{Tr} Y$, where $\|\cdot\|$ denotes the operator norm. Due to statement S , it follows that $\operatorname{Tr} Y \pm\langle u| Y|u\rangle=$ $\operatorname{Tr}\left(I_{\tilde{d}} \pm|u\rangle\langle u|\right) Y \geq 0$ for every unit vector $u \in \mathcal{H}\left(d_{1}, \ldots, d_{n}\right)$. Thus, $\|Y\| \leq \operatorname{Tr} Y$.

The remainder is almost the same as the proof of Proposition 2.2. The difference between the proof of Proposition 2.2 and this proof is only 2.1 . We must change (2.1) as follows:

$$
\begin{aligned}
m & =\sum_{i=1}^{m} \operatorname{Tr} \rho_{i} M_{i} \leq \begin{cases}\sum_{i=1}^{m}\left\|\rho_{i}\right\|_{1}\left\|M_{i}\right\| & \mathcal{K} \subset \operatorname{PSD}(\tilde{d}), \\
\sum_{i=1}^{m}\left\|\rho_{i}\right\|\left\|M_{i}\right\|_{1} & \mathcal{K} \supset \operatorname{PSD}(\tilde{d})\end{cases} \\
& \leq \sum_{i=1}^{m}\left(\operatorname{Tr} \rho_{i}\right)\left(\operatorname{Tr} M_{i}\right)=\sum_{i=1}^{m} \operatorname{Tr} M_{i}=\operatorname{Tr} I_{\tilde{d}}=\tilde{d}
\end{aligned}
$$

where $\|\cdot\|_{1}$ denotes the trace norm.

## Chapter 3

## Weaker statement: Maximum dimension of subspaces with no product basis

In this chapter, we prove Theorem 1.1. Although the scalar field in Theorem 1.1 is the field of complex numbers, the same theorem holds in replacing it with an arbitrary infinite field. We also address the case of finite fields.

### 3.1 More general proposition and proof

Let $\left(e_{i}^{[j]}\right)_{i=1}^{d_{j}}$ be the standard basis of $\mathbb{C}^{d_{j}}$ for $j=1, \ldots, n$. Denote by $\operatorname{span}(\mathcal{S})$ the linear span of a subset $\mathcal{S}$, and by $\mathcal{L}^{\perp}$ the orthogonal complement of a subspace $\mathcal{L}$. Although product vectors have been already defined in the case $n \geq 2$, all vectors are regarded as product vectors in the case $n=1$.

Now, we prove the following proposition which is more general than statement WS.

Proposition 3.1. Let $n \geq 1$ and $d_{1}, \ldots, d_{n} \geq 2$ be integers, and $r$ be an integer in the interval $\left[0, \min \left\{d_{1}, \ldots, d_{n}\right\}\right]$. If the dimension of a subspace $\mathcal{L}$ of $\mathcal{H}\left(d_{1}, \ldots, d_{n}\right)$ is greater than or equal to $\tilde{d}-r$, then $\mathcal{L}$ has a $\left(\tilde{d}-r^{n}\right)$-tuple $\left(u_{i}\right)_{i=1}^{\tilde{d}-r^{n}}$ of linearly independent product vectors.

To prove Proposition 3.1, we need two lemmas. The first one is basic in algebra, and the second one is proved by using the first one.

Lemma 3.2. Let $\mathcal{F}$ be an infinite field, $n \geq 1$ be an integer, and $f\left(x_{1}, \ldots, x_{n}\right)$ be a polynomial over $\mathcal{F}$. Then the following conditions are equivalent:

1. $f\left(\alpha_{1}, \ldots, \alpha_{n}\right)=0$ for all $\alpha_{1}, \ldots, \alpha_{n} \in \mathcal{F}$;
2. $f\left(x_{1}, \ldots, x_{n}\right)=0$ as a polynomial.

Proof. See [34, Theorem 2.19].
Lemma 3.3. Let $m, n \geq 1, d_{1}, \ldots, d_{n} \geq 2$ and $r \in[1, \tilde{d}]$ be integers, and let $\left(u_{k, l}\right)_{l=1}^{\tilde{d}-r}$, $k=1, \ldots, m$, be $(\tilde{d}-r)$-tuples of linearly independent vectors in $\mathcal{H}\left(d_{1}, \ldots, d_{n}\right)$. Then there exist product vectors $v_{1}, \ldots, v_{r} \in \mathcal{H}\left(d_{1}, \ldots, d_{n}\right)$ such that

$$
\operatorname{det}\left[u_{k, 1}, \ldots, u_{k, \tilde{d}-r}, v_{1}, \ldots, v_{r}\right] \neq 0
$$

for all $k=1, \ldots, m$.
Proof. Let $\left(e_{i}\right)_{i=1}^{\tilde{d}}$ be the standard basis of $\mathcal{H}\left(d_{1}, \ldots, d_{n}\right)$, and let $n^{\prime}=r\left(d_{1}+\cdots+d_{n}\right)$. Define the $m$ polynomials $f_{k}\left(x_{1}, \ldots, x_{n^{\prime}}\right)$ over $\mathbb{C}$ as

$$
\begin{gathered}
f_{k}\left(x_{1}, \ldots, x_{n^{\prime}}\right)=\operatorname{det}\left[u_{k, 1}, \ldots, u_{k, \tilde{d}-r}, v_{1}, \ldots, v_{r}\right] \quad(k=1, \ldots, m), \\
v_{i}=v_{i}^{[1]} \otimes \cdots \otimes v_{i}^{[n]} \quad(i=1, \ldots, r)
\end{gathered}
$$

where the variables $x_{1}, \ldots, x_{n^{\prime}}$ correspond to the $n^{\prime}$ entries of $v_{i}^{[j]}, 1 \leq i \leq r$, $1 \leq j \leq n$. Since $\left(u_{k, l}\right)_{l=1}^{d-r}, k=1, \ldots, m$, are tuples of linearly independent vectors, we have $f_{k}\left(\alpha_{1}, \ldots, \alpha_{n^{\prime}}\right) \neq 0$ for some $\alpha_{1}, \ldots, \alpha_{n^{\prime}} \in \mathbb{C}$ corresponding to $v_{1}, \ldots, v_{r} \in\left\{e_{1}, \ldots, e_{\tilde{d}}\right\}$ (note that $e_{1}, \ldots, e_{\tilde{d}}$ are product vectors). Therefore, for every $k=1, \ldots, m$, the polynomial $f_{k}\left(x_{1}, \ldots, x_{n^{\prime}}\right)$ is not zero as a polynomial. Since the polynomial ring $\mathbb{C}\left[x_{1}, \ldots, x_{n^{\prime}}\right]$ is an integral domain, the product $f\left(x_{1}, \ldots, x_{n^{\prime}}\right):=$ $\prod_{k=1}^{m} f_{k}\left(x_{1}, \ldots, x_{n^{\prime}}\right)$ is not also zero as a polynomial. Thus, Lemma 3.2 implies that $f\left(\beta_{1}, \ldots, \beta_{n^{\prime}}\right) \neq 0$ for some $\beta_{1}, \ldots, \beta_{n^{\prime}} \in \mathbb{C}$. Taking the vectors $v_{i}^{[j]}$ corresponding to $\beta_{1}, \ldots, \beta_{n^{\prime}} \in \mathbb{C}$, we obtain desired product vectors $v_{1}, \ldots, v_{r}$.

Proof of Proposition 3.1. Since the case $r=0$ is clear, we assume the condition $r \geq 1$ in this proof. We show the proposition by induction on $n \geq 1$. First, the case $n=1$ is trivial. Let $n \geq 2$ and assume that the proposition is true for $n-1$. Then we show that the proposition is also true for $n$. Let the dimension of a subspace $\mathcal{L}$ be greater than or equal to $\tilde{d}-r$. For some $w_{1}, \ldots, w_{r} \in \mathcal{L}^{\perp}$, the subspace $\mathcal{L}$ can be expressed as

$$
\mathcal{L}=\left\{u \in \mathcal{H}\left(d_{1}, \ldots, d_{n}\right): \forall i=1, \ldots, r,\left\langle w_{i} \mid u\right\rangle=0\right\} .
$$

Now, take a basis $\left(u_{k}^{[1]}\right)_{k=1}^{d_{1}}$ of $\mathbb{C}^{d_{1}}$ and set $\tilde{d}^{\prime}=\tilde{d} / d_{1}$. Since the dimension of the subspace

$$
\mathcal{L}_{k}^{[2: n]}:=\left\{u^{[2: n]} \in \mathcal{H}\left(d_{2}, \ldots, d_{n}\right): \forall i=1, \ldots, r,\left\langle w_{i} \mid u_{k}^{[1]} \otimes u^{[2: n]}\right\rangle=0\right\}
$$

is greater than or equal to $\tilde{d^{\prime}}-r$ for every $k=1, \ldots, d_{1}$, the induction hypothesis implies that $\mathcal{L}_{k}^{[2: n]}$ has a $\left(\tilde{d}^{\prime}-r^{n-1}\right)$-tuple $\left(u_{k, l}^{[2: n]}\right)_{l=1}^{\tilde{d}^{\prime}-r^{n-1}}$ of linearly independent product vectors. Also, due to Lemma 3.3, we can take an $r^{n-1}$-tuple $\left(v_{s}^{[2: n]}\right)_{s=1}^{n-1}$ of product vectors with the following condition:

$$
\begin{equation*}
\forall k=1, \ldots, d_{1}, \operatorname{det}\left[u_{k, 1}^{[2: n]}, \ldots, u_{k, \tilde{d}^{\prime}-r^{n-1}}^{[2: n]}, v_{1}^{[2: n]}, \ldots, v_{r^{n-1}}^{[2: n]}\right] \neq 0 \tag{3.1}
\end{equation*}
$$

Moreover, for every $s=1, \ldots, r^{n-1}$, take a $\left(d_{1}-r\right)$-tuple $\left(v_{s, t}^{[1]}\right)_{t=1}^{d_{1}-r}$ of linearly independent vectors in the subspace

$$
\mathcal{L}_{s}^{[1]}:=\left\{u^{[1]} \in \mathbb{C}^{d_{1}}: \forall i=1, \ldots, r,\left\langle w_{i} \mid u^{[1]} \otimes v_{s}^{[2: n]}\right\rangle=0\right\}
$$

Note that the $r^{n-1}\left(d_{1}-r\right)$ vectors $v_{s, t}^{[1]} \otimes v_{s}^{[2: n]}$ are linearly independent.
Let us show that the $\tilde{d}-r^{n}$ product vectors of $\mathcal{L}$

$$
u_{k}^{[1]} \otimes u_{k, l}^{[2: n]}, \quad v_{s, t}^{[1]} \otimes v_{s}^{[2: n]} \quad\left(\begin{array}{c}
1 \leq k \leq d_{1},  \tag{3.2}\\
1 \leq s \leq r^{n-1}, \\
1 \leq l \leq \tilde{d}^{\prime}-r^{n-1} \\
1 \leq t \leq d_{1}-r
\end{array}\right)
$$

are linearly independent. Suppose that $\tilde{d}-r^{n}$ scalars $\alpha_{k, l}$ and $\beta_{s, t}$ satisfy

$$
\begin{equation*}
\sum_{k, l} \alpha_{k, l} u_{k}^{[1]} \otimes u_{k, l}^{[2: n]}+\sum_{s, t} \beta_{s, t} v_{s, t}^{[1]} \otimes v_{s}^{[2: n]}=0 \tag{3.3}
\end{equation*}
$$

Since $\left(u_{k}^{[1]}\right)_{k=1}^{d_{1}}$ is a basis of $\mathbb{C}^{d_{1}}$, for all $s$ and $t$, there exist scalars $\gamma_{s, t, k}$ such that $v_{s, t}^{[1]}=\sum_{k} \gamma_{s, t, k} u_{k}^{[1]}$. Thus, (3.3) can be rewritten as follows:

$$
\begin{aligned}
0 & =\sum_{k, l} \alpha_{k, l} u_{k}^{[1]} \otimes u_{k, l}^{[2: n]}+\sum_{s, t, k} \beta_{s, t} \gamma_{s, t, k} u_{k}^{[1]} \otimes v_{s}^{[2: n]} \\
& =\sum_{k} u_{k}^{[1]} \otimes\left(\sum_{l} \alpha_{k, l} u_{k, l}^{[2: n]}+\sum_{s, t} \beta_{s, t} \gamma_{s, t, k} v_{s}^{[2: n]}\right) .
\end{aligned}
$$

Since $\left(u_{k}^{[1]}\right)_{k=1}^{d_{1}}$ is a basis of $\mathbb{C}^{d_{1}}$, we have

$$
\sum_{l} \alpha_{k, l} u_{k, l}^{[2: n]}+\sum_{s, t} \beta_{s, t} \gamma_{s, t, k} v_{s}^{[2: n]}=0
$$

for every $k=1, \ldots, d_{1}$. This and (3.1) imply that $\alpha_{k, l}=0$ for all $k$ and $l$. Thus, (3.3) turns to $\sum_{s, t} \beta_{s, t} v_{s, t}^{[1]} \otimes v_{s}^{[2: n]}=0$. Since the $r^{n-1}\left(d_{1}-r\right)$ vectors $v_{s, t}^{[1]} \otimes v_{s}^{[2: n]}$ are linearly independent, it follows that $\beta_{s, t}=0$ for all $s$ and $t$. Therefore, the vectors (3.2) are linearly independent, and the proposition is also true for $n$.

Next, we construct a ( $\tilde{d}-2$ )-dimensional subspace with no product basis by using the case $n=2$.

Proposition 3.4. For all integers $n \geq 2$ and $d_{1}, \ldots, d_{n} \geq 2$, there exists a $(\tilde{d}-2)$ dimensional subspace of $\mathcal{H}\left(d_{1}, \ldots, d_{n}\right)$ with no product basis.

Proof. First, assuming $n=2$, we show that the $(\tilde{d}-2)$-dimensional subspace

$$
\begin{equation*}
\mathcal{L}^{[1: 2]}:=\operatorname{span}\left(\left\{e_{1}^{[1]} \otimes e_{1}^{[2]}+e_{2}^{[1]} \otimes e_{2}^{[2]}\right\} \cup\left\{e_{i}^{[1]} \otimes e_{j}^{[2]}:(i, j) \neq(1,1),(2,1),(2,2)\right\}\right) \tag{3.4}
\end{equation*}
$$

has no product basis. Take an arbitrary product vector $u=u^{[1]} \otimes u^{[2]} \in \mathcal{L}^{[1: 2]}$ with the expressions $u^{[k]}=\sum_{i=1}^{d_{k}} \alpha_{i}^{[k]} e_{i}^{[k]}, \alpha_{i}^{[k]} \in \mathbb{C}, k=1,2$. Then $u$ is expressed in two ways:

$$
\begin{aligned}
u & =\sum_{i, j} \alpha_{i}^{[1]} \alpha_{j}^{[2]} e_{i}^{[1]} \otimes e_{j}^{[2]} \\
& =\beta\left(e_{1}^{[1]} \otimes e_{1}^{[2]}+e_{2}^{[1]} \otimes e_{2}^{[2]}\right)+\sum_{(i, j) \neq(1,1),(2,1),(2,2)} \alpha_{i}^{[1]} \alpha_{j}^{[2]} e_{i}^{[1]} \otimes e_{j}^{[2]} \quad(\exists \beta \in \mathbb{C}),
\end{aligned}
$$

where the second equality follows from the basis $(3.4)$ of $\mathcal{L}^{[1: 2]}$. For readability, write the matrix representation of the above expression:

$$
\left[\begin{array}{cccc}
\alpha_{1}^{[1]} \alpha_{1}^{[2]} & \alpha_{1}^{[1]} \alpha_{2}^{[2]} & \cdots & \alpha_{1}^{[1]} \alpha_{d_{2}}^{[2]} \\
\alpha_{2}^{[1]} \alpha_{1}^{[2]} & \alpha_{2}^{[1]} \alpha_{2}^{[2]} & \cdots & \alpha_{2}^{[1]} \alpha_{d_{2}}^{[2]} \\
\vdots & \vdots & & \vdots \\
\alpha_{d_{1}}^{[1]} \alpha_{1}^{[2]} & \alpha_{d_{1}}^{[1]} \alpha_{2}^{[2]} & \cdots & \alpha_{d_{1}}^{[1]} \alpha_{d_{2}}^{[2]}
\end{array}\right]=\left[\begin{array}{cccc}
\beta & \alpha_{1}^{[1]} \alpha_{2}^{[2]} & \cdots & \alpha_{1}^{[1]} \alpha_{d_{2}}^{[2]} \\
0 & \beta & \cdots & \alpha_{2}^{[1]} \alpha_{d_{2}}^{[2]} \\
\vdots & \vdots & & \vdots \\
\alpha_{d_{1}}^{[1]} \alpha_{1}^{[2]} & \alpha_{d_{1}}^{[1]} \alpha_{2}^{[2]} & \cdots & \alpha_{d_{1}}^{[1]} \alpha_{d_{2}}^{[2]}
\end{array}\right] .
$$

This yields that $\alpha_{2}^{[1]} \alpha_{1}^{[2]}=0$ and $\alpha_{1}^{[1]} \alpha_{1}^{[2]}=\alpha_{2}^{[1]} \alpha_{2}^{[2]}$. Since $\left(\alpha_{1}^{[1]} \alpha_{1}^{[2]}\right)^{2}=\left(\alpha_{1}^{[1]} \alpha_{1}^{[2]}\right)\left(\alpha_{2}^{[1]} \alpha_{2}^{[2]}\right)=$ 0 , it turns out that $\left\langle e_{1}^{[1]} \otimes e_{1}^{[2]} \mid u\right\rangle=\alpha_{1}^{[1]} \alpha_{1}^{[2]}=0$. That is, $u$ is orthogonal to $e_{1}^{[1]} \otimes e_{1}^{[2]}$. However, the vector $e_{1}^{[1]} \otimes e_{1}^{[2]}+e_{2}^{[1]} \otimes e_{2}^{[2]} \in \mathcal{L}^{[1: 2]}$ is not orthogonal to $e_{1}^{[1]} \otimes e_{1}^{[2]}$, which implies that $\mathcal{L}^{[1: 2]}$ has no product basis.

Next, consider the case $n \geq 3$. We show that the ( $\tilde{d}-2)$-dimensional subspace

$$
\mathcal{L}=\mathcal{L}^{[1: 2]} \otimes \operatorname{span}\left(u_{0}^{[3: n]}\right)+\mathcal{H}\left(d_{1}, d_{2}\right) \otimes \operatorname{span}\left(u_{0}^{[3: n]}\right)^{\perp}
$$

has no product basis, where $u_{0}^{[3: n]}:=e_{1}^{[3]} \otimes \cdots \otimes e_{1}^{[n]}$. Take an arbitrary product vector $u \in \mathcal{L}$. Then $u$ is expressed as $u=u^{[1: 2]} \otimes u_{0}^{[3: n]}+v^{[1: 2]} \otimes v^{[3: n]}$ with suitable vectors $u^{[1: 2]} \in \mathcal{L}^{[1: 2]}, v^{[1: 2]} \in \mathcal{H}\left(d_{1}, d_{2}\right)$ and $v^{[3: n]} \in \operatorname{span}\left(u_{0}^{[3: n]}\right)^{\perp}$. Since $u$ is a product vector, so is $\left(I^{[1: 2]} \otimes\left\langle u_{0}^{[3: n]}\right|\right) u=u^{[1: 2]}$, where $I^{[1: 2]}$ denotes the identity matrix on $\mathcal{H}\left(d_{1}, d_{2}\right)$. As already proved, the product vector $u^{[1: 2]} \in \mathcal{L}^{[1: 2]}$ is orthogonal to $e_{1}^{[1]} \otimes e_{1}^{[2]}$. Thus, $u$ is orthogonal to $e_{1}^{[1]} \otimes e_{1}^{[2]} \otimes u_{0}^{[3: n]}$. However, the vector $\left(e_{1}^{[1]} \otimes e_{1}^{[2]}+e_{2}^{[1]} \otimes e_{2}^{[2]}\right) \otimes u_{0}^{[3: n]} \in \mathcal{L}$ is not orthogonal to $e_{1}^{[1]} \otimes e_{1}^{[2]} \otimes u_{0}^{[3: n]}$, which implies that $\mathcal{L}$ has no product basis.

Proof of Theorem 1.1. Proposition 3.1 with $r=0,1$ and Proposition 3.4 yield the theorem immediately.

Remark 3.5. Let us consider the case when the scalar field $\mathbb{C}$ and the Hermitian inner product $\langle v \mid u\rangle=\sum_{i} \overline{v(i)} u(i)$ are replaced with an arbitrary field $\mathcal{F}$ and the non-degenerate bilinear form $\langle v, u\rangle=\sum_{i} v(i) u(i)$, respectively. In this case, the proof of Proposition 3.4 works well. Moreover, if $\mathcal{F}$ is infinite, then the proof of Proposition 3.1 also works well because (i) $\operatorname{dim} \mathcal{L}+\operatorname{dim} \mathcal{L}^{\perp}=d$ and $\left(\mathcal{L}^{\perp}\right)^{\perp}=\mathcal{L}$ for every subspace $\mathcal{L}$ of $\mathcal{F}^{d}$ and (ii) Lemma 3.2 holds. The fact (i) is also true for every finite field $\mathcal{F}$, but (ii) is false for every finite field even if the polynomial $f\left(x_{1}, \ldots, x_{n}\right)$ is homogeneous. Nevertheless, a modified version of Lemma 3.2 holds for every finite field (see Lemma 3.10).

Finally, we verify that statement S implies statement WS, which follows from the following proposition immediately.

Proposition 3.6. For a subspace $\mathcal{L}$ of $\mathcal{H}\left(d_{1}, \ldots, d_{n}\right)$, consider the following conditions:

1. the orthogonal projection $P_{\mathcal{L}}$ onto $\mathcal{L}$ lies in $\operatorname{Sep}\left(d_{1}, \ldots, d_{n}\right)$;
2. $\mathcal{L}$ has a product basis.

The one direction " $1 \Rightarrow 2$ " holds for all integers $n \geq 2$ and $d_{1}, \ldots, d_{n} \geq 2$ and subspaces $\mathcal{L}$, but the converse does not necessarily hold.

Proof. $1 \Rightarrow 2$. See [32, Theorem 2] (which is only the case $n=2$, but the case $n \geq 3$ is also proved in the same way).
$2 \nRightarrow 1$. Let $n \geq 2$ and $d_{1}, \ldots, d_{n} \geq 2$ be integers. Choose $\mathcal{L}$ as the subspace spanned by the two vectors $e_{1}^{[1]} \otimes e_{1}^{[2]} \otimes u^{[3: n]}$ and $\left(e_{1}^{[1]}+e_{2}^{[1]}\right) \otimes\left(e_{1}^{[2]}+e_{2}^{[2]}\right) \otimes u^{[3: n]}$, where $u^{[3: n]}$ be an arbitrary unit product vector in $\mathcal{H}\left(d_{3}, \cdots, d_{n}\right)$. Then $\mathcal{L}$ has the
product basis composed of $e_{1}^{[1]} \otimes e_{1}^{[2]} \otimes u^{[3: n]}$ and $\left(e_{1}^{[1]}+e_{2}^{[1]}\right) \otimes\left(e_{1}^{[2]}+e_{2}^{[2]}\right) \otimes u^{[3: n]}$. Also, the orthogonal projection $P_{\mathcal{L}}$ is equal to

$$
P_{\mathcal{L}}=\left(\left|e_{1}^{[1]}\right\rangle\left\langle e_{1}^{[1]}\right| \otimes\left|e_{1}^{[2]}\right\rangle\left\langle e_{1}^{[2]}\right|+\left|u^{[1: 2]}\right\rangle\left\langle u^{[1: 2]}\right|\right) \otimes\left|u^{[3: n]}\right\rangle\left\langle u^{[3: n]}\right|,
$$

where $u^{[1: 2]}$ is the unit vector $\left(e_{1}^{[1]} \otimes e_{2}^{[2]}+e_{2}^{[1]} \otimes e_{1}^{[2]}+e_{2}^{[1]} \otimes e_{2}^{[2]}\right) / \sqrt{3}$. Since the matrix $\left|e_{1}^{[1]}\right\rangle\left\langle e_{1}^{[1]}\right| \otimes\left|e_{1}^{[2]}\right\rangle\left\langle e_{1}^{[2]}\right|+\left|u^{[1: 2]}\right\rangle\left\langle u^{[1: 2]}\right|$ is not separable (we can use the positive partial transpose criterion), the orthogonal projection $P_{\mathcal{L}}$ is not also separable.

### 3.2 Case of finite fields

As already stated in Remark 3.5, Theorem 1.1 holds for every infinite field. In this section, we consider the case of finite fields. Let $\mathcal{F}$ be a finite field of order $q,\langle\cdot, \cdot\rangle$ be the non-degenerate bilinear form $\langle v, u\rangle=\sum_{i} v(i) u(i), \tilde{d}$ be the dimension $d_{1} d_{2} \cdots d_{n}$ of $\mathcal{F}^{d_{1}} \otimes \cdots \otimes \mathcal{F}^{d_{n}}$, and $\left(e_{i}^{[j]}\right)_{i=1}^{d_{j}}$ be the standard basis of $\mathcal{F}^{d_{j}}$ for $j=1, \ldots, n$. We denote by $\mathcal{L}^{\perp}$ the orthogonal compliment of a subspace $\mathcal{L}$ with respect to the non-degenerate bilinear form $\langle\cdot, \cdot\rangle$.

Proposition 3.7. Let $d_{1}, d_{2} \geq 2$ be integers. Then every $\left(d_{1} d_{2}-1\right)$-dimensional subspace of $\mathcal{F}^{d_{1}} \otimes \mathcal{F}^{d_{2}}$ has a product basis.

Proof. Let $\mathcal{L}$ be a $\left(d_{1} d_{2}-1\right)$-dimensional subspace of $\mathcal{F}^{d_{1}} \otimes \mathcal{F}^{d_{2}}$. Taking a non-zero $w \in \mathcal{L}^{\perp}$, we have $\mathcal{L}=\left\{u \in \mathcal{F}^{d_{1}} \otimes \mathcal{F}^{d_{2}}:\langle w, u\rangle=0\right\}$.

Step 1. Let us consider the case $w=w_{r}:=\sum_{i=1}^{r} e_{i}^{[1]} \otimes e_{i}^{[2]}$ with $1 \leq r \leq$ $\min \left\{d_{1}, d_{2}\right\}$. Set $u_{0}^{[2]}=\sum_{i=1}^{r} e_{i}^{[2]}$. In this case, the $d_{1} d_{2}-1$ product vectors of $\mathcal{L}$

$$
\begin{equation*}
e_{i}^{[1]} \otimes e_{j}^{[2]}, \quad\left(e_{k}^{[1]}-e_{k+1}^{[1]}\right) \otimes u_{0}^{[2]} \quad((i, j) \neq(1,1),(2,2), \ldots,(r, r), 1 \leq k \leq r-1) \tag{3.5}
\end{equation*}
$$

are linearly independent, which is proved as follows. First, it is easily checked that all the vectors (3.5) are orthogonal to $w_{r}$, where we note that if $r=1$, then the vectors 3.5 are only $e_{i}^{[1]} \otimes e_{j}^{[2]},(i, j) \neq(1,1)$. Next, suppose that $d_{1} d_{2}-1$ scalars $\alpha_{i, j}$ and $\beta_{k}$ satisfy

$$
\begin{equation*}
\sum_{(i, j) \neq(1,1),(2,2), \ldots,(r, r)} \alpha_{i, j} e_{i}^{[1]} \otimes e_{j}^{[2]}+\sum_{k=1}^{r-1} \beta_{k}\left(e_{k}^{[1]}-e_{k+1}^{[1]}\right) \otimes u_{0}^{[2]}=0 \tag{3.6}
\end{equation*}
$$

Taking the inner product of (3.6) and $e_{l}^{[1]} \otimes e_{l}^{[2]}$ for $l=1, \ldots, r$, we obtain that $\beta_{1}=0$ and $\beta_{l}-\beta_{l-1}=0$ for all $l=2, \ldots, r$. Thus, all $\beta_{k}$ are zero. Since the vectors
$e_{i}^{[1]} \otimes e_{j}^{[2]},(i, j) \neq(1,1),(2,2), \ldots,(r, r)$, are linearly independent, all $\alpha_{i, j}$ are also zero. Therefore, the vectors (3.5) are linearly independent, and $\mathcal{L}$ has a product basis.

Step 2. Let us reduce the case of general $w$ to Step 1. For a matrix $A=\left(\alpha_{i, j}\right) \in$ $\mathcal{F}^{d_{1} \times d_{2}}$, define the vector $\operatorname{vec}(A) \in \mathcal{F}^{d_{1}} \otimes \mathcal{F}^{d_{2}}$ as $\operatorname{vec}(A)=\sum_{i, j} \alpha_{i, j} e_{i}^{[1]} \otimes e_{j}^{[2]}$. Then $\operatorname{vec}\left(P A Q^{\top}\right)=(P \otimes Q) \operatorname{vec}(A)$ for all matrices $A \in \mathcal{F}^{d_{1} \times d_{2}}, P \in \mathcal{F}^{d_{1} \times d_{1}}$ and $Q \in$ $\mathcal{F}^{d_{2} \times d_{2}}$, where $Q^{\top}$ denotes the transpose of $Q$. Now, express $w$ as $w=\operatorname{vec}(A)$ with $A \in \mathcal{F}^{d_{1} \times d_{2}}$. For $r=0,1, \ldots, \min \left\{d_{1}, d_{2}\right\}$, define the matrix $B_{r}=\left(\beta_{i, j}\right) \in \mathcal{F}^{d_{1} \times d_{2}}$ as $\beta_{i, j}=1$ if $(i, j)=(1,1),(2,2), \ldots,(r, r)$ and $\beta_{i, j}=0$ otherwise. Since $A$ can be factorized as $A=P^{\top} B_{r} Q$ with an integer $r \in\left[1, \min \left\{d_{1}, d_{2}\right\}\right]$ and invertible matrices $P \in \mathcal{F}^{d_{1} \times d_{1}}$ and $Q \in \mathcal{F}^{d_{2} \times d_{2}}$, it follows that

$$
w=\operatorname{vec}(A)=\left(P^{\top} \otimes Q^{\top}\right) \operatorname{vec}\left(B_{r}\right)=\left(P^{\top} \otimes Q^{\top}\right) w_{r}
$$

Letting $\left(u_{i}\right)_{i=1}^{d_{1} d_{2}-1}$ be the product basis (3.5), we find that $\left(\left(P^{-1} \otimes Q^{-1}\right) u_{i}\right)_{i=1}^{d_{1} d_{2}-1}$ is a product basis of $\mathcal{L}$.

Proposition 3.8. Let $n \geq 3$ and $d_{1}, \ldots, d_{n} \geq 2$ be integers. If $q>\max \left\{d_{i}: i \neq\right.$ $\left.n_{1}, n_{2}\right\}$ for some $n_{1}$ and $n_{2}$, then every $(\tilde{d}-1)$-dimensional subspace of $\mathcal{F}^{d_{1}} \otimes \cdots \otimes \mathcal{F}^{d_{n}}$ has a product basis.

Since Proposition 3.4 holds for every finite field (see Remark 3.5), we obtain the following theorem.

Theorem 3.9. Let $n \geq 2$ and $d_{1}, \ldots, d_{n} \geq 2$ be integers. If either (i) $n=2$ or (ii) $n \geq 3$ and $q>\max \left\{d_{i}: i \neq n_{1}, n_{2}\right\}$ for some $n_{1}$ and $n_{2}$, then

$$
\max \left\{\operatorname{dim} \mathcal{L}: \begin{array}{l}
\mathcal{L} \text { is a subspace of } \mathcal{F}^{d_{1}} \otimes \cdots \otimes \mathcal{F}^{d_{n}} \text { and } \\
\text { has no product basis }
\end{array}\right\}=\tilde{d}-2
$$

To prove Proposition 3.8, we use the following lemmas instead of Lemmas 3.2 and 3.3.

Lemma 3.10. Let $n \geq 1$ be an integer, $f\left(x_{1}, \ldots, x_{n}\right)$ be a polynomial over $\mathcal{F}$, and $d_{i}$ be the degree of $f\left(x_{1}, \ldots, x_{n}\right)$ in $x_{i}$. If $q>\max \left\{d_{1}, \ldots, d_{n}\right\}$, then the following conditions are equivalent:

1. $f\left(\alpha_{1}, \ldots, \alpha_{n}\right)=0$ for all $\alpha_{1}, \ldots, \alpha_{n} \in \mathcal{F}$;
2. $f\left(x_{1}, \ldots, x_{n}\right)=0$ as a polynomial.

Proof. See [34, Theorem 2.19] (where only Lemma 3.2 is proved but the proof works well for Lemma 3.10).

Lemma 3.11. Let $m \in[1, q-1], n \geq 1$ and $d_{1}, \ldots, d_{n} \geq 2$ be integers, and let $\left(u_{k, l}\right)_{l=1}^{\tilde{d}-1}, k=1, \ldots, m$, be $(\tilde{d}-1)$-tuples of linearly independent vectors in $\mathcal{F}^{d_{1}} \otimes$ $\cdots \otimes \mathcal{F}^{d_{n}}$. Then there exists a product vector $v \in \mathcal{F}^{d_{1}} \otimes \cdots \otimes \mathcal{F}^{d_{n}}$ such that

$$
\operatorname{det}\left[u_{k, 1}, \ldots, u_{k, \tilde{d}-1}, v\right] \neq 0
$$

for all $k=1, \ldots, m$.
Proof. This proof is almost the same as the proof of Lemma 3.3. Let $n^{\prime}=d_{1}+\cdots+d_{n}$. Define the $m$ polynomials $f_{k}\left(x_{1}, \ldots, x_{n^{\prime}}\right)$ over $\mathcal{F}$ as

$$
\begin{gathered}
f_{k}\left(x_{1}, \ldots, x_{n^{\prime}}\right)=\operatorname{det}\left[u_{k, 1}, \ldots, u_{k, \tilde{d}-1}, v\right] \quad(k=1, \ldots, m), \\
v=v^{[1]} \otimes \cdots \otimes v^{[n]}
\end{gathered}
$$

where the variables $x_{1}, \ldots, x_{n^{\prime}}$ correspond to the $n^{\prime}$ entries of $v^{[j]}$, $1 \leq j \leq n$. Then, for every $k=1, \ldots, m$, the polynomial $f_{k}\left(x_{1}, \ldots, x_{n^{\prime}}\right)$ is not zero as a polynomial. Since the polynomial ring $\mathcal{F}\left[x_{1}, \ldots, x_{n^{\prime}}\right]$ is an integral domain, the product $f\left(x_{1}, \ldots, x_{n^{\prime}}\right):=\prod_{k=1}^{m} f_{k}\left(x_{1}, \ldots, x_{n^{\prime}}\right)$ is not also zero as a polynomial. Also, for every $i=1, \ldots, n^{\prime}$, the degree of the polynomial $f\left(x_{1}, \ldots, x_{n^{\prime}}\right)$ in $x_{i}$ is less than or equal to $m \leq q-1$. Thus, Lemma 3.10 implies that $f\left(\beta_{1}, \ldots, \beta_{n^{\prime}}\right) \neq 0$ for some $\beta_{1}, \ldots, \beta_{n^{\prime}} \in \mathcal{F}$. Taking the vectors $v^{[j]}$ corresponding to $\beta_{1}, \ldots, \beta_{n^{\prime}} \in \mathcal{F}$, we obtain a desired product vector $v$.

Proof of Proposition 3.8. Without loss of generality, we may assume that $2 \leq d_{1} \leq$ $\cdots \leq d_{n}$ and $q>d_{n-2}\left(\left\{n_{1}, n_{2}\right\}=\{n-1, n\}\right.$ in this case $)$. The proposition is proved in the same way as the proof of Proposition 3.1 by using Lemmas 3.10 and 3.11 instead of Lemmas 3.2 and 3.3 .

## Chapter 4

## Classical-quantum differential privacy

In this chapter, we introduce linear mappings used in quantum information theory, prove Theorems 1.5 and 1.6 , and give concrete $\mathrm{CQ} \varepsilon$-DP $n$-tuples that do not lie in $\mathrm{EC}_{n}(\varepsilon)$.

### 4.1 Linear mappings in quantum information theory

In this section, we discuss linear mappings from $\operatorname{Herm}(d)$ into Herm $\left(d^{\prime}\right)$ which are used in quantum information theory. First, let us begin with several basic terms (see also a textbook in quantum information theory, e.g., [26,27]). For two linear mapping $\Lambda_{i}$, $i=1,2$, from $\operatorname{Herm}\left(d_{i}\right)$ into $\operatorname{Herm}\left(d_{i}^{\prime}\right)$, the tensor product $\Lambda_{1} \otimes \Lambda_{2}$ is a linear mapping from $\operatorname{Herm}\left(d_{1}\right) \otimes \operatorname{Herm}\left(d_{2}\right)$ into $\operatorname{Herm}\left(d_{1}^{\prime}\right) \otimes \operatorname{Herm}\left(d_{2}^{\prime}\right)$. Since $\operatorname{Herm}\left(d_{1}\right) \otimes \operatorname{Herm}\left(d_{2}\right)$ can be regarded as $\operatorname{Herm}\left(d_{1} d_{2}\right)$ (see Section 2.1), the tensor product $\Lambda_{1} \otimes \Lambda_{2}$ is also a linear mapping from $\operatorname{Herm}\left(d_{1} d_{2}\right)$ into $\operatorname{Herm}\left(d_{1}^{\prime} d_{2}^{\prime}\right)$. Let $\mathrm{id}_{d}$ be the identity mapping on $\operatorname{Herm}(d)$.

- A linear mapping $\Lambda$ from $\operatorname{Herm}(d)$ into $\operatorname{Herm}\left(d^{\prime}\right)$ is called positive if $\Lambda(\operatorname{PSD}(d)) \subset$ $\operatorname{PSD}\left(d^{\prime}\right)$.
- A linear mapping $\Lambda$ from $\operatorname{Herm}(d)$ into $\operatorname{Herm}\left(d^{\prime}\right)$ is called completely positive if $\Lambda \otimes \mathrm{id}_{k}$ is positive for every integer $k \geq 2$.
- A linear mapping $\Lambda$ from Herm $(d)$ into $\operatorname{Herm}\left(d^{\prime}\right)$ is called trace-preserving if $\operatorname{Tr} \Lambda(X)=\operatorname{Tr} X$ for every $X \in \operatorname{Herm}(d)$.
- A linear mapping $\Lambda$ from $\operatorname{Herm}(d)$ into $\operatorname{Herm}\left(d^{\prime}\right)$ is called CPTP if $\Lambda$ is completely positive and trace-preserving.

In quantum information theory, a quantum channel is a CPTP map.
Example 4.1 (Entanglement breaking channel). The linear mapping $\Lambda$ below is a CPTP map called entanglement breaking channel. Let $\sigma_{1}, \ldots, \sigma_{m}$ be density matrices on $\mathbb{C}^{d^{\prime}}$ and $\left(M_{k}\right)_{k=1}^{m}$ be a POVM, i.e., $M_{1}, \ldots, M_{m} \geq 0$ and $\sum_{k=1}^{m} M_{k}=I_{d}$. For example, $\left(\left|e_{k}\right\rangle\left\langle e_{k}\right|\right)_{k=1}^{d}$ is a POVM, where $\left(e_{k}\right)_{k=1}^{d}$ denotes the standard basis of $\mathbb{C}^{d}$. Define the linear mapping $\Lambda$ from $\operatorname{Herm}(d)$ into $\operatorname{Herm}\left(d^{\prime}\right)$ as $\Lambda(X)=\sum_{k=1}^{m}\left(\operatorname{Tr} M_{k} X\right) \sigma_{k}$. It can easily be checked that $\Lambda$ is a CPTP map. This fact is used implicitly in this section.

Using an entanglement breaking channel, we prove Proposition 1.7 ,
Proof of Proposition 1.7. First, assume that $\left(\rho_{i}\right)_{i=1}^{n}$ lies in $\mathrm{EC}_{n}(\varepsilon)$, i.e., there exists an $\varepsilon$-DP $n$-tuple $\left(p_{i}\right)_{i=1}^{n} \in \mathrm{C}_{n}(\varepsilon)$ and a CPTP map $\Lambda$ such that $\rho_{i}=\Lambda\left(\operatorname{diag}\left(p_{i}\right)\right)$ for all $i=1, \ldots, n$. Denote by $\left(e_{k}\right)_{k=1}^{d}$ the standard basis of $\mathbb{C}^{d}$, where $d$ is the dimension of the vector space that $p_{1}, \ldots, p_{n}$ inhabit. Then $\rho_{i}=\sum_{k=1}^{d} p_{i}(k) \Lambda\left(\left|e_{k}\right\rangle\left\langle e_{k}\right|\right)$ for all $i=1, \ldots, n$. Since all $\Lambda\left(\left|e_{k}\right\rangle\left\langle e_{k}\right|\right)$ are density matrices, the $n$-tuple $\left(\rho_{i}\right)_{i=1}^{n}$ lies in the right-hand side of (1.7).

Conversely, assume that $\left(\rho_{i}\right)_{i=1}^{n}$ lies in the right-hand side of (1.7): there exist $\left(p_{i}\right)_{i=1}^{n} \in \mathrm{C}_{n}(\varepsilon)$ and density matrices $\sigma_{k}$ such that $\rho_{i}=\sum_{k=1}^{d} p_{i}(k) \sigma_{k}$ for all $i=$ $1, \ldots, n$, where $d$ is the dimension of the vector space that $p_{1}, \ldots, p_{n}$ inhabit. Define the CPTP map $\Lambda$ as $\Lambda(X)=\sum_{k=1}^{d}\left\langle e_{k}\right| X\left|e_{k}\right\rangle \sigma_{k}$. Then $\Lambda\left(\operatorname{diag}\left(p_{i}\right)\right)=\rho_{i}$ for all $i=1, \ldots, n$. Therefore, $\left(\rho_{i}\right)_{i=1}^{n}$ lies in $\mathrm{EC}_{n}(\varepsilon)$.

### 4.2 Case $n=2$

In this section, we show the following theorem that is stronger than Theorem 1.5 .
Theorem 4.2. Let $\bar{p}_{1}$ and $\bar{p}_{2}$ be the probability vectors defined as

$$
\bar{p}_{1}=\left[\frac{e^{\varepsilon}}{e^{\varepsilon}+1}, \frac{1}{e^{\varepsilon}+1}\right], \quad \bar{p}_{2}=\left[\frac{1}{e^{\varepsilon}+1}, \frac{e^{\varepsilon}}{e^{\varepsilon}+1}\right] .
$$

Then, for every $\left(\rho_{i}\right)_{i=1}^{2} \in \mathrm{CQ}_{2}(\varepsilon)$, there exists a CPTP map $\Lambda$ such that $\Lambda\left(\operatorname{diag}\left(\bar{p}_{i}\right)\right)=$ $\rho_{i}$ for $i=1,2$.

That is, the pair $\left(\bar{p}_{1}, \bar{p}_{2}\right)$ generates all CQ $\varepsilon$-DP 2-tuples by CPTP maps. The pair $\left(\bar{p}_{1}, \bar{p}_{2}\right)$ also plays an important role in the classical setting [30,38]. Since $\mathrm{EC}_{2}(\varepsilon)$ is a subset of $\mathrm{CQ}_{2}(\varepsilon)$, Theorem 1.5 follows from Theorem 4.2 and Proposition 1.7 immediately.

In general, many researchers in quantum information theory are interested in whether, given $2 n$ density matrices $\rho_{1}, \ldots, \rho_{n}, \sigma_{1}, \ldots, \sigma_{n}$, there exists a CPTP map $\Lambda$ such that $\Lambda\left(\rho_{i}\right)=\sigma_{i}$ for all $i=1, \ldots, n$. In this context, the CPTP map $\Lambda$ is called a physical transformation [9,33]. If the above $\Lambda$ exists, we write $\left(\rho_{1}, \ldots, \rho_{n}\right) \rightarrow$ $\left(\sigma_{1}, \ldots, \sigma_{n}\right)$. Then the assertion of Theorem 4.2 can be rewritten as

$$
\forall\left(\rho_{i}\right)_{i=1}^{2} \in \mathrm{CQ}_{2}(\varepsilon), \quad\left(\operatorname{diag}\left(\bar{p}_{i}\right)\right)_{i=1}^{2} \rightarrow\left(\rho_{i}\right)_{i=1}^{2} .
$$

If inputs $\rho_{1}$ and $\rho_{2}$ are of rank one, it is easy to investigate whether $\left(\rho_{1}, \rho_{2}\right) \rightarrow\left(\sigma_{1}, \sigma_{2}\right)$, because $\left(\rho_{1}, \rho_{2}\right) \rightarrow\left(\sigma_{1}, \sigma_{2}\right)$ if and only if $F\left(\rho_{1}, \rho_{2}\right) \leq F\left(\sigma_{1}, \sigma_{2}\right)$ [9, 33, 49], where $F$ denotes the fidelity defined as $F(\rho, \sigma)=\operatorname{Tr}\left|\rho^{1 / 2} \sigma^{1 / 2}\right|$ for density matrices $\rho$ and $\sigma$. However, if inputs $\rho_{1}$ and $\rho_{2}$ are not necessarily of rank one, it is difficult to investigate whether $\left(\rho_{1}, \rho_{2}\right) \rightarrow\left(\sigma_{1}, \sigma_{2}\right)$ in general. Since both the inputs in Theorem 4.2 are of rank two, Theorem 4.2 is never trivial.

Now, we show the following preliminary lemma before proving Theorem 4.2.
Lemma 4.3. If density matrices $\rho_{1}$ and $\rho_{2}$ are orthogonal to each other, i.e., $\operatorname{Tr} \rho_{1} \rho_{2}=$ 0 , then for all density matrices $\sigma_{1}$ and $\sigma_{2}$ there exists a CPTP map $\Lambda$ such that $\Lambda\left(\rho_{1}\right)=\sigma_{1}$ and $\Lambda\left(\rho_{2}\right)=\sigma_{2}$.

Proof. Take the orthogonal projection $P_{1}$ onto the support of $\rho_{1}$. Put $P_{2}=I_{d}-P_{1}$. Defining $\Lambda(X)=\left(\operatorname{Tr} X P_{1}\right) \sigma_{1}+\left(\operatorname{Tr} X P_{2}\right) \sigma_{2}$ for $X \in \operatorname{Herm}(d)$, we find that $\Lambda$ is a CPTP map satisfying that $\Lambda\left(\rho_{1}\right)=\sigma_{1}$ and $\Lambda\left(\rho_{2}\right)=\sigma_{2}$.

Proof of Theorem 4.2. Let $\left(\rho_{i}\right)_{i=1}^{2}$ be a CQ $\varepsilon$-DP 2-tuple. Since $e^{\varepsilon} \rho_{1}-\rho_{2}$ and $e^{\varepsilon} \rho_{2}-\rho_{1}$ are positive semi-definite, they are expressed as

$$
e^{\varepsilon} \rho_{1}-\rho_{2}=\left(e^{\varepsilon}-1\right) \sigma_{1}, \quad e^{\varepsilon} \rho_{2}-\rho_{1}=\left(e^{\varepsilon}-1\right) \sigma_{2}
$$

where $\sigma_{1}$ and $\sigma_{2}$ are density matrices. Lemma 4.3 implies that there exists a CPTP map $\Lambda$ such that $\Lambda\left(\left|e_{1}\right\rangle\left\langle e_{1}\right|\right)=\sigma_{1}$ and $\Lambda\left(\left|e_{2}\right\rangle\left\langle e_{2}\right|\right)=\sigma_{2}$. Thus,

$$
\begin{gather*}
e^{\varepsilon} \operatorname{diag}\left(\bar{p}_{1}\right)-\operatorname{diag}\left(\bar{p}_{2}\right)=\left(e^{\varepsilon}-1\right)\left|e_{1}\right\rangle\left\langle e_{1}\right|, \\
e^{\varepsilon} \operatorname{diag}\left(\bar{p}_{2}\right)-\operatorname{diag}\left(\bar{p}_{1}\right)=\left(e^{\varepsilon}-1\right)\left|e_{2}\right\rangle\left\langle e_{2}\right|, \\
e^{\varepsilon} \Lambda\left(\operatorname{diag}\left(\bar{p}_{1}\right)\right)-\Lambda\left(\operatorname{diag}\left(\bar{p}_{1}\right)\right)=\left(e^{\varepsilon}-1\right) \Lambda\left(\left|e_{1}\right\rangle\left\langle e_{1}\right|\right)=\left(e^{\varepsilon}-1\right) \sigma_{1}=e^{\varepsilon} \rho_{1}-\rho_{2},  \tag{4.1}\\
e^{\varepsilon} \Lambda\left(\operatorname{diag}\left(\bar{p}_{2}\right)\right)-\Lambda\left(\operatorname{diag}\left(\bar{p}_{1}\right)\right)=\left(e^{\varepsilon}-1\right) \Lambda\left(\left|e_{2}\right\rangle\left\langle e_{2}\right|\right)=\left(e^{\varepsilon}-1\right) \sigma_{2}=e^{\varepsilon} \rho_{2}-\rho_{1} . \tag{4.2}
\end{gather*}
$$

Solving the simultaneous equations (4.1) and (4.2), we obtain $\Lambda\left(\operatorname{diag}\left(\bar{p}_{i}\right)\right)=\rho_{i}$ for all $i=0,1$.

### 4.3 Case $n \geq 3$ : Main idea and result

We describe the main idea and result in the case $n \geq 3$. A key point to distinguish elements in $\mathrm{EC}_{n}(\varepsilon)$ and $\mathrm{CQ}_{n}(\varepsilon)$ is monotonicity for CPTP maps (see Definition 1.3). Let $\Phi$ be a real-valued function of $n$ density matrices satisfying monotonicity for CPTP maps. As stated in Section 1.2 , the equality $S_{n}^{\mathrm{EC}}(\varepsilon ; \Phi)=S_{n}^{\mathrm{C}}(\varepsilon ; \Phi)$ holds. This fact gives us an idea to distinguish elements in $\mathrm{EC}_{n}(\varepsilon)$ and $\mathrm{CQ}_{n}(\varepsilon)$ : if $\Phi\left(\rho_{1}, \ldots, \rho_{n}\right)$ with $\left(\rho_{i}\right)_{i=1}^{n} \in \mathrm{CQ}_{n}(\varepsilon)$ is greater than $S_{n}^{\mathrm{C}}(\varepsilon ; \Phi)$, then $\left(\rho_{i}\right)_{i=1}^{n}$ does not lie in $\mathrm{EC}_{n}(\varepsilon)$. To use this criterion, we must solve optimization problem (1.6). Fortunately, for special $\Phi$, optimization problem (1.6) can be reduced to a linear program [38, Theorem 4], which is stated in Section 4.4.

Before stating our main result, we define the RLD Fisher information of a oneparameter family, which satisfies monotonicity for CPTP maps.

Definition 4.4 (RLD Fisher information [26, p. 260]). For density matrices $\rho$ and $\sigma$ with full rank, we denote the RLD Fisher information of the one-parameter family $((1-\theta) \rho+\theta \sigma)_{\theta \in[0,1]}$ at the point $\theta$ as

$$
J_{\theta}(\rho, \sigma)=\operatorname{Tr}(\sigma-\rho)^{2}((1-\theta) \rho+\theta \sigma)^{-1}
$$

For probability vectors $p$ and $q$, we set $J_{\theta}(p, q)=J_{\theta}(\operatorname{diag}(p), \operatorname{diag}(q))$.
If $\left(\rho_{i}\right)_{i=1}^{n}$ is CQ $\varepsilon$-DP, we may assume that all $\rho_{i}$ have full rank (see Section 1.2), and hence, we can consider the value $J_{\theta}\left(\rho_{i}, \rho_{j}\right)$ for all $i, j=1, \ldots, n$. Also, for probability vectors $p$ and $q$, the value $J_{\theta}(p, q)$ is the Fisher information in the classical sense. From now on, we denote by $\operatorname{avg}_{i \neq j} \alpha_{i, j}$ the arithmetic mean of real numbers $\alpha_{i, j}, i \neq j$. Now, our main result is as follows.

Theorem 4.5. For real numbers $\theta \in[0,1]$ and $\varepsilon>0$ and an integer $n \geq 2$, we define the suprema $M_{n}^{\mathrm{C}}\left(\varepsilon ; J_{\theta}\right), M_{n}^{\mathrm{EC}}\left(\varepsilon ; J_{\theta}\right)$ and $M_{n}^{\mathrm{CQ}}\left(\varepsilon ; J_{\theta}\right)$ as

$$
\begin{aligned}
& M_{n}^{\mathrm{C}}\left(\varepsilon ; J_{\theta}\right)=\sup _{\left(p_{i}\right)_{i=1}^{n} \in \mathrm{C}_{n}(\varepsilon)} \min _{i \neq j} J_{\theta}\left(p_{i}, p_{j}\right), \\
& M_{n}^{\mathrm{X}}\left(\varepsilon ; J_{\theta}\right)=\sup _{\left(\rho_{i}\right)_{i=1}^{n} \in \mathrm{X}_{n}(\varepsilon)} \min _{i \neq j} J_{\theta}\left(\rho_{i}, \rho_{j}\right) \quad(\mathrm{X}=\mathrm{EC}, \mathrm{CQ}) .
\end{aligned}
$$

Then, for all $\theta \in[0,1], \varepsilon>0$ and $n \geq 2$, we have $M_{n}^{\mathrm{CQ}}\left(\varepsilon ; J_{\theta}\right)=M_{2}^{\mathrm{C}}\left(\varepsilon ; J_{\theta}\right)$ and

$$
\begin{aligned}
M_{n}^{\mathrm{EC}}\left(\varepsilon ; J_{\theta}\right) & =M_{n}^{\mathrm{C}}\left(\varepsilon ; J_{\theta}\right)=\sup _{\left(p_{i}\right)_{i=1}^{\mathrm{E}} \in \mathrm{C}_{n}(\varepsilon)} \operatorname{avg}_{i \neq j} J_{\theta}\left(p_{i}, p_{j}\right) \\
& =\frac{f_{\theta}\left(e^{\varepsilon}, 1\right)+f_{\theta}\left(1, e^{\varepsilon}\right)}{n-1} \max _{1 \leq k \leq n / 2} \frac{k(n-k)}{k e^{\varepsilon}+n-k},
\end{aligned}
$$

where $f_{\theta}(\alpha, \beta):=(\alpha-\beta)^{2} /((1-\theta) \alpha+\theta \beta)$ for $\alpha, \beta>0$. Moreover, $M_{n}^{\mathrm{EC}}\left(\varepsilon ; J_{\theta}\right)<$ $M_{n}^{\mathrm{CQ}}\left(\varepsilon ; J_{\theta}\right)$ for all $\theta \in[0,1], \varepsilon>0$ and $n \geq 3$.

Theorem 4.5 implies Theorem 1.6 and the following corollary immediately.
Corollary 4.6. Let $\varepsilon>0$ be a real number and $n \geq 3$ be an integer. If $\left(\rho_{i}\right)_{i=1}^{n} \in$ $\mathrm{CQ}_{n}(\varepsilon)$ satisfies that $M_{n}^{\mathrm{C}}\left(\varepsilon ; J_{\theta}\right)<\operatorname{avg}_{i \neq j} J_{\theta}\left(\rho_{i}, \rho_{j}\right)$ for some $\theta \in[0,1]$, then $\left(\rho_{i}\right)_{i=1}^{n}$ does not lie in $\mathrm{EC}_{n}(\varepsilon)$.

### 4.4 Case $n \geq 3$ : Proof

First, we begin with the classical optimization, for which we need the following definition and lemma [38, Theorem 4].
Definition 4.7 (Sublinear function). We say that a function $\phi:(0, \infty)^{n} \rightarrow \mathbb{R}$ is sublinear if $\phi(x+y) \leq \phi(x)+\phi(y)$ and $\phi(\alpha x)=\alpha \phi(x)$ for all $x, y \in(0, \infty)^{n}$ and $\alpha>0$.

Lemma 4.8. Let $\Phi_{\mathrm{C}}$ be a real-valued function of $n$ probability vectors with the following condition: there exists a sublinear function $\phi:(0, \infty)^{n} \rightarrow \mathbb{R}$ such that

$$
\begin{equation*}
\Phi_{\mathrm{C}}\left(p_{1}, \ldots, p_{n}\right)=\sum_{p_{1}(k), \ldots, p_{n}(k)>0} \phi\left(p_{1}(k), \ldots, p_{n}(k)\right), \tag{4.3}
\end{equation*}
$$

where the above sum is taken all over $k$ with $p_{1}(k), \ldots, p_{n}(k)>0$. Then, for all $\varepsilon>0$ and $n \geq 2$,

$$
\begin{aligned}
& \sup _{\left(p_{i}\right)_{i=1}^{n} \in \mathrm{C}_{n}(\varepsilon)} \Phi_{\mathrm{C}}\left(p_{1}, \ldots, p_{n}\right) \\
= & \max \left\{\sum_{v \in \mathcal{S}_{n}(\varepsilon)} \phi(v(1), \ldots, v(n)) \alpha_{v}: \begin{array}{l}
\sum_{v \in \mathcal{S}_{n}(\varepsilon)} \alpha_{v} v=\mathbf{1}_{n}, \\
\forall v \in \mathcal{S}_{n}(\varepsilon), \alpha_{v} \geq 0
\end{array}\right\},
\end{aligned}
$$

where $\mathcal{S}_{n}(\varepsilon):=\left\{1, e^{\varepsilon}\right\}^{n}$ and $\mathbf{1}_{n}:=[1, \ldots, 1]^{\top} \in \mathbb{R}^{n}$.
Many information-theoretic quantities can be expressed as 4.3). Such examples are relative entropy, Fisher information, total variation distance. Especially, $J_{\theta}(p, q)$ is expressed as

$$
J_{\theta}(p, q)=\sum_{p(k), q(k)>0} f_{\theta}(p(k), q(k)),
$$

where the above sum is taken all over $k$ with $p(k), q(k)>0$, and the function $f_{\theta}$ defined in Theorem 4.5 is sublinear. We now prove the following lemma by using Lemma 4.8.

Lemma 4.9. Let $\psi:(0, \infty)^{2} \rightarrow \mathbb{R}$ be a sublinear function with $\psi(1,1)=0$, and $\Psi$ be the function $\Psi(p, q)=\sum_{p(k), q(k)>0} \psi(p(k), q(k))$ of two probability vectors. Then, for all $\varepsilon>0$ and $n \geq 2$,

$$
\begin{aligned}
M_{n}^{\mathrm{C}}(\varepsilon ; \Psi) & :=\sup _{\left(p_{i}\right)_{i=1}^{n} \in \mathrm{C}_{n}(\varepsilon)} \min _{i \neq j} \Psi\left(p_{i}, p_{j}\right)=\sup _{\left(p_{i}\right)_{i=1}^{n} \in \mathrm{C}_{n}(\varepsilon)} \operatorname{avg} \Psi\left(p_{i}, p_{j}\right) \\
& =\frac{\psi\left(e^{\varepsilon}, 1\right)+\psi\left(1, e^{\varepsilon}\right)}{n-1} \max _{1 \leq k \leq n / 2} \frac{k(n-k)}{k e^{\varepsilon}+n-k} .
\end{aligned}
$$

Proof. Let $\varepsilon>0$ be a real number and $n \geq 2$ be an integer. The following inequality holds:

$$
\begin{equation*}
M_{n}^{\mathrm{C}}(\varepsilon ; \Psi) \leq \sup _{\left(p_{i}\right)_{i=1}^{n} \in \mathrm{C}_{n}(\varepsilon)} \operatorname{avg}_{i \neq j}^{\operatorname{avg}} \Psi\left(p_{i}, p_{j}\right) . \tag{4.4}
\end{equation*}
$$

Recall the definition $\mathcal{S}_{n}(\varepsilon)=\left\{1, e^{\varepsilon}\right\}^{n}$. Set $\Phi_{\mathrm{C}}\left(p_{1}, \ldots, p_{n}\right)=\sum_{i \neq j} \Psi\left(p_{i}, p_{j}\right)$ and $\phi(x)=\sum_{i \neq j} \psi(x(i), x(j))$ for $x \in(0, \infty)^{n}$. Lemma 4.8 yields that

$$
\left.\begin{array}{rl} 
& \sup _{\left(p_{i}\right)_{i=1}^{n} \in \mathrm{C}_{n}(\varepsilon)} \sum_{i \neq j} \Psi\left(p_{i}, p_{j}\right) \\
= & \max \left\{\sum_{v \in \mathcal{S}_{n}(\varepsilon)} \sum_{i \neq j} \psi(v(i), v(j)) \alpha_{v}: \sum_{v \in \mathcal{S}_{n}(\varepsilon)} \alpha_{v} v=\mathbf{1}_{n},\right.  \tag{4.5}\\
\forall v \in \mathcal{S}_{n}(\varepsilon), \alpha_{v} \geq 0
\end{array}\right\} . \quad .
$$

Consider the partition of $\mathcal{S}_{n}(\varepsilon)$ into the $n+1$ subsets

$$
\mathcal{S}_{n, k}(\varepsilon):=\left\{v \in \mathcal{S}_{n}(\varepsilon): \text { the number of } i \text { with } v(i)=e^{\varepsilon} \text { is } k\right\} \quad(k=0,1, \ldots, n) .
$$

If $v \in \mathcal{S}_{n, k}(\varepsilon)$, then $\sum_{i \neq j} \psi(v(i), v(j))=\left(\psi\left(e^{\varepsilon}, 1\right)+\psi\left(1, e^{\varepsilon}\right)\right) k(n-k)$ due to the assumption $\psi(1,1)=0$. Thus, for every $k=0,1, \ldots, n$, we have

$$
\begin{aligned}
\sum_{v \in \mathcal{S}_{n}(\varepsilon)} \sum_{i \neq j} \psi(v(i), v(j)) \alpha_{v} & =\sum_{k=0}^{n} \sum_{v \in \mathcal{S}_{n, k}(\varepsilon)} \sum_{i \neq j} \psi(v(i), v(j)) \alpha_{v} \\
& =\sum_{k=0}^{n}\left(\psi\left(e^{\varepsilon}, 1\right)+\psi\left(1, e^{\varepsilon}\right)\right) k(n-k) \beta_{k}
\end{aligned}
$$

where $\beta_{k}=\sum_{v \in \mathcal{S}_{n, k}(\varepsilon)} \alpha_{v}$. Since the equality $\sum_{v \in \mathcal{S}_{n}(\varepsilon)} \alpha_{v} v=\mathbf{1}_{n}$ yields

$$
\sum_{k=0}^{n}\left(k e^{\varepsilon}+n-k\right) \beta_{k}=\sum_{v \in \mathcal{S}_{n}(\varepsilon)} \alpha_{v}\left\langle\mathbf{1}_{n} \mid v\right\rangle=\left\langle\mathbf{1}_{n} \mid \mathbf{1}_{n}\right\rangle=n
$$

the right-hand side in 4.5 is bounded above by

$$
\begin{align*}
& \max \left\{\sum_{k=0}^{n}\left(\psi\left(e^{\varepsilon}, 1\right)+\psi\left(1, e^{\varepsilon}\right)\right) k(n-k) \beta_{k}: \begin{array}{c}
\sum_{k=0}^{n}\left(k e^{\varepsilon}+n-k\right) \beta_{k}=n, \\
\beta_{0}, \ldots, \beta_{n} \geq 0
\end{array}\right\} \\
= & \left(\psi\left(e^{\varepsilon}, 1\right)+\psi\left(1, e^{\varepsilon}\right)\right) n \max _{0 \leq k \leq n} \frac{k(n-k)}{k e^{\varepsilon}+n-k} \\
= & \left(\psi\left(e^{\varepsilon}, 1\right)+\psi\left(1, e^{\varepsilon}\right)\right) n \max _{1 \leq k \leq n / 2} \frac{k(n-k)}{k e^{\varepsilon}+n-k} . \tag{4.6}
\end{align*}
$$

From (4.4), (4.5) and (4.6), it follows that

$$
\begin{equation*}
M_{n}^{\mathrm{C}}(\varepsilon ; \Psi) \leq \sup _{\left(p_{i}\right)_{i=1}^{n} \in \mathrm{C}_{n}(\varepsilon)} \operatorname{avg}_{i \neq j} \Psi\left(p_{i}, p_{j}\right) \leq \frac{\psi\left(e^{\varepsilon}, 1\right)+\psi\left(1, e^{\varepsilon}\right)}{n-1} \max _{1 \leq k \leq n / 2} \frac{k(n-k)}{k e^{\varepsilon}+n-k} \tag{4.7}
\end{equation*}
$$

Fix an arbitrary integer $1 \leq k \leq n / 2$. Let $d$ be the number of elements in $\mathcal{S}_{n, k}(\varepsilon)$, i.e., $d=\binom{n}{k}$. Then the vector space $\mathbb{R}^{\mathcal{S}_{n, k}(\varepsilon)}$ is isomorphic to $\mathbb{R}^{d}$. Define the probability vectors $p_{1}, \ldots, p_{n} \in \mathbb{R}^{\mathcal{S}_{n, k}(\varepsilon)}$ as

$$
p_{i}(v)=\frac{v(i)}{\binom{n-1}{k-1} e^{\varepsilon}+\binom{n-1}{k}} \quad\left(v \in \mathcal{S}_{n, k}(\varepsilon) ; i=1, \ldots, n\right)
$$

Then $\left(p_{i}\right)_{i=1}^{n}$ is $\varepsilon$-DP, and moreover,

$$
\begin{aligned}
& \min _{i \neq j} \Psi\left(p_{i}, p_{j}\right)=\min _{i \neq j} \frac{1}{\binom{n-1}{k-1} e^{\varepsilon}+\binom{n-1}{k}} \sum_{v \in \mathcal{S}_{n, k}(\varepsilon)} \psi(v(i), v(j)) \\
= & \frac{1}{\binom{n-1}{k-1} e^{\varepsilon}+\binom{n-1}{k}} \cdot\left(\psi\left(e^{\varepsilon}, 1\right)+\psi\left(1, e^{\varepsilon}\right)\right)\binom{n-2}{k-1} \\
= & \frac{\psi\left(e^{\varepsilon}, 1\right)+\psi\left(1, e^{\varepsilon}\right)}{((n-1) /(n-k)) e^{\varepsilon}+(n-1) / k}=\frac{\psi\left(e^{\varepsilon}, 1\right)+\psi\left(1, e^{\varepsilon}\right)}{n-1} \cdot \frac{k(n-k)}{k e^{\varepsilon}+n-k} .
\end{aligned}
$$

Since $1 \leq k \leq n / 2$ is arbitrary, the inequalities in 4.7) turn to equality.
Next, we consider the quantum optimization, for which we need the following lemmas.

Lemma 4.10. Let $n \geq 2$ be an integer and $c \in[0,1]$ be a real number. There exists an $n$-tuple $\left(u_{i}\right)_{i=1}^{n}$ of unit vectors in $\mathbb{R}^{n}$ such that $\left\langle u_{i} \mid u_{j}\right\rangle=c$ for all $i \neq j$.
Proof. Since the matrix $A=(1-c) I_{n}+c\left|\mathbf{1}_{n}\right\rangle\left\langle\mathbf{1}_{n}\right|$ is positive semi-definite, there exists a real square matrix $B$ of order $n$ such that $A=B^{\top} B$. The column vectors $u_{1}, \ldots, u_{n} \in \mathbb{R}^{n}$ of $B$ satisfy that $\left\langle u_{i} \mid u_{j}\right\rangle=1$ if $i=j$ and $\left\langle u_{i} \mid u_{j}\right\rangle=c$ if $i \neq j$.

Lemma 4.11. Let $c \in[0,1)$ and $\varepsilon, t>0$ be real numbers, $\left(u_{i}\right)_{i=1}^{n}$ be the $n$-tuple in Lemma 4.10, and $\rho_{1}, \ldots, \rho_{n}$ be the density matrices defined as

$$
\rho_{i}=\frac{1}{n+t}\left(I_{n}+t\left|u_{i}\right\rangle\left\langle u_{i}\right|\right) \quad(i=1, \ldots, n),
$$

where $I_{n}$ denotes the identity matrix of order $n$. If $D=\left(e^{\varepsilon}-1\right)^{2}+4\left(1-c^{2}\right) e^{\varepsilon}$ and

$$
0<t \leq t_{\max }:=\frac{2\left(e^{\varepsilon}-1\right)}{\sqrt{D}+1-e^{\varepsilon}}
$$

then $\left(\rho_{i}\right)_{i=1}^{n}$ is $C Q \varepsilon-D P$.
Proof. Let $i \neq j$. We show that

$$
\begin{equation*}
\left|u_{i}\right\rangle\left\langle u_{i}\right|-e^{\varepsilon}\left|u_{j}\right\rangle\left\langle u_{j}\right| \leq \frac{1-e^{\varepsilon}+\sqrt{D}}{2} I_{n} . \tag{4.8}
\end{equation*}
$$

Take an orthonormal system $\left(e_{i}\right)_{i=1}^{2}$ of $\mathbb{C}^{n}$ such that $u_{i}=e_{1}, u_{j}=\alpha e_{1}+\beta e_{2}, \alpha=c$ and $\beta=\sqrt{1-\alpha^{2}}$. Then the matrix $\left|u_{i}\right\rangle\left\langle u_{i}\right|-e^{\varepsilon}\left|u_{j}\right\rangle\left\langle u_{j}\right|$ can be expressed as a square matrix of order 2:

$$
\left[\begin{array}{ll}
1 & 0  \tag{4.9}\\
0 & 0
\end{array}\right]-e^{\varepsilon}\left[\begin{array}{cc}
\alpha^{2} & \alpha \beta \\
\alpha \beta & \beta^{2}
\end{array}\right]=\left[\begin{array}{cc}
1-e^{\varepsilon} \alpha^{2} & -e^{\varepsilon} \alpha \beta \\
-e^{\varepsilon} \alpha \beta & -e^{\varepsilon} \beta^{2}
\end{array}\right]=: A .
$$

Since $\operatorname{Tr} A=1-e^{\varepsilon}$, $\operatorname{det} A=-e^{\varepsilon} \beta^{2}$ and $D=(\operatorname{Tr} A)^{2}-4 \operatorname{det} A$, the greatest eigenvalue of $A$ is equal to $\left(1-e^{\varepsilon}+\sqrt{D}\right) / 2$. Therefore, inequality (4.8) holds. Consequently, we obtain

$$
t\left(\left|u_{i}\right\rangle\left\langle u_{i}\right|-e^{\varepsilon}\left|u_{j}\right\rangle\left\langle u_{j}\right|\right) \leq t_{\max } \frac{1-e^{\varepsilon}+\sqrt{D}}{2} I_{n}=\left(e^{\varepsilon}-1\right) I_{n}
$$

for all $0<t \leq t_{\text {max }}$. This implies $\rho_{i} \leq e^{\varepsilon} \rho_{j}$.
Recalling the definition of $M_{n}^{\mathrm{X}}\left(\varepsilon ; J_{\theta}\right)$, we have the monotonicity

$$
M_{2}^{\mathrm{X}}\left(\varepsilon ; J_{\theta}\right) \geq M_{3}^{\mathrm{X}}\left(\varepsilon ; J_{\theta}\right) \geq \cdots
$$

for $\mathrm{X}=\mathrm{C}, \mathrm{EC}, \mathrm{CQ}$. This monotonicity is used below.
Lemma 4.12. For all $\theta \in[0,1], \varepsilon>0$ and $n \geq 2, M_{n}^{\mathrm{CQ}}\left(\varepsilon ; J_{\theta}\right)=M_{2}^{\mathrm{C}}\left(\varepsilon ; J_{\theta}\right)$.

Proof. Let $\theta \in[0,1], c \in[0,1)$ and $\varepsilon, t>0$ be real numbers, and $n \geq 2$ be an integer. Take a CQ $\varepsilon$-DP tuple $\left(\rho_{i}\right)_{i=1}^{n}$ in Lemma 4.11, i.e.,

$$
\rho_{i}=\frac{1}{n+t}\left(I_{n}+t\left|u_{i}\right\rangle\left\langle u_{i}\right|\right), \quad u_{i} \in \mathbb{R}^{n} \quad(i=1, \ldots, n)
$$

and $\left\langle u_{i} \mid u_{j}\right\rangle=c$ for all $i \neq j$. Fix $i \neq j$ arbitrarily. The matrix $\left|u_{i}\right\rangle\left\langle u_{i}\right|-\left|u_{j}\right\rangle\left\langle u_{j}\right|$ can be expressed as a square matrix of order 2 in the same way as 4.9):

$$
\left[\begin{array}{ll}
1 & 0 \\
0 & 0
\end{array}\right]-\left[\begin{array}{cc}
\alpha^{2} & \alpha \beta \\
\alpha \beta & \beta^{2}
\end{array}\right]=\left[\begin{array}{cc}
\beta^{2} & -\alpha \beta \\
-\alpha \beta & -\beta^{2}
\end{array}\right]=-\beta\left[\begin{array}{cc}
-\beta & \alpha \\
\alpha & \beta
\end{array}\right],
$$

where $\alpha=c$ and $\beta=\sqrt{1-\alpha^{2}}$. Moreover, $\left(\left|u_{i}\right\rangle\left\langle u_{i}\right|-\left|u_{j}\right\rangle\left\langle u_{j}\right|\right)^{2}$ is expressed as

$$
\beta^{2}\left[\begin{array}{cc}
-\beta & \alpha \\
\alpha & \beta
\end{array}\right]^{2}=\beta^{2} I_{2}
$$

Denote by $\lambda_{1}$ and $\lambda_{2}$ two eigenvalues of the matrix

$$
A:=I_{2}+t\left((1-\theta)\left[\begin{array}{ll}
1 & 0 \\
0 & 0
\end{array}\right]+\theta\left[\begin{array}{cc}
\alpha^{2} & \alpha \beta \\
\alpha \beta & \beta^{2}
\end{array}\right]\right)=\left[\begin{array}{cc}
1+t\left(1-\theta \beta^{2}\right) & t \theta \alpha \beta \\
t \theta \alpha \beta & 1+t \theta \beta^{2}
\end{array}\right] .
$$

It follows that

$$
\begin{aligned}
& J_{\theta}\left(\rho_{i}, \rho_{j}\right)=\frac{t^{2}}{n+t} \operatorname{Tr}\left(\left|u_{i}\right\rangle\left\langle u_{i}\right|-\left|u_{j}\right\rangle\left\langle u_{j}\right|\right)^{2}\left(I_{n}+t\left((1-\theta)\left|u_{i}\right\rangle\left\langle u_{i}\right|+\theta\left|u_{j}\right\rangle\left\langle u_{j}\right|\right)\right)^{-1} \\
= & \frac{t^{2}}{n+t} \beta^{2} \operatorname{Tr} A^{-1}=\frac{t^{2}}{n+t} \beta^{2}\left(1 / \lambda_{1}+1 / \lambda_{2}\right)=\frac{t^{2}}{n+t} \beta^{2} \frac{\lambda_{1}+\lambda_{2}}{\lambda_{1} \lambda_{2}}=\frac{t^{2}}{n+t} \beta^{2} \frac{\operatorname{Tr} A}{\operatorname{det} A} .
\end{aligned}
$$

Since $\operatorname{Tr} A=2+t$ and

$$
\begin{aligned}
\operatorname{det} A & =\left(1+t\left(1-\theta \beta^{2}\right)\right)\left(1+t \theta \beta^{2}\right)-(t \theta \alpha \beta)^{2} \\
& =1+t+t^{2}\left(1-\theta \beta^{2}\right) \theta \beta^{2}-(t \theta \beta)^{2} \alpha^{2} \\
& =1+t+t^{2} \theta \beta^{2}-(t \theta \beta)^{2}=1+t+t^{2} \theta(1-\theta) \beta^{2} \\
& =1+t+t^{2} \theta(1-\theta)\left(1-c^{2}\right),
\end{aligned}
$$

we have

$$
\begin{equation*}
J_{\theta}\left(\rho_{i}, \rho_{j}\right)=\frac{t^{2}}{n+t} \beta^{2} \frac{\operatorname{Tr} A}{\operatorname{det} A}=\frac{t^{2}}{n+t} \cdot\left(1-c^{2}\right) \cdot \frac{2+t}{1+t+t^{2} \theta(1-\theta)\left(1-c^{2}\right)} \tag{4.10}
\end{equation*}
$$

Finally, putting $t=t_{\text {max }}$, we take the limit $c \rightarrow 1-0$. Set $s=e^{\varepsilon}-1>0$. Then $D=s^{2}+4\left(1-c^{2}\right) e^{\varepsilon}$ and

$$
t_{\max }=\frac{2 s}{\sqrt{D}-s}=\frac{2 s(\sqrt{D}+s)}{D-s^{2}}=\frac{s(\sqrt{D}+s)}{2\left(1-c^{2}\right) e^{\varepsilon}}
$$

Thus, the positive number $t=t_{\text {max }}$ diverges to $+\infty$ as $c \rightarrow 1-0$, and moreover,

$$
\begin{gathered}
\lim _{c \rightarrow 1-0} \frac{t}{n+t}=1, \quad \lim _{c \rightarrow 1-0}\left(1-c^{2}\right) t=\frac{s^{2}}{e^{\varepsilon}} \\
\lim _{c \rightarrow 1-0} \frac{2+t}{1+t+t^{2} \theta(1-\theta)\left(1-c^{2}\right)}=\lim _{c \rightarrow 1-0} \frac{t}{t+t^{2} \theta(1-\theta)\left(1-c^{2}\right)} \\
=\lim _{c \rightarrow 1-0} \frac{1}{1+t \theta(1-\theta)\left(1-c^{2}\right)}=\frac{1}{1+\theta(1-\theta) s^{2} / e^{\varepsilon}}=\frac{e^{\varepsilon}}{e^{\varepsilon}+\theta(1-\theta) s^{2}} .
\end{gathered}
$$

Since Theorem 1.5 and Lemma 4.9 yield that

$$
\begin{aligned}
& M_{2}^{\mathrm{CQ}}\left(\varepsilon ; J_{\theta}\right)=M_{2}^{\mathrm{C}}\left(\varepsilon ; J_{\theta}\right)=\frac{f_{\theta}\left(e^{\varepsilon}, 1\right)+f_{\theta}\left(1, e^{\varepsilon}\right)}{e^{\varepsilon}+1} \\
= & \frac{1}{e^{\varepsilon}+1} \cdot \frac{\left(e^{\varepsilon}-1\right)^{2}\left(e^{\varepsilon}+1\right)}{\left((1-\theta) e^{\varepsilon}+\theta\right)\left(\theta e^{\varepsilon}+1-\theta\right)}=\frac{s^{2}}{((1-\theta) s+1)(\theta s+1)},
\end{aligned}
$$

it turns out that

$$
\begin{aligned}
& M_{2}^{\mathrm{C}}\left(\varepsilon ; J_{\theta}\right)=M_{2}^{\mathrm{CQ}}\left(\varepsilon ; J_{\theta}\right) \geq M_{n}^{\mathrm{CQ}}\left(\varepsilon ; J_{\theta}\right) \geq \lim _{c \rightarrow 1-0} J\left(\rho_{1}, \rho_{2}\right) \\
= & \lim _{c \rightarrow 1-0} \frac{t^{2}}{n+t} \cdot\left(1-c^{2}\right) \cdot \frac{2+t}{1+t+t^{2} \theta(1-\theta)\left(1-c^{2}\right)} \\
= & \frac{s^{2}}{e^{\varepsilon}} \cdot \frac{e^{\varepsilon}}{e^{\varepsilon}+\theta(1-\theta) s^{2}}=\frac{s^{2}}{e^{\varepsilon}+\theta(1-\theta) s^{2}}=\frac{s^{2}}{((1-\theta) s+1)(\theta s+1)} \\
= & M_{2}^{\mathrm{C}}\left(\varepsilon ; J_{\theta}\right) .
\end{aligned}
$$

Proof of Theorem 4.5. The former assertion follows from Lemmas 4.9 and 4.12 immediately. Let us show the latter assertion. Let $\theta \in[0,1]$ and $\varepsilon>0$ be real numbers, and $n \geq 3$ be an integer. Since Lemmas 4.9 and 4.12 imply that

$$
\begin{aligned}
M_{n}^{\mathrm{C}}\left(\varepsilon ; J_{\theta}\right) & \leq M_{3}^{\mathrm{C}}\left(\varepsilon ; J_{\theta}\right)=\frac{f_{\theta}\left(e^{\varepsilon}, 1\right)+f_{\theta}\left(1, e^{\varepsilon}\right)}{2} \cdot \frac{2}{e^{\varepsilon}+2} \\
& =\frac{e^{\varepsilon}+1}{e^{\varepsilon}+2} M_{2}^{\mathrm{C}}\left(\varepsilon ; J_{\theta}\right)<M_{2}^{\mathrm{C}}\left(\varepsilon ; J_{\theta}\right)=M_{n}^{\mathrm{CQ}}\left(\varepsilon ; J_{\theta}\right),
\end{aligned}
$$

we obtain the latter assertion.

### 4.5 Concrete CQ $\varepsilon$-DP $n$-tuples that do not lie in $\mathrm{EC}_{n}(\varepsilon)$

We construct CQ $\varepsilon$-DP $n$-tuples that do not lie in $\mathrm{EC}_{n}(\varepsilon)$. In this section, we use the following lemmas instead of Lemmas 4.10 and 4.11 .

Lemma 4.13. Let $d \geq 2$ be an integer and $c \in[1 / d, 1]$ be a real number. There exists a $(d+1)$-tuple $\left(u_{i}\right)_{i=1}^{d+1}$ of unit vectors in $\mathbb{C}^{d}$ such that $\left|\left\langle u_{i} \mid u_{j}\right\rangle\right|=c$ for all $i \neq j$.

Proof. For $z \in \mathbb{C}$, define the Hermitian matrix $A(z)=\left(\alpha_{i, j}\right)$ of order $d+1$ as $\alpha_{i, j}=1$ if $i=j$ and $\alpha_{i, j}=z$ if $i<j$. Denote by eig $A(z)$ the set of all eigenvalues of $A(z)$. Then $\operatorname{eig} A(c)=\{1+d c, 1-c\}$, eig $A(-c)=\{1-d c, 1+c\}$, min eig $A(c)=1-c \geq 0$ and min eig $A(-c)=1-d c \leq 0$. Since the minimum eigenvalue of $A(z)$ can be expressed as

$$
\min \operatorname{eig} A(z)=\min _{\|u\|=1}\langle u| A(z)|u\rangle
$$

it follows that

$$
\left|\min \operatorname{eig} A(z)-\min \operatorname{eig} A\left(z^{\prime}\right)\right| \leq\left\|A(z)-A\left(z^{\prime}\right)\right\|
$$

for all $z, z^{\prime} \in \mathbb{C}$, which shows that min eig $A(z)$ is continuous in $z$. Thus, the intermediate value theorem implies that min eig $A\left(z_{0}\right)=0$ for some $z_{0} \in \mathbb{C}$ of magnitude c. Therefore, there exists a $d \times(d+1)$ complex matrix $B$ such that $A\left(z_{0}\right)=B^{*} B$. The column vectors $u_{1}, \ldots, u_{d+1} \in \mathbb{C}^{d}$ of $B$ satisfy that $\left\langle u_{i} \mid u_{j}\right\rangle=1$ if $i=j$ and $\left\langle u_{i} \mid u_{j}\right\rangle=z_{0}$ if $i<j$.

Lemma 4.14. Let $c \in[1 / d, 1)$ and $\varepsilon, t>0$ be real numbers, $\left(u_{i}\right)_{i=1}^{d+1}$ be the $(d+1)$ tuple in Lemma 4.13, and $\rho_{1}, \ldots, \rho_{d+1}$ be the density matrices defined as

$$
\rho_{i}=\frac{1}{d+t}\left(I_{d}+t\left|u_{i}\right\rangle\left\langle u_{i}\right|\right) \quad(i=1, \ldots, d+1) .
$$

If $D=\left(e^{\varepsilon}-1\right)^{2}+4\left(1-c^{2}\right) e^{\varepsilon}$ and

$$
0<t \leq t_{\max }:=\frac{2\left(e^{\varepsilon}-1\right)}{\sqrt{D}+1-e^{\varepsilon}}
$$

then $\left(\rho_{i}\right)_{i=1}^{d+1}$ is $C Q \varepsilon-D P$.
Proof. See the proof of Lemma 4.11.

Theorem 4.15. Let $\left(\rho_{i}\right)_{i=1}^{3}$ be a $C Q \varepsilon$-DP 3-tuple in Lemma 4.14 with $d=2$ and $t=t_{\mathrm{max}}$. Then $\left(\rho_{i}\right)_{i=1}^{3}$ does not lie in $\mathrm{EC}_{3}(\varepsilon)$.
Proof. Set $s=e^{\varepsilon}-1>0$. Then $D=s^{2}+4\left(1-c^{2}\right)(s+1)$ and $t=t_{\max }=2 s /(\sqrt{D}-s)$. Let $i \neq j$. We show that $M_{3}\left(\varepsilon ; J_{1 / 2}\right)<J_{1 / 2}\left(\rho_{i}, \rho_{j}\right)$. First, Lemma 4.9 and the equality

$$
f_{1 / 2}\left(e^{\varepsilon}, 1\right)+f_{1 / 2}\left(1, e^{\varepsilon}\right)=\frac{4\left(e^{\varepsilon}-1\right)^{2}}{e^{\varepsilon}+1}
$$

imply that

$$
M_{3}\left(\varepsilon ; J_{1 / 2}\right)=\frac{4\left(e^{\varepsilon}-1\right)^{2}}{2\left(e^{\varepsilon}+1\right)} \cdot \frac{2}{e^{\varepsilon}+2}=\frac{4 s^{2}}{(s+2)(s+3)}
$$

Also, it follows form the same calculation as 4.10) that

$$
J_{1 / 2}\left(\rho_{i}, \rho_{j}\right)=\frac{t^{2}}{2+t} \cdot\left(1-c^{2}\right) \cdot \frac{2+t}{1+t+(t / 2)^{2}\left(1-c^{2}\right)}=\frac{\left(1-c^{2}\right) t^{2}}{1+t+(t / 2)^{2}\left(1-c^{2}\right)}
$$

where we must replace $n$ in 4.10 with the dimension $d=2$. Thus,

$$
\begin{aligned}
& M_{3}\left(\varepsilon ; J_{1 / 2}\right)<J_{1 / 2}\left(\rho_{i}, \rho_{j}\right) \Longleftrightarrow \frac{4 s^{2}}{(s+2)(s+3)}<\frac{\left(1-c^{2}\right) t^{2}}{1+t+(t / 2)^{2}\left(1-c^{2}\right)} \\
& \Longleftrightarrow \frac{(s+2)(s+3)}{s^{2}}>\frac{4}{1-c^{2}} \cdot \frac{1+t+(t / 2)^{2}\left(1-c^{2}\right)}{t^{2}} \\
& \Longleftrightarrow 1+\frac{5 s+6}{s^{2}}>\frac{4}{1-c^{2}}\left(\frac{1+t}{t^{2}}+\frac{1-c^{2}}{4}\right) \Longleftrightarrow \frac{5 s+6}{s^{2}}>\frac{4}{1-c^{2}} \cdot \frac{1+t}{t^{2}}
\end{aligned}
$$

Recalling that $D=s^{2}+4\left(1-c^{2}\right)(s+1)$ and $t=t_{\max }=2 s /(\sqrt{D}-s)$, we have

$$
\begin{aligned}
\frac{t^{2}}{1+t} & =t-1+\frac{1}{1+t}=\frac{2 s}{\sqrt{D}-s}-1+\frac{\sqrt{D}-s}{\sqrt{D}+s}=\frac{2 s}{\sqrt{D}-s}+\frac{-2 s}{\sqrt{D}+s} \\
& =\frac{4 s^{2}}{D-s^{2}}=\frac{s^{2}}{\left(1-c^{2}\right)(s+1)}
\end{aligned}
$$

Therefore,

$$
M_{3}\left(\varepsilon ; J_{1 / 2}\right)<J_{1 / 2}\left(\rho_{i}, \rho_{j}\right) \Longleftrightarrow \frac{5 s+6}{s^{2}}>\frac{4(s+1)}{s^{2}}
$$

Since the right inequality always holds, so does the left inequality. From Corollary 4.6, it follows that $\left(\rho_{i}\right)_{i=1}^{3}$ does not lie in $\mathrm{EC}_{3}(\varepsilon)$.
Corollary 4.16. Let $n \geq 3$ be an integer, and $\left(\rho_{i}\right)_{i=1}^{3}$ be a $C Q \varepsilon$-DP 3-tuple in Lemma 4.14 with $d=2$ and $t=t_{\max }$. Then every $\left(\sigma_{i}\right)_{i=1}^{n} \in \mathrm{CQ}_{n}(\varepsilon)$ with $\sigma_{i}=\rho_{i}$, $i=1,2,3$, does not lie in $\mathrm{EC}_{n}(\varepsilon)$.
Proof. This corollary follows from Theorem 4.15 and Definition 1.4 .

## Chapter 5

## Perfect discrimination of two states in approximate quantum theory

In this chapter, we construct two kinds of one-parameter family of proper cones. They define locally quantum systems that are sufficiently close to quantum theory. We show sufficient conditions for two special states to be perfectly distinguishable in the locally quantum systems.

### 5.1 One-parameter families of proper cones

We begin with a few basic terms. Let $d_{1}, d_{2} \geq 2$ be integers. For $v \in \mathbb{C}^{d_{1}} \otimes \mathbb{C}^{d_{2}}$, there exist unique coefficients $\lambda_{1} \geq \cdots \geq \lambda_{d} \geq 0, d=\min \left\{d_{1}, d_{2}\right\}$, such that $v$ is expressed as $v=\sum_{i=1}^{d} \lambda_{i} v_{i}^{[1]} \otimes v_{i}^{[2]}$ with two orthonormal systems $\left(v_{i}^{[1]}\right)_{i=1}^{d}$ and $\left(v_{i}^{[2]}\right)_{i=1}^{d}$. The coefficients $\lambda_{1} \geq \cdots \geq \lambda_{d}$ are called the Schmidt coefficients of $v$. If $v \neq 0$, they are closely related to the entanglement of the quantum state $\|v\|^{-2}|v\rangle\langle v|$ (see a textbook in quantum information theory, e.g., [26, 27]). Also, let $\mathrm{id}_{d}$ be the identity mapping on $\operatorname{Herm}(d)$, and $\tau_{d}$ be the transposition on $\operatorname{Herm}(d)$. We denote by $\Gamma$ the partial transpose on the second system, i.e., $\Gamma=\operatorname{id}_{d_{1}} \otimes \tau_{d_{2}}$. If a density matrix $\rho$ lies in $\operatorname{Sep}\left(d_{1}, d_{2}\right)$, then $\rho$ has positive partial transpose, i.e., $\Gamma(\rho) \geq 0$. This fact is called the positive partial transpose criterion. The converse is true for $\left(d_{1}, d_{2}\right)=(2,2),(2,3)$ [31], but is false for $\left(d_{1}, d_{2}\right)=(3,3),(2,4)$ 32].

Now, we consider $\left(d_{1}, d_{2}\right)$-level locally quantum systems. Let us construct two kinds of one-parameter family of proper cones.

Definition 5.1 (One-parameter family of proper cones, I). For $s \geq 0$, we define the proper cone $\mathcal{K}_{s}^{\text {neg }}$ as

$$
\mathcal{K}_{s}^{\mathrm{neg}}=\left\{X \in \operatorname{Sep}\left(d_{1}, d_{2}\right)^{*}: \operatorname{neg}(X) \leq s \operatorname{Tr} X\right\}
$$

where for $X \in \operatorname{Herm}\left(d_{1} d_{2}\right)$ the value $\operatorname{neg}(X)$ is defined as

$$
\operatorname{neg}(X)=\max _{\lambda \text { eigenvalue of } X}\{-\lambda, 0\} .
$$

Note that the inequality

$$
\operatorname{neg}(X+Y) \leq \operatorname{neg}(X)+\operatorname{neg}(Y)
$$

holds for all $X, Y \in \operatorname{Herm}\left(d_{1} d_{2}\right)$.
Definition 5.2 (One-parameter family of proper cones, II). For a vector $v \in \mathbb{C}^{d_{1}} \otimes$ $\mathbb{C}^{d_{2}}$, let $\operatorname{sc}(v)$ be the value

$$
\operatorname{sc}(v)= \begin{cases}\lambda_{1} \lambda_{2} & v \neq 0 \\ 0 & v=0\end{cases}
$$

where $\lambda_{1} \geq \lambda_{2} \geq \cdots \geq \lambda_{d}$, $d=\min \left\{d_{1}, d_{2}\right\}$, denote the Schmidt coefficients of $v /\|v\|$. Then, for $s \geq 0$, we define the proper cones $\mathcal{K}_{s}^{(0)}$ and $\mathcal{K}_{s}^{\mathrm{sc}}$ as

$$
\begin{aligned}
\mathcal{K}_{s}^{(0)} & =\operatorname{conv}\left\{|v\rangle\langle v|: v \in \mathbb{C}^{d_{1}} \otimes \mathbb{C}^{d_{2}}, \operatorname{sc}(v) \leq s\right\} \\
\mathcal{K}_{s}^{\mathrm{sc}} & =\operatorname{PSD}\left(d_{1} d_{2}\right)+\Gamma\left(\mathcal{K}_{s}^{(0)}\right)
\end{aligned}
$$

where $\mathcal{K}_{1}+\mathcal{K}_{2}$ denotes the Minkowski sum of convex cones $\mathcal{K}_{1}$ and $\mathcal{K}_{2}$.
The value $\operatorname{sc}(v)$ is closely related to negative eigenvalues of $\Gamma(|v\rangle\langle v|)$ : the equality

$$
\begin{equation*}
\|v\|^{2} \operatorname{sc}(v)=\operatorname{neg}(\Gamma(|v\rangle\langle v|)) \tag{5.1}
\end{equation*}
$$

holds for every $v \in \mathbb{C}^{d_{1}} \otimes \mathbb{C}^{d_{2}}$. This is because if $v \neq 0$, then the set of all eigenvalues of $\Gamma(|v\rangle\langle v|) /\|v\|^{2}$ is

$$
\left\{ \pm \lambda_{i} \lambda_{j}, \lambda_{k}^{2}: 1 \leq i<j \leq d, 1 \leq k \leq d\right\}
$$

where $\lambda_{1} \geq \cdots \geq \lambda_{d}$ denote the Schmidt coefficients of $v /\|v\|$. An element in $\mathcal{K}_{s}^{\text {sc }}$ may have negative eigenvalues, but they are restricted by the definition of $\mathcal{K}_{s}^{\text {sc }}$.

Once the parameter $s$ increases, the proper cones $\mathcal{K}_{s}^{\text {neg }}, \mathcal{K}_{s}^{(0)}$, and $\mathcal{K}_{s}^{\text {sc }}$ become larger. Thus, the following inclusion relations hold:

$$
\begin{aligned}
& \operatorname{PSD}\left(d_{1} d_{2}\right)=\mathcal{K}_{0}^{\mathrm{neg}} \subset \mathcal{K}_{s}^{\mathrm{neg}} \subset \operatorname{Sep}\left(d_{1}, d_{2}\right)^{*} \\
& \operatorname{Sep}\left(d_{1}, d_{2}\right)=\mathcal{K}_{0}^{(0)} \subset \mathcal{K}_{s}^{(0)} \subset \operatorname{PSD}\left(d_{1} d_{2}\right) \\
& \operatorname{PSD}\left(d_{1} d_{2}\right)=\mathcal{K}_{0}^{\mathrm{sc}} \subset \mathcal{K}_{s}^{\mathrm{sc}} \subset \operatorname{Sep}\left(d_{1}, d_{2}\right)^{*}
\end{aligned}
$$

Since $\mathcal{K}_{s}^{(0)}$ satisfies local unitary invariance, i.e., $\left(U^{[1]} \otimes U^{[2]}\right) \mathcal{K}_{s}^{(0)}\left(U^{[1]} \otimes U^{[2]}\right)^{*}=\mathcal{K}_{s}^{(0)}$ for all unitary matrices $U^{[1]}$ and $U^{[2]}$, no proper cones $\mathcal{K}_{s}^{\text {sc }}$ depend on the orthonormal basis of $\mathbb{C}^{d_{2}}$ that defines the partial transpose $\Gamma$.

### 5.2 Proofs

Actually, it suffices to prove Theorems 1.8 and 1.9 only for $\left(d_{1}, d_{2}\right)=(2,2)$ due to the following reason. Let $\rho_{1}^{[1]}, \rho_{1}^{[2]}, \rho_{2}^{[1]}, \rho_{2}^{[2]}$ be the density matrices in Theorems 1.8 and 1.9. Then there exist orthonormal bases $\mathcal{B}^{[i]}$ of $\mathbb{C}^{d_{i}}, i=1,2$, such that the matrix representations of $\rho_{1}^{[1]}, \rho_{1}^{[2]}, \rho_{2}^{[1]}, \rho_{2}^{[2]}$ with respect to $\mathcal{B}^{[1]}$ and $\mathcal{B}^{[2]}$ are given below:

$$
\begin{aligned}
& \rho_{1}^{[1]}=\left[\begin{array}{ccccc}
1 & 0 & & & \\
0 & 0 & & & \\
& & 0 & & \\
& & & \ddots & \\
& & & & 0
\end{array}\right], \quad \rho_{1}^{[2]}=\left[\begin{array}{ccccc}
1 & 0 & & & \\
0 & 0 & & & \\
& & 0 & & \\
& & \ddots & \\
& & & & 0
\end{array}\right], \\
& \rho_{2}^{[1]}=\left[\begin{array}{ccccc}
1-\alpha_{1} & \beta_{1} & & & \\
\beta_{1} & \alpha_{1} & & & \\
& & 0 & & \\
& & & \ddots & \\
& & & 0
\end{array}\right], \quad \rho_{2}^{[2]}=\left[\begin{array}{cccc}
1-\alpha_{2} & \beta_{2} & & \\
\beta_{2} & \alpha_{2} & & \\
& & 0 & \\
& & & \ddots \\
& & & 0
\end{array}\right],
\end{aligned}
$$

where $\alpha_{1}, \alpha_{2} \in[0,1]$ and $\beta_{i}=\sqrt{\alpha_{i}\left(1-\alpha_{i}\right)}$. That is, each representation matrix above is equal to the direct sum of a $2 \times 2$ matrix and the zero matrix. Since the proper cones $\mathcal{K}_{s}^{\text {sc }}$ are independent of an orthonormal basis of $\mathbb{C}^{d_{2}}$ that defines the partial transpose $\Gamma$, the general case is reduced to the case $\left(d_{1}, d_{2}\right)=(2,2)$. Therefore, it suffices to prove the following lemmas.

Lemma 5.3 (Perfect discrimination with $\left.\mathcal{M}\left(\mathcal{K}_{s}^{\text {neg }}\right)\right)$. Let $\rho_{1}$ and $\rho_{2}$ be the rank-one density matrices in $\operatorname{Sep}(2,2)$ given as

$$
\rho_{1}=\left[\begin{array}{ll}
1 & 0  \tag{5.2}\\
0 & 0
\end{array}\right] \otimes\left[\begin{array}{ll}
1 & 0 \\
0 & 0
\end{array}\right], \quad \rho_{2}=\left[\begin{array}{cc}
1-\alpha_{1} & \beta_{1} \\
\beta_{1} & \alpha_{1}
\end{array}\right] \otimes\left[\begin{array}{cc}
1-\alpha_{2} & \beta_{2} \\
\beta_{2} & \alpha_{2}
\end{array}\right],
$$

where $\alpha_{1}, \alpha_{2} \in[0,1]$ and $\beta_{i}=\sqrt{\alpha_{i}\left(1-\alpha_{i}\right)}$. If $s \in[0,1 / 4]$ and

$$
\begin{equation*}
\left(1-\alpha_{1}\right)\left(1-\alpha_{2}\right) \leq 16 s^{2} \alpha_{1} \alpha_{2}, \tag{5.3}
\end{equation*}
$$

then the two states $\rho_{1}$ and $\rho_{2}$ are perfectly distinguishable by some measurement $\left\{T_{i}+\Gamma\left(T_{i}\right)\right\}_{i=1,2} \in \mathcal{M}\left(\mathcal{K}_{s}^{\mathrm{neg}}\right)$. The measurement $\left\{T_{i}+\Gamma\left(T_{i}\right)\right\}_{i=1,2} \in \mathcal{M}\left(\mathcal{K}_{s}^{\mathrm{neg}}\right)$ is given below except for the trivial cases $\alpha_{1}=1$ and $\alpha_{2}=1$ : If $\gamma:=\alpha_{1}+\alpha_{2}>1$, then

$$
\begin{gather*}
2 \gamma T_{1}=\gamma\left|v_{1}\right\rangle\left\langle v_{1}\right|+(\gamma-1)\left|v_{2}\right\rangle\left\langle v_{2}\right|+(\gamma-1)\left|v_{3}\right\rangle\left\langle v_{3}\right|,  \tag{5.4}\\
v_{1}=\left[\begin{array}{l}
1 \\
0
\end{array}\right] \otimes\left[\begin{array}{l}
1 \\
0
\end{array}\right]-\frac{\beta_{1} \beta_{2}}{\alpha_{1} \alpha_{2}}\left[\begin{array}{l}
0 \\
1
\end{array}\right] \otimes\left[\begin{array}{l}
0 \\
1
\end{array}\right],  \tag{5.5}\\
v_{2}=\left[\begin{array}{c}
1 \\
-\beta_{1} / \alpha_{1}
\end{array}\right] \otimes\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad v_{3}=\left[\begin{array}{l}
0 \\
1
\end{array}\right] \otimes\left[\begin{array}{c}
1 \\
-\beta_{2} / \alpha_{2}
\end{array}\right],  \tag{5.6}\\
T_{2}=\left(U^{[1]} \otimes U^{[2]}\right) T_{1}\left(U^{[1]} \otimes U^{[2]}\right)^{*},  \tag{5.7}\\
U^{[1]}=\frac{1}{\sqrt{\alpha_{1}}}\left[\begin{array}{cc}
\beta_{1} & \alpha_{1} \\
\alpha_{1} & -\beta_{1}
\end{array}\right], \quad U^{[2]}=\frac{1}{\sqrt{\alpha_{2}}}\left[\begin{array}{cc}
\beta_{2} & \alpha_{2} \\
\alpha_{2} & -\beta_{2}
\end{array}\right] ; \tag{5.8}
\end{gather*}
$$

if $\gamma=1$, then

$$
T_{1}=\frac{1}{2}\left[\begin{array}{cccc}
1 & 0 & 0 & -1 \\
0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 \\
-1 & 0 & 0 & 1
\end{array}\right], \quad T_{2}=\frac{1}{2}\left[\begin{array}{cccc}
0 & 0 & 0 & 0 \\
0 & 1 & 1 & 0 \\
0 & 1 & 1 & 0 \\
0 & 0 & 0 & 0
\end{array}\right]
$$

Lemma 5.4 (Perfect discrimination with $\mathcal{M}\left(\mathcal{K}_{s}^{\mathrm{sc}}\right)$ ). Let $\rho_{1}$ and $\rho_{2}$ be the rank-one density matrices in $\operatorname{Sep}(2,2)$ given as (5.2). If $s=\sqrt{t} /(1+t), t \in[0,1]$, and

$$
\begin{equation*}
\left(1-\alpha_{1}\right)\left(1-\alpha_{2}\right) \leq t \alpha_{1} \alpha_{2}, \tag{5.9}
\end{equation*}
$$

then the two states $\rho_{1}$ and $\rho_{2}$ are perfectly distinguishable by the measurement $\left\{T_{i}+\right.$ $\left.\Gamma\left(T_{i}\right)\right\}_{i=1,2} \in \mathcal{M}\left(\mathcal{K}_{s}^{\mathrm{neg}}\right)$ given in Lemma 5.3 .

We prove Lemmas 5.4 and 5.3 in this order.

Proof of Lemma 5.4. Assume that $s=\sqrt{t} /(1+t), t \in[0,1]$, and (5.9). All we need is to show that
(i) $T_{1}+T_{2}+\Gamma\left(T_{1}+T_{2}\right)=I_{4}$,
(ii) $T_{i} \in \mathcal{K}_{s}^{(0)}$ for $i=1,2$,
(iii) $\operatorname{Tr} \rho_{1} T_{2}=\operatorname{Tr} \rho_{2} T_{1}=0$.

Indeed, if (i) and (ii) hold, then $\left\{T_{i}+\Gamma\left(T_{i}\right)\right\}_{i=1,2} \in \mathcal{M}\left(\mathcal{K}_{s}^{\text {sc }}\right)$. Also, if (i) and (iii) hold, then the equations $\Gamma\left(\rho_{i}\right)=\rho_{i}, i=1,2$, imply that $\operatorname{Tr} \rho_{i}\left(T_{j}+\Gamma\left(T_{j}\right)\right)=2 \operatorname{Tr} \rho_{i} T_{j}=\delta_{i j}$ for all $i, j \in\{1,2\}$. Therefore, if (i)-(iii) hold, then Lemma 5.4 follows. Also, note that $\left(1-\alpha_{2}\right)\left(1-\alpha_{1}\right) \leq t \alpha_{1} \alpha_{2} \leq \alpha_{1} \alpha_{2}$ due to $t \in[0,1]$ and (5.9). Thus, $\gamma=\alpha_{1}+\alpha_{2} \geq 1$. If $\alpha_{1} \alpha_{2}=0$, then $\alpha_{1}=1$ or $\alpha_{2}=1$, which is a trivial case. Therefore, without loss of generality, we may assume $\alpha_{1} \alpha_{2}>0$.

Proof of (i). First, assume $\gamma=1$. Then

$$
T_{1}+T_{2}+\Gamma\left(T_{1}+T_{2}\right)=\frac{1}{2}\left[\begin{array}{cccc}
1 & 0 & 0 & -1 \\
0 & 1 & 1 & 0 \\
0 & 1 & 1 & 0 \\
-1 & 0 & 0 & 1
\end{array}\right]+\frac{1}{2}\left[\begin{array}{cccc}
1 & 0 & 0 & 1 \\
0 & 1 & -1 & 0 \\
0 & -1 & 1 & 0 \\
1 & 0 & 0 & 1
\end{array}\right]=I_{4}
$$

Next, assume $\gamma>1$. Put $w_{i}=\left(U^{[1]} \otimes U^{[2]}\right) v_{i}$ for $i=1,2,3$. Then $w_{i}, i=1,2,3$, can be calculated as follows:

$$
\begin{gathered}
w_{1}=\frac{1}{\sqrt{\alpha_{1} \alpha_{2}}}\left(\left[\begin{array}{l}
\beta_{1} \\
\alpha_{1}
\end{array}\right] \otimes\left[\begin{array}{l}
\beta_{2} \\
\alpha_{2}
\end{array}\right]-\frac{\beta_{1} \beta_{2}}{\alpha_{1} \alpha_{2}}\left[\begin{array}{c}
\alpha_{1} \\
-\beta_{1}
\end{array}\right] \otimes\left[\begin{array}{c}
\alpha_{2} \\
-\beta_{2}
\end{array}\right]\right) \\
=\frac{1}{\sqrt{\alpha_{1} \alpha_{2}}}\left(\left[\begin{array}{c}
\beta_{1} \beta_{2} \\
\beta_{1} \alpha_{2} \\
\alpha_{1} \beta_{2} \\
\alpha_{1} \alpha_{2}
\end{array}\right]-\frac{\beta_{1} \beta_{2}}{\alpha_{1} \alpha_{2}}\left[\begin{array}{c}
\alpha_{1} \alpha_{2} \\
-\alpha_{1} \beta_{2} \\
-\beta_{1} \alpha_{2} \\
\beta_{1} \beta_{2}
\end{array}\right]\right)=\frac{1}{\sqrt{\alpha_{1} \alpha_{2}}}\left[\begin{array}{c}
0 \\
\beta_{1} \\
\beta_{2} \\
\gamma-1
\end{array}\right], \\
w_{2}=\frac{1}{\sqrt{\alpha_{1} \alpha_{2}}}\left[\begin{array}{l}
0 \\
1
\end{array}\right] \otimes\left[\begin{array}{c}
\alpha_{2} \\
-\beta_{2}
\end{array}\right]=\sqrt{\frac{\alpha_{2}}{\alpha_{1}}} v_{3}, \quad w_{3}=\frac{1}{\sqrt{\alpha_{1} \alpha_{2}}}\left[\begin{array}{c}
\alpha_{1} \\
-\beta_{1}
\end{array}\right] \otimes\left[\begin{array}{l}
0 \\
1
\end{array}\right]=\sqrt{\frac{\alpha_{1}}{\alpha_{2}}} v_{2} .
\end{gathered}
$$

Thus, putting $\xi=\beta_{1} \beta_{2} / \alpha_{1} \alpha_{2}$, we have

$$
\begin{aligned}
& T_{1}+T_{2}=\frac{1}{2}\left(\left|v_{1}\right\rangle\left\langle v_{1}\right|+\left|w_{1}\right\rangle\left\langle w_{1}\right|\right)+\frac{\gamma-1}{2 \gamma}\left(\left|v_{2}\right\rangle\left\langle v_{2}\right|+\left|v_{3}\right\rangle\left\langle v_{3}\right|+\left|w_{2}\right\rangle\left\langle w_{2}\right|+\left|w_{3}\right\rangle\left\langle w_{3}\right|\right) \\
= & \frac{1}{2}\left(\left|v_{1}\right\rangle\left\langle v_{1}\right|+\left|w_{1}\right\rangle\left\langle w_{1}\right|\right)+\frac{\gamma-1}{2 \gamma}\left(\left|v_{2}\right\rangle\left\langle v_{2}\right|+\left|v_{3}\right\rangle\left\langle v_{3}\right|+\frac{\alpha_{2}}{\alpha_{1}}\left|v_{3}\right\rangle\left\langle v_{3}\right|+\frac{\alpha_{1}}{\alpha_{2}}\left|v_{2}\right\rangle\left\langle v_{2}\right|\right) \\
= & \frac{1}{2}\left(\left|v_{1}\right\rangle\left\langle v_{1}\right|+\left|w_{1}\right\rangle\left\langle w_{1}\right|\right)+\frac{\gamma-1}{2}\left(\frac{1}{\alpha_{2}}\left|v_{2}\right\rangle\left\langle v_{2}\right|+\frac{1}{\alpha_{1}}\left|v_{3}\right\rangle\left\langle v_{3}\right|\right)
\end{aligned}
$$

and

$$
\left.\begin{array}{rl}
T_{1}+T_{2}= & \frac{1}{2}\left[\begin{array}{cccc}
1 & 0 & 0 & -\xi \\
0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 \\
-\xi & 0 & 0 & \xi^{2}
\end{array}\right]+\frac{1}{2 \alpha_{1} \alpha_{2}}\left[\begin{array}{ccc}
0 & 0 & 0 \\
0 & \beta_{1}^{2} & \beta_{1} \beta_{2}
\end{array} \begin{array}{c}
(\gamma-1) \beta_{1} \\
0 \\
\beta_{1} \beta_{2} \\
0 \\
(\gamma-1) \beta_{1} \\
(\gamma-1) \beta_{2}^{2}
\end{array}\right. \\
& (\gamma-1) \beta_{2} \\
& +\frac{\gamma-1}{2 \alpha_{1} \alpha_{2}}\left(\left[\begin{array}{cc}
\alpha_{1} & -\beta_{1} \\
-\beta_{1} & 1-\alpha_{1}
\end{array}\right] \otimes\left[\begin{array}{ll}
0 & 0 \\
0 & 1
\end{array}\right]+\left[\begin{array}{ll}
0 & 0 \\
0 & 1
\end{array}\right] \otimes\left[\begin{array}{cc}
\alpha_{2} & -\beta_{2} \\
-\beta_{2} & 1-\alpha_{2}
\end{array}\right]\right) \\
= & \frac{1}{2}\left[\begin{array}{cccc}
1 & 0 & 0 & -\xi \\
0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 \\
-\xi & 0 & 0 & \xi^{2}
\end{array}\right]+\frac{1}{2 \alpha_{1} \alpha_{2}}\left[\begin{array}{ccc}
0 & 0 & 0
\end{array}\right] 0 \\
0 & \beta_{1}^{2} \\
0 & \beta_{1} \beta_{2} \\
\beta_{1} \beta_{2} & (\gamma-1) \beta_{1} \\
0 & (\gamma-1) \beta_{1}^{2} \\
(\gamma-1) \beta_{2} & (\gamma-1) \beta_{2} \\
0-1)^{2}
\end{array}\right] .
$$

When $t_{i j}$ denotes the $(i, j)$-th entry of $T_{1}+T_{2}$, it follows that $t_{11}=1 / 2, t_{12}=t_{21}=$ $t_{13}=t_{31}=0, t_{14}=t_{41}=-\xi / 2, t_{23}=t_{32}=\xi / 2$,

$$
\begin{aligned}
t_{24}=t_{42}= & \frac{(\gamma-1) \beta_{1}}{2 \alpha_{1} \alpha_{2}}-\frac{(\gamma-1) \beta_{1}}{2 \alpha_{1} \alpha_{2}}=0, \quad t_{34}=t_{43}=\frac{(\gamma-1) \beta_{2}}{2 \alpha_{1} \alpha_{2}}-\frac{(\gamma-1) \beta_{2}}{2 \alpha_{1} \alpha_{2}}=0, \\
t_{22} & =\frac{\beta_{1}^{2}}{2 \alpha_{1} \alpha_{2}}+\frac{(\gamma-1) \alpha_{1}}{2 \alpha_{1} \alpha_{2}}=\frac{1-\alpha_{1}}{2 \alpha_{2}}+\frac{\gamma-1}{2 \alpha_{2}}=1 / 2, \\
t_{33} & =\frac{\beta_{2}^{2}}{2 \alpha_{1} \alpha_{2}}+\frac{(\gamma-1) \alpha_{2}}{2 \alpha_{1} \alpha_{2}}=\frac{1-\alpha_{2}}{2 \alpha_{1}}+\frac{\gamma-1}{2 \alpha_{1}}=1 / 2, \\
t_{44} & =\frac{\xi^{2}}{2}+\frac{(\gamma-1)^{2}}{2 \alpha_{1} \alpha_{2}}+\frac{(\gamma-1)\left(1-\alpha_{1}\right)}{2 \alpha_{1} \alpha_{2}}+\frac{(\gamma-1)\left(1-\alpha_{2}\right)}{2 \alpha_{1} \alpha_{2}} \\
& =\frac{\left(1-\alpha_{1}\right)\left(1-\alpha_{2}\right)}{2 \alpha_{1} \alpha_{2}}+\frac{(\gamma-1)^{2}}{2 \alpha_{1} \alpha_{2}}+\frac{(\gamma-1)(2-\gamma)}{2 \alpha_{1} \alpha_{2}} \\
& =\frac{1-\gamma+\alpha_{1} \alpha_{2}}{2 \alpha_{1} \alpha_{2}}+\frac{\gamma-1}{2 \alpha_{1} \alpha_{2}}=1 / 2 .
\end{aligned}
$$

Therefore,

$$
T_{1}+T_{2}+\Gamma\left(T_{1}+T_{2}\right)=\frac{1}{2}\left[\begin{array}{cccc}
1 & 0 & 0 & -\xi \\
0 & 1 & \xi & 0 \\
0 & \xi & 1 & 0 \\
-\xi & 0 & 0 & 1
\end{array}\right]+\frac{1}{2}\left[\begin{array}{cccc}
1 & 0 & 0 & \xi \\
0 & 1 & -\xi & 0 \\
0 & -\xi & 1 & 0 \\
\xi & 0 & 0 & 1
\end{array}\right]=I_{4}
$$

Proof of (ii). First, assume $\gamma=1$. Then

$$
\alpha_{1} \alpha_{2} \stackrel{\gamma=1}{=}\left(1-\alpha_{2}\right)\left(1-\alpha_{1}\right) \stackrel{\sqrt{5.9}}{\leq} t \alpha_{1} \alpha_{2} \stackrel{t \in[0,1]}{\leq} \alpha_{1} \alpha_{2}
$$

which shows that $t=1$ and $s=1 / 2$. Since it is easily checked that $T_{i} \in \mathcal{K}_{1 / 2}^{(0)}$ for $i=1,2$, we obtain (ii). Next, assume $\gamma>1$. Since the function $\sqrt{t^{\prime}} /\left(1+t^{\prime}\right)$, $t^{\prime} \in[0,1]$, is increasing, from (5.5) and (5.6), it follows that $\operatorname{sc}\left(v_{2}\right)=\operatorname{sc}\left(v_{3}\right)=0$ and

$$
\operatorname{sc}\left(v_{1}\right)=\frac{\beta_{1} \beta_{2} / \alpha_{1} \alpha_{2}}{1+\left(\beta_{1} \beta_{2} / \alpha_{1} \alpha_{2}\right)^{2}} \stackrel{\sqrt[5.9]{\leq}}{1+t} \frac{\sqrt{t}}{1+t}=s
$$

Thus, $T_{1} \in \mathcal{K}_{s}^{(0)}$. Thanks to (5.7), we also have $T_{2} \in \mathcal{K}_{s}^{(0)}$. Therefore, (ii) holds.
Proof of (iii). First, assume $\gamma=1$. Then it is easily checked that $\operatorname{Tr} \rho_{1} T_{2}=0$ and

$$
\operatorname{Tr} \rho_{2} T_{1}=\left(1-\alpha_{1}\right)\left(1-\alpha_{2}\right)+\alpha_{1} \alpha_{2}-2 \beta_{1} \beta_{2} \stackrel{\gamma=1}{=} \alpha_{1} \alpha_{2}+\alpha_{1} \alpha_{2}-2 \alpha_{1} \alpha_{2}=0
$$

which are just (iii). Next, assume $\gamma>1$. Using (5.4), $\rho_{2}\left|v_{2}\right\rangle=\rho_{2}\left|v_{3}\right\rangle=0$, and

$$
\begin{aligned}
\left\langle v_{1}\right| \rho_{2}\left|v_{1}\right\rangle & =\left(1-\alpha_{1}\right)\left(1-\alpha_{2}\right)+\left(\frac{\beta_{1} \beta_{2}}{\alpha_{1} \alpha_{2}}\right)^{2} \alpha_{1} \alpha_{2}-\frac{\beta_{1} \beta_{2}}{\alpha_{1} \alpha_{2}} \cdot 2 \beta_{1} \beta_{2} \\
& =\left(1-\alpha_{1}\right)\left(1-\alpha_{2}\right)-\frac{\left(\beta_{1} \beta_{2}\right)^{2}}{\alpha_{1} \alpha_{2}}=0,
\end{aligned}
$$

we have

$$
\operatorname{Tr} \rho_{2} T_{1}=\frac{1}{2}\left\langle v_{1}\right| \rho_{2}\left|v_{1}\right\rangle+\frac{\gamma-1}{2 \gamma}\left\langle v_{2}\right| \rho_{2}\left|v_{2}\right\rangle+\frac{\gamma-1}{2 \gamma}\left\langle v_{3}\right| \rho_{2}\left|v_{3}\right\rangle=0 .
$$

Moreover, since $\rho_{2}=\left(U^{[1]} \otimes U^{[2]}\right)^{*} \rho_{1}\left(U^{[1]} \otimes U^{[2]}\right)$ holds, we obtain

$$
\begin{aligned}
\operatorname{Tr} \rho_{1} T_{2} & \stackrel{5.77}{=} \operatorname{Tr} \rho_{1}\left(U^{[1]} \otimes U^{[2]}\right) T_{1}\left(U^{[1]} \otimes U^{[2]}\right)^{*} \\
& =\operatorname{Tr}\left(U^{[1]} \otimes U^{[2]}\right)^{*} \rho_{1}\left(U^{[1]} \otimes U^{[2]}\right) T_{1}=\operatorname{Tr} \rho_{2} T_{1}=0 .
\end{aligned}
$$

Therefore, (iii) holds.
Proof of Lemma 5.3. Assume that $s \in[0,1 / 4]$ and (5.3). All we need is to show that
(i) $T_{1}+T_{2}+\Gamma\left(T_{1}+T_{2}\right)=I_{4}$,
(ii) $\operatorname{neg}\left(T_{i}+\Gamma\left(T_{i}\right)\right) \leq s \operatorname{Tr}\left(T_{i}+\Gamma\left(T_{i}\right)\right)$ for $i=1,2$,
(iii) $\operatorname{Tr} \rho_{1} T_{2}=\operatorname{Tr} \rho_{2} T_{1}=0$,
by the same reason as the proof of Lemma 5.4. Since (i) and (iii) have been already proved, we show only (ii). Also, the inequality $\gamma=\alpha_{1}+\alpha_{2} \geq 1$ holds, and we may assume $\alpha_{1} \alpha_{2}>0$, by the same reason as the proof of Lemma 5.4. Moreover, (5.7) and (i) yield that $\operatorname{Tr} T_{1}=\operatorname{Tr} T_{2}$ and $\operatorname{Tr} T_{1}+\operatorname{Tr} T_{2}=2$, which implies $\operatorname{Tr} T_{1}=\operatorname{Tr} T_{2}=1$. Proof of (ii). First, assume $\gamma=1$. Then

$$
\alpha_{1} \alpha_{2} \stackrel{\gamma=1}{=}\left(1-\alpha_{2}\right)\left(1-\alpha_{1}\right) \stackrel{\boxed{55.3}}{\leq} 16 s^{2} \alpha_{1} \alpha_{2} \stackrel{s \in[0,1 / 4]}{\leq} \alpha_{1} \alpha_{2},
$$

which implies $s=1 / 4$. Since it is easily checked that $\operatorname{neg}\left(T_{i}+\Gamma\left(T_{i}\right)\right)=1 / 2$ for $i=1$, 2 , we obtain (ii). Next, assume $\gamma>1$. Then

$$
\begin{aligned}
\operatorname{neg}\left(T_{1}+\Gamma\left(T_{1}\right)\right) & \leq \operatorname{neg}\left(\Gamma\left(T_{1}\right)\right) \stackrel{\sqrt[5.4 \mid]{\leq}}{\leq} \operatorname{neg}\left(\Gamma\left(\frac{1}{2}\left|v_{1}\right\rangle\left\langle v_{1}\right|\right)\right) \\
& \stackrel{5.11}{=} \frac{1}{2}\left\|v_{1}\right\|^{2} \operatorname{sc}\left(v_{1}\right) \stackrel{(5.5)}{=} \frac{\beta_{1} \beta_{2}}{2 \alpha_{1} \alpha_{2}} \stackrel{\sqrt[5.3 \mid]{\leq}}{\leq} 2 s=s \operatorname{Tr}\left(T_{1}+\Gamma\left(T_{1}\right)\right)
\end{aligned}
$$

By (5.7), we have

$$
\operatorname{neg}\left(T_{2}+\Gamma\left(T_{2}\right)\right) \leq \operatorname{neg}\left(\Gamma\left(T_{2}\right)\right)=\operatorname{neg}\left(\Gamma\left(T_{1}\right)\right) \leq 2 s=s \operatorname{Tr}\left(T_{2}+\Gamma\left(T_{2}\right)\right)
$$

Therefore, (ii) holds.

## Chapter 6

## Conclusion

We have studied three topics associated with quantum systems and locally quantum systems: (i) capacity, (ii) differential privacy, and (iii) perfect discrimination of two states.

### 6.1 Statement determining the capacities of GPTs, and weaker statement

As stated in Section 1.1, statement S is useful to determine the capacities of special locally quantum systems. We have shown statement WS that is weaker than statement $S$, but we still do not know whether statement $S$ holds. If statement $S$ is true, it is a separability criterion in the $n$-partite case. Other separability criteria are known, e.g., the positive partial transpose criterion (see Section 5.1) and $k$-extendability [14]. However, both of them are criteria in the bipartite case. In general, it is difficult to prove a separability criterion in the $n$-partite case.

In Chapter 3, we have shown that the maximum dimension of subspaces of $\mathcal{F}^{d_{1}} \otimes$ $\cdots \otimes \mathcal{F}^{d_{n}}$ with no product basis is equal to $\tilde{d}-2$ if either (i) $n=2$ or (ii) $n \geq 3$ and $\# \mathcal{F}>\max \left\{d_{i}: i \neq n_{1}, n_{2}\right\}$ for some $n_{1}$ and $n_{2}$. When $n \geq 3$, assumption (ii) is maybe unnecessary, but we do not know its proof. Also, we have shown that a subspace $\mathcal{L}$ of $\mathcal{F}^{d_{1}} \otimes \cdots \otimes \mathcal{F}^{d_{n}}$ contains a $\left(\tilde{d} \tilde{\tilde{d}}^{n}\right)$-dimensional subspace with no product basis if the dimension of $\mathcal{L}$ is equal to $\tilde{d}-r$ and if $\mathcal{F}$ is infinite. This is best possible for $r=0,1$, but we do not know the case $r \geq 2$, which is an interesting problem.

### 6.2 Differential privacy

In Chapter 4, we have investigated the difference between the sets $\mathrm{EC}_{n}(\varepsilon)$ and $\mathrm{CQ}_{n}(\varepsilon)$. Elements in $\mathrm{CQ}_{n}(\varepsilon)$ correspond to $\mathrm{CQ} \varepsilon$-DP mechanisms in the local privacy context, and elements in $\mathrm{EC}_{n}(\varepsilon)$ correspond to essentially classical ones. Although we have not fixed the dimension $d$ of the finite-dimensional vector spaces $\mathbb{C}^{d}$ and $\mathbb{R}^{d}$, it is also important to study the case when $d$ is fixed. For instance, it is an interesting problem to find extreme points of $\mathrm{CQ}_{n}^{(d)}(\varepsilon)$. For the classical case, Holohan et al. 29] studied extreme points of $\mathrm{C}_{n}^{(d)}(\varepsilon)$. Also, CQ-DP mechanisms are expected to be superior to classical ones in information processing. Hence, we hope that such information processing will be found in the near future.

We have used Lemma 4.13 to construct CQ $\varepsilon$-DP $n$-tuples that do not lie in $\mathrm{EC}_{n}(\varepsilon)$. Instead of Lemma 4.13, one might use symmetric, informationally complete, positive-operator-valued measures (SIC-POVMs). In this case, one can probably prove that the CQ- $\varepsilon$-DP $d^{2}$-tuple $\left(\rho_{i}\right)_{i=1}^{d^{2}}$ of density matrices on $\mathbb{C}^{d}$ constructed by a SIC-POVM does not lie in $\mathrm{EC}_{d^{2}}(\varepsilon)$. However, we can prove this statement only for large $\varepsilon>0$ if using Corollary 4.6. Hence, one needs an alternative criterion instead of Corollary 4.6 to prove the above statement.

### 6.3 Perfect discrimination of two states

In Chapter 5, we have constructed two kinds of one-parameter family of proper cones in a certain natural manner. However, there are other natural manners to construct one-parameter families. Given a concrete one-parameter family, one will show a similar result to Theorems 1.8 and 1.9 . However, it is difficult to show such a result in the general case. Hence, we need another idea for further research. For instance, it is an interesting problem whether equivalence E can characterize quantum theory in some sense. Since this problem can be regarded as that to characterize a symmetric cone, it might also be mathematically interesting.

Since we have had no chance to introduce symmetric cones and Jordan algebras, we mention them briefly here. A convex cone $\mathcal{K}$ of a real Hilbert space $\mathcal{V}$ is called symmetric if $\mathcal{K}$ is

- self-dual, i.e., $\mathcal{K}^{*}=\mathcal{K}$ and
- homogeneous, i.e., $\forall x, x^{\prime}$ in the interior of $\mathcal{K}, \exists g \in G(\mathcal{K}), g x=x^{\prime}$,
where the $\operatorname{group} G(\mathcal{K})$ is defined as

$$
G(\mathcal{K})=\{g \in \mathrm{GL}(\mathcal{V}): g \mathcal{K}=\mathcal{K}\} .
$$

Table 6.1: Euclidean Jordan algebras.

| Symmetric cone | Euclidean <br> Jordan algebra | Multiplication | Inner product |
| :---: | :---: | :---: | :---: |
| Orthant <br> $[0, \infty)^{d}$ | $\mathbb{R}^{d}$ | $x \circ y=\left[\begin{array}{c}x(1) y(1) \\ \vdots \\ x(d) y(d)\end{array}\right]$ | $\sum_{i=1}^{d} x(i) y(i)$ |
| $\operatorname{PSD}(d)$ | $\operatorname{Herm}(d)$ | $X \circ Y=(X Y+Y X) / 2$ | $\operatorname{Tr} X Y$ |
| Lorentz cone |  | $\left[\begin{array}{c}\xi \\ x\end{array}\right] \circ\left[\begin{array}{c}\eta \\ y\end{array}\right]$ |  |
| $\left\{\left[\begin{array}{c}\xi \\ x\end{array}\right]:\\|x\\| \leq \xi\right\}$ | $\mathbb{R}^{1+d}$ | $=\frac{1}{\sqrt{2}\left[\begin{array}{c}\xi \eta+x^{\top} y \\ \xi y+\eta x\end{array}\right]} \quad$ |  |

Mathematically, a symmetric cone $\mathcal{K}$ is characterized by a Euclidean Jordan algebra $\mathcal{V}$ [39, 50 :
$\mathcal{K}$ symmetric cone $\Longleftrightarrow \exists \mathcal{V}$ Euclidean Jordan algebra, $\mathcal{K}=\left\{x^{2}: x \in \mathcal{V}\right\}$.
A real Hilbert space $\mathcal{V}$ with multiplication is called a Euclidean Jordan algebra if $\mathcal{V}$ satisfies the following conditions.

- $\mathcal{V}$ has an identity element.
- Distributive property: $x(y+z)=x y+x z$ and $(y+z) x=y x+z x$ for all $x, y, z \in \mathcal{V}$.
- Compatible property with scalars: $(\alpha x)(\beta y)=\alpha \beta(x y)$ for all $x, y \in \mathcal{V}$ and $\alpha, \beta \in \mathbb{R}$.
- Commutative property: $x y=y x$ for all $x, y \in \mathcal{V}$.
- Jordan identity: $x(y(x x))=(x y)(x x)$ for all $x, y \in \mathcal{V}$.
- Euclidean property: $\langle x, y z\rangle=\langle x y, z\rangle$ for all $x, y, z \in \mathcal{V}$.

Table 6.1 has three types of Euclidean Jordan algebras. Every Euclidean Jordan algebra can be decomposed to simple ones, and the classification of simple Euclidean

Jordan algebras is known 37]. In this way, convex cones are often studied in mathematics with an algebraic view. However, it is maybe more useful to study a topic like perfect discrimination in order to further understand GPTs. We hope that such studies will increase much more.

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[^0]:    ${ }^{1}$ (C)2020 IEEE
    ${ }^{2}$ (C)2020 American Physical Society

[^1]:    ${ }^{1}$ Actually, the convex cones $[0, \infty)^{d}$ and $\operatorname{PSD}(d)$ are special cones called symmetric cones. Hence, GPTs are also related to symmetric cones and Euclidean Jordan algebras 6, 37, 39, 44, 50.

[^2]:    ${ }^{2}$ This is consistent with the definition of $\mathcal{K}_{\text {max }}$, since $\mathcal{K}^{* *}=\mathcal{K}$ for every non-empty closed convex cone 8, p. 53], [20, p. 63].

[^3]:    ${ }^{1}$ [53 © 2020 American Physical Society

