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# Efficient Classical Simulation of a Variant of Cluster State Quantum Computation 

Sahar Atallah

1830214

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Brunel University London
Department of Mathematics
Uxbridge
Middlesex
UB8 3PH
United Kingdom
T: +44 1895203397
F: +44 (0) 1895251686

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## Abstract

Quantum computers are known for their ability to solve some computational problems faster than classical computers. There is a race to build quantum computers because it is believed they might be better than classical; but it remains unknown whether quantum computers are in fact better than conventional computers. To understand this problem, we develop a new method of classically simulating certain types of quantum system that are previously unknown to be efficiently simulatable on classical computers.

We adjust a part of cluster state quantum computation to study the computational power and we demonstrate that there is a finite region of pure states $|\psi\rangle$ around the $Z$-eigenstates for which the setup can be efficiently simulated classically, given that the measurements are limited to $Z$ and $X-Y$ plane measurements. This classical simulation works by considering alternative local state spaces that we called "cylinders" and different notion of entanglement to normal quantum entanglement.

Then, we work out similar regions for states created using other diagonal gates instead of the $C Z$. These diagonal gates are represented by $V(\theta)=|0\rangle\langle 0| \otimes I+|1\rangle\langle 1| \otimes$ $Z_{\theta}$ where $Z_{\theta}=|0\rangle\langle 0|+e^{i \theta}|1\rangle\langle 1|$. It turns out that almost all inputs are classically simulatable when $\theta$ is small.

In addition, we find that classical simulation also works by considering new type of non-quantum state spaces other than cylinders and maintaining non-entangled representation by growing the size of these state spaces. We search over some state spaces to try optimize our classical simulation and it turns out that, among the state spaces that we searched through, the cylinder is the most optimal state space.

And finally, we will look at a coarse graining version of construction which increases the efficiently simulatable region.

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## Chapter 1

## Introduction

Quantum computing is the use of quantum phenomena, such as the so-called entanglement and superposition, to perform computation. Computers that performs quantum computations are known as quantum computers. These computers are able to solve some computational problems, such as factorization in polynomial time, exponentially faster than the best known algorithm on classical computers. Algorithms that are performed on quantum computers, are called quantum algorithm [22, 23, $44,42,56]$. For example, Shor [57] who works on factoring integers; this problem is believed to be hard on a classical computer but efficient algorithms are given for this problem on a quantum computer. An additional example is that Grover [30] discovers an algorithm that can solve a problem on a quantum computer using quadratic speedup operations.

Before mentioning the goal of the research, let's start by revealing why we are interested in simulating quantum systems on a classical computer. The first reason is that it is unknown if quantum computers are better than conventional computers. To understand this problem, we can proceed one of two ways. We can try to find that quantum computers can't be efficiently simulated classically or we can try to develop an algorithm that can classically simulate a quantum system.

The motivation of simulating a complex quantum system efficiently classically is that people think that quantum computers, built with quantum systems, can solve problems such as factorization that can't be done classically; but if it is found that simulating quantum computer can be done classically, then it means that any algorithm can be done on quantum computers, it can be efficiently simulated classically. For instance, there would be an efficient algorithm for factorization on a classical computer.

Progress has been made in efficiently simulating a quantum system on a classical computer. For example, Gottesman-Knill theorem [27] shows that a quantum computation, that involves only the state preparations in the computational basis,

Hadamard gates, Phase gates, Controlled-NOT gates, Pauli gates and measurements in the Pauli group, may be efficiently simulated on a classical computer. So, it is an algorithm that only simulate certain type of quantum system. Gottesman-Knill theorem displays that some quantum computations, including entangled states (but not all types of entanglement) may be efficiently simulated classically. Hence, GottesmanKnill theorem is very interesting because there are connections to error correction and only single particle operations are needed to turn stabilizer circuits into a form of quantum computation.

Most people in this field believe that it is not possible to efficiently simulate quantum systems classically. One reason is that factorization can be done efficiently on quantum computer, in spite of the fact that over many years people have failed to find an efficient factorization algorithm on a classical computer. But there is no proof of this fact nor the fact that quantum system cannot be efficiently simulated classically.

In this thesis, we aim to develop a new method of classically simulating certain types of quantum system that are previously unknown to be efficiently simulatable on classical computers.

### 1.1 Basic mathematics

In this section, we are going to explain the definitions and the formulas of all the mathematical terms used in this research.

### 1.1.1 Hilbert space

We begin with the definition of the fundamental mathematical concept, "the Hilbert Space" which is denoted by $H$.

The Hilbert spaces permit generalising the linear algebra and calculus methods from the two and three dimensional Euclidean spaces to other dimensions spaces and to spaces that have an infinite dimension, but we will only be concerned with finite dimensional systems. Hilbert space is a complex inner product space which means that it is a complex vector space endowed with an inner product operation, we will shortly explain the meaning of these terms.

The vector state space is a set of elements, called vectors such that the set must satisfies the following properties:

1. Any vector $|v\rangle$ can be multiplied by any number called scalar. Then the vector is denoted by $\alpha|v\rangle$, where $\alpha$ can be a real or complex scalar depending upon the vector space. If the scalar is a real number, then the space is called a real vector space, and if the scalar is a complex number, then the space is a complex vector space.
2. For any vector $\left|v_{1}\right\rangle$ and any two scalars $\alpha, \beta$, it holds that $(\alpha+\beta)\left|v_{1}\right\rangle=$ $\alpha\left|v_{1}\right\rangle+\beta\left|v_{1}\right\rangle$.
3. For any vector $\alpha|v\rangle$ and any scalar $\beta$, there is a vector denoted by $\beta(\alpha|v\rangle)=$ $(\beta \alpha)|v\rangle=(\alpha \beta)|v\rangle$.
4. The vectors can be added together. For instance, for any two vectors $\left|v_{1}\right\rangle$ and $\left|v_{2}\right\rangle$, there is a vector $|w\rangle$ denoted by $|w\rangle=\left|v_{1}\right\rangle+\left|v_{2}\right\rangle$.
5. The sum of vectors must be associative and commutative. For instance, by associative and commutative, we mean $\left(\left|v_{1}\right\rangle+\left|v_{2}\right\rangle\right)+\left|v_{3}\right\rangle=\left|v_{1}\right\rangle+\left(\left|v_{2}\right\rangle+\left|v_{3}\right\rangle\right)$ and $\left|v_{1}\right\rangle+\left|v_{2}\right\rangle=\left|v_{2}\right\rangle+\left|v_{1}\right\rangle$ respectively.
6. The multiplication by $\alpha$ is distributive such that $\alpha\left(\left|v_{1}\right\rangle+\left|v_{2}\right\rangle\right)=\alpha\left|v_{1}\right\rangle+\alpha\left|v_{2}\right\rangle$.
7. By multiplying any vector $|v\rangle$ by $1,|v\rangle$ remains the same.
8. There exists a zero vector 0 such that by adding any $|v\rangle$ to 0 , the same $|v\rangle$ will be given, i.e. $|v\rangle+0=|v\rangle$.
9. For any vector $|v\rangle$, there is an inverse for this vector denoted $-|v\rangle$, such that $|v\rangle+(-|v\rangle)=0$.

A set of vectors $\left|a_{1}\right\rangle,\left|a_{2}\right\rangle, \ldots,\left|a_{d}\right\rangle$ is a spanning set of a vector space $E$ if $\forall u \in$ $E, \exists \beta_{1}, \beta_{2}, \ldots, \beta_{d} \in \mathbb{K}$ such that $u=\beta_{1}\left|a_{1}\right\rangle+\beta_{2}\left|a_{2}\right\rangle+\ldots+\beta_{d}\left|a_{d}\right\rangle$. A set of vectors is said to be Linearly Independent if one of the following two equivalent cases holds:

- if the basis vectors give a unique expansion,
- if the zero vector 0 has a unique expansion, $0=m\left|a_{1}\right\rangle+n\left|a_{2}\right\rangle+p\left|a_{3}\right\rangle+\ldots+q\left|a_{d}\right\rangle$, where the coefficients, $m, n, p, \ldots, q$ must be zeros.

Hence, a set of vectors is a basis if the set is a linearly independent spanning set.
The total number of vectors, which is $d$ in the above basis set, represents the dimension of a vector space.

Some vector spaces can be given an inner product which is a map of two input vectors to a complex output number. Using the Dirac notation or bra-ket notation, which is a language that match the needs of expressing states in quantum mechanics, an inner product is obtained by combining two vectors, the bra and the ket vectors, denoted by $\langle g|$ and $|h\rangle$ respectively. The bra vector $\langle g|$, is a linear map that takes a vector $|h\rangle$ and gives $\langle g \mid h\rangle$ which represents the inner product.

The inner product is defined to satisfy the following axioms:

1. The inner product $\langle g \mid h\rangle$ must be equal to its adjoint $\langle g \mid h\rangle^{*}=\langle h \mid g\rangle$. This equality is called the conjugate symmetry. By adjoint, we mean that it is a linear transformation from a vector space to itself and this linear map leads to the transpose conjugate.
2. When a vector $\langle v|$ acts on a sum of two vectors $(|h\rangle+|w\rangle),\langle v \mid h\rangle+\langle v \mid w\rangle$ will be given.
3. When a vector $|h\rangle$ is multiplied by a scalar $t$, then $\langle g|(t|h\rangle)=t\langle g \mid h\rangle$.
4. For all vectors $|h\rangle,\langle h \mid h\rangle \geq 0$.
5. If $|h\rangle$ is a zero vector, then $\langle h \mid h\rangle=0$.
6. If $\langle h|$ acts on a zero vector $|0\rangle$, then $\langle h \mid 0\rangle=\langle 0 \mid h\rangle=0$.

If the inner product of two vectors, $|g\rangle$ and $|h\rangle$ is 0 , which means that $\langle g \mid h\rangle=$ $\langle g \mid h\rangle^{*}=\langle h \mid g\rangle=0$, then these two vectors are defined to be orthogonal.

In vector spaces, the norm, which is a non-negative real valued function, is used to define lengths of vectors. The norm which is denoted by $\||v\rangle \|$, satisfies three conditions; for instance, by taking a vector $|x\rangle$ in $H$ :

1. For all vectors $|x\rangle, \||x\rangle \| \geq 0$ which means $\|-|x\rangle\|=\||x\rangle \|$ with $\||x\rangle$ $\|=0$ if $|x\rangle=0$.
2. By multiplying the vector $|x\rangle$ by a scalar w , then the norm is given by $\| w|x\rangle \|$ $=|w| \times \||x\rangle \|$.
3. By taking two vectors $|x\rangle$ and $|y\rangle$, the norm of the sum of these two vectors is given by $\||x\rangle+|y\rangle\|\leq\||x\rangle\|+\||y\rangle \|$ which represents the triangle inequality.

The inner product automatically leads to a type of norm. Then the norm of a vector space that has an inner product is given by:

$$
\||x\rangle \|_{2}=\sqrt{\langle x \mid x\rangle}=\sqrt{\sum_{i=1}\left|x_{i}\right|^{2}} .
$$

Let's consider a basis of vectors $\left|a_{1}\right\rangle,\left|a_{2}\right\rangle,\left|a_{3}\right\rangle, \ldots$. If the inner product of two vectors

$$
\left\langle a_{i} \mid a_{j}\right\rangle=0
$$

for all $i \neq j$ and if the norm of all vectors

$$
\sqrt{\left\langle a_{i} \mid a_{i}\right\rangle}=1=\left\langle a_{i} \mid a_{i}\right\rangle
$$

then the set is an orthonormal basis.
Vectors can be described in terms of expansion coefficients in an orthonormal basis. Let's consider any two vectors $|h\rangle=\sum_{i} \alpha_{i}\left|h_{i}\right\rangle$ and $|g\rangle=\sum_{i} \beta_{i}\left|g_{i}\right\rangle$ where $\left|g_{i}\right\rangle=\left|h_{i}\right\rangle$, then the inner product between the two vectors is given by:

$$
\begin{aligned}
\langle g \mid h\rangle & =\sum_{i} \beta_{i}^{*} \alpha_{i}\left\langle g_{i} \mid h_{i}\right\rangle \\
& =\sum_{i} \beta_{i}^{*} \alpha_{i} ; \text { where }\left\langle g_{i} \mid h_{i}\right\rangle=1 \text { because they are normalised. }
\end{aligned}
$$

The inner product can be written as a product of two matrices:

$$
\langle g \mid h\rangle=\left(\begin{array}{lll}
\beta_{1}^{*} & \beta_{2}^{*} & \ldots
\end{array}\right)\left(\begin{array}{c}
\alpha_{1} \\
\alpha_{2} \\
\vdots
\end{array}\right)
$$

Therefore, the two vectors $|h\rangle$ and $|g\rangle$, also can be written as matrices. $|h\rangle$ which is the ket vector, is a column vector in a complex vector space and the bra vector $\langle g|$ is the adjoint of $|g\rangle$, these vectors are given by

$$
|h\rangle=\left(\begin{array}{c}
h_{1} \\
h_{2} \\
\vdots
\end{array}\right)
$$

and

$$
\langle g|=\left(\begin{array}{lll}
g_{1}^{*} & g_{2}^{*} & \ldots
\end{array}\right)
$$

### 1.1.2 Operator

An operator $O$ maps input vectors into output vectors that are in two different Hilbert spaces, $O: H_{\text {in }} \rightarrow H_{\text {out }}$. For instance, we denote this transformation by the following equation

$$
|w\rangle=O|v\rangle
$$

- An operator is linear, if the output vector $|w\rangle$ is the sum of all the output vectors, if $|v\rangle=a\left|v_{1}\right\rangle+b\left|v_{2}\right\rangle+\ldots$, then

$$
|w\rangle=O|v\rangle=a O\left|v_{1}\right\rangle+b O\left|v_{2}\right\rangle+\ldots
$$

- By considering an operator $A$ that maps two vectors in two different Hilbert spaces, from $H_{1} \rightarrow H_{2}$, its adjoint operator denoted $A^{\dagger}$ maps from $H_{2} \rightarrow H_{1}$ such that $\left\langle A v_{1} \mid v_{2}\right\rangle=\left\langle v_{1} \mid A^{\dagger} v_{2}\right\rangle$. We will only be concerned with an adjoint operator that maps from a Hibert space to itself.
- A Hermitian operator is an operator $A$ that satisfies $A^{\dagger}=A$ where $A^{\dagger}$ is the adjoint operator of $A$ which also can be called Hermitian conjugate.
- Unitary operators $U$ are linear transformation that maps an input vector $|\phi\rangle$ in $H_{1}$ to an output vector $|\psi\rangle$ in $H_{2}$. The two Hilbert spaces $H_{1}$ and $H_{2}$, must be equivalent, which means $H_{1}=H_{2}$ and $U$ has an inverse $U^{\dagger}$ which is the Hermitian conjugate, then

$$
U^{\dagger} U=I
$$

- A projector or a projection operator $P$ is a linear operator that maps a vector in vector space $V$ to itself such that $P^{2}=P$ which means that if $P^{2}$ is applied to a vector, the same outcome is obtained as if only one $P$ is applied.

Operators can be represented as matrices when their actions are expressed in terms of orthonormal basis. Let's consider a linear operator $O$ that maps an input vector $|v\rangle$ to an output vector $|w\rangle$. The input vector $|v\rangle$ can be expanded in terms of
an orthonormal basis:

$$
\begin{aligned}
|v\rangle & =\sum_{i} \alpha_{i}\left|v_{i}\right\rangle \\
& =\alpha_{1}\left|v_{1}\right\rangle+\alpha_{2}\left|v_{2}\right\rangle+\ldots,
\end{aligned}
$$

The coefficients column vector is given by

$$
\left(\begin{array}{c}
\alpha_{1} \\
\alpha_{2} \\
\vdots
\end{array}\right)
$$

After applying the operator on $|v\rangle$, it is the same as applying $O$ on each vector $\left|v_{1}\right\rangle,\left|v_{2}\right\rangle, \ldots$, these output vectors that also can be expanded as a sum of basis vectors, are presented by

$$
\begin{aligned}
|w\rangle & =O|v\rangle \\
& =\sum_{i} \alpha_{i} O\left|v_{i}\right\rangle \\
& =\alpha_{1} O\left|v_{1}\right\rangle+\alpha_{2} O\left|v_{2}\right\rangle+\ldots \\
& =\alpha_{1}\left(\beta_{1,1}\left|w_{1}\right\rangle+\beta_{2,1}\left|w_{2}\right\rangle+\ldots\right)+\alpha_{2}\left(\beta_{1,2}\left|w_{1}\right\rangle+\beta_{2,2}\left|w_{2}\right\rangle+\ldots\right)+\ldots \\
& =\alpha_{1}\left(\sum_{j} \beta_{j, 1}\left|w_{j}\right\rangle\right)+\alpha_{2}\left(\sum_{j} \beta_{j, 2}\left|w_{j}\right\rangle\right)+\ldots \\
& =\sum_{i} \alpha_{i}\left(\sum_{j} \beta_{j, i}\left|w_{j}\right\rangle\right) \\
& =\sum_{j}\left(\sum_{i} \beta_{j, i} \alpha_{i}\right)\left|w_{j}\right\rangle \\
& =\sum_{j} \gamma_{j}\left|w_{j}\right\rangle
\end{aligned}
$$

where $\gamma_{j}=\sum_{i} \beta_{j, i} \alpha_{i}$, it can be written as

$$
\left(\begin{array}{c}
\gamma_{1} \\
\gamma_{2} \\
\gamma_{3} \\
\vdots
\end{array}\right)=\left(\begin{array}{cccc}
\beta_{1,1} & \beta_{1,2} & \beta_{1,3} & \cdots \\
\beta_{2,1} & \beta_{2,2} & \beta_{2,3} & \ldots \\
\beta_{3,1} & \beta_{3,2} & \beta_{3,3} & \ldots \\
\vdots & & & \\
& \ldots & &
\end{array}\right)\left(\begin{array}{c}
\alpha_{1} \\
\alpha_{2} \\
\alpha_{3} \\
\vdots
\end{array}\right)
$$

Then $|w\rangle$ can be represented by

$$
(|w\rangle)=\left(\begin{array}{cccc} 
& & & \\
& \\
\left|w_{1}\right\rangle & \left|w_{2}\right\rangle & \left|w_{3}\right\rangle & \ldots \\
& & & \\
&
\end{array}\right)
$$

Hence, the linear operator $O$ can be represented by a matrix that takes input column vectors and gives output column vectors.

### 1.1.3 Tensor product

The tensor product is a map that take a pair of vector spaces, $V$ and $W$ with dimensions $c$ and $d$ respectively to a single vector space $V \otimes W$ that has a dimension of $c d$. The sign " $\otimes$ " is the tensor product symbol.

The tensor product satisfies these axioms:

1. By taking two vectors $|v\rangle \in V$ and $|w\rangle \in W$, it holds that: $a(|v\rangle \otimes|w\rangle)=$ $a(|v\rangle) \otimes|w\rangle=|v\rangle \otimes a(|w\rangle)$ where $a$ is a scalar.
2. By adding two vectors $\left|v_{1}\right\rangle+\left|v_{2}\right\rangle$, for a vector $|w\rangle \in W$, we have $\left(\left|v_{1}\right\rangle+\left|v_{2}\right\rangle\right) \otimes$ $|w\rangle=\left|v_{1}\right\rangle \otimes|w\rangle+\left|v_{2}\right\rangle \otimes|w\rangle$.
3. If we consider a vector $|v\rangle \in V$ and $\left(\left|w_{1}\right\rangle+\left|w_{2}\right\rangle\right)$, then $|v\rangle \otimes\left(\left|w_{1}\right\rangle+\left|w_{2}\right\rangle\right)=$ $|v\rangle \otimes\left|w_{1}\right\rangle+|v\rangle \otimes\left|w_{2}\right\rangle$.

Let's consider $|v\rangle$ and $|w\rangle$ two column vectors in vector space $V$ and $W$. These
two vectors have $x_{i}$ and $y_{i}$ as components where $i=1,2, \ldots, n$,

$$
|v\rangle=\left(\begin{array}{c}
x_{1} \\
x_{2} \\
\vdots \\
x_{n}
\end{array}\right) \in \mathbb{R}^{n} \text { and }|w\rangle=\left(\begin{array}{c}
y_{1} \\
y_{2} \\
\vdots \\
y_{n}
\end{array}\right) \in \mathbb{R}^{n}
$$

When a tensor product is applied between these two vectors, a joint vector has been formed:

$$
|v\rangle \otimes|w\rangle=\left(\begin{array}{c}
x_{1} \times\left(\begin{array}{c}
y_{1} \\
y_{2} \\
\vdots \\
y_{n}
\end{array}\right) \\
\vdots \\
x_{n} \times\left(\begin{array}{c}
y_{1} \\
y_{2} \\
\vdots \\
y_{n}
\end{array}\right)
\end{array}\right)=\left(\begin{array}{c}
x_{1} y_{1} \\
x_{1} y_{2} \\
\vdots \\
x_{1} y_{n} \\
\vdots \\
x_{n} y_{1} \\
x_{n} y_{2} \\
\vdots \\
x_{n} y_{n}
\end{array}\right)
$$

In general, if we have $E \in M_{1}$ and $F \in M_{2}$, where $M_{1}$ and $M_{2}$ are the spaces of $c \times d$ and $m \times n$ matrices respectively, then $E \otimes F$ is a $c m \times d n$ matrix.

$$
E \otimes F=\left(\begin{array}{cccc}
E_{11} F & E_{12} F & \ldots & E_{1 d} F \\
E_{21} F & E_{22} F & \ldots & E_{2 d} F \\
\vdots & \vdots & \vdots & \vdots \\
E_{c 1} F & E_{c 2} F & \ldots & E_{c d} F
\end{array}\right)
$$

For example, by taking $c=d=3$ then

$$
E \otimes F=\left(\begin{array}{lll}
E_{11} F & E_{12} F & E_{13} F \\
E_{21} F & E_{22} F & E_{23} F \\
E_{31} F & E_{32} F & E_{33} F
\end{array}\right)
$$

To simplify many calculations, the fact about tensor product is that

$$
(E \otimes F)(G \otimes H)=E G \otimes F H
$$

### 1.1.4 Quantum states

An observable is a physical property of a system that can be measured in the lab. In quantum mechanics, an observable is described by a Hermitian operator and the eigenvalues of this operator represent the possible outcomes of measurement that correspond to the vectors in the orthonormal basis. The observable is said to be degenerate when eigenvectors correspond to one eigenvalue and when there is one eigenvector for only one eigenvalue then the observable is non-degenerate.

A projective measurement can be constructed from an observable denoted $M$, which is a Hermitian operator. The observable $M$ has a decomposition:

$$
\begin{equation*}
M=\sum_{m} m P_{m} \tag{1.1}
\end{equation*}
$$

where $P_{m}$ is the projector on the eigenspace of the observable that has an eigenvalue $m$. By eigenspace we mean the set of eigenvectors that correspond to the same eigenvalue.

In quantum information theory, we are only interested in labelling these $m$ eigenvalues, then we are going to use $P_{i}$ where $i$ are integer numbers that represent the possible outcomes. $P_{i}$ must satisfy the following properties:

- $P_{i}^{\dagger}=P_{i}$,
- $P_{i} P_{i^{\prime}}=\delta_{i, i^{\prime}} P_{i}$
- $\sum_{i}^{d} P_{i}=I$,

The probability of getting a measurement outcome $i$ is given by the following Born Rule.

## The Born Rule

The Born Rule or the Probability Rule is the rule that gives the probabilities of getting possible outcomes of measurement.

Let's consider an operator $O$ that has real eigenvalues $\lambda_{i}$ that correspond to eigenvectors $\left|\lambda_{i}\right\rangle$, then the operator $O$ can be defined by

$$
O=\sum_{i} \lambda_{i} P_{i}
$$

where $P_{i}=\sum_{k}\left|\lambda_{k}\right\rangle\left\langle\lambda_{k}\right|$ is the projector. If a measurement acts on a quantum system that is in a state $|\psi\rangle$, the probability of getting an outcome $i$ is presented by:

$$
\begin{equation*}
\operatorname{prob}(i)=\langle\psi| P_{i}|\psi\rangle \tag{1.2}
\end{equation*}
$$

Instead of taking a single state $|\psi\rangle$, we consider an ensemble of state $\left|\psi_{j}\right\rangle$ with probabilities $p_{j}$, then the probability of getting an outcome $i$ is given by:

$$
\begin{equation*}
\operatorname{prob}(i)=\sum_{j} p_{j}\left\langle\psi_{j}\right| P_{i}\left|\psi_{j}\right\rangle \tag{1.3}
\end{equation*}
$$

We can derive from equation (1.3) the following:

$$
\begin{aligned}
\operatorname{prob}(i) & =\sum_{j} p_{j}\left\langle\psi_{j}\right| P_{i}\left|\psi_{j}\right\rangle \\
& =\sum_{j} p_{j} \sum_{k}\left\langle\psi_{j} \mid \lambda_{k}\right\rangle\left\langle\lambda_{k} \mid \psi_{j}\right\rangle \\
& =\sum_{j} p_{j} \sum_{k}\left\langle\lambda_{k} \mid \psi_{j}\right\rangle\left\langle\psi_{j} \mid \lambda_{k}\right\rangle \\
& =\sum_{k}\left\langle\lambda_{k}\right| \sum_{j} p_{j}\left|\psi_{j}\right\rangle\left\langle\psi_{j} \mid \lambda_{k}\right\rangle \\
& =\sum_{k}\left\langle\lambda_{k}\right| \rho\left|\lambda_{k}\right\rangle \\
& =\operatorname{tr}\left(\rho \sum_{k}\left|\lambda_{k}\right\rangle\left\langle\lambda_{k}\right|\right) \\
& =\operatorname{tr}\left(\rho P_{i}\right)
\end{aligned}
$$

where "tr" is the trace of an operator matrix. $\rho=\sum_{j} p_{j}\left|\psi_{j}\right\rangle\left\langle\psi_{j}\right|$ is the density matrix (or a density operator) which describes the quantum state and it has the following properties:

- Hermitian, it means that $\rho=\rho^{\dagger}$,
- a positive semi-definite which means that the eigenvalues are non-negative. It can be denoted as $\rho \geq 0$,
- has trace one.

Using the Born's rule, a quantum state is a mathematical structure that allows us to calculate probabilities of the outcomes of each possible measurement on a system.

States that can't be written as a mixture of states are called pure quantum states but others are called mixed quantum states.

If $\rho=\sum_{j}^{n} p_{j}\left|\psi_{j}\right\rangle\left\langle\psi_{j}\right|=|\psi\rangle\langle\psi|$, where $n=1$ and $p_{1}=1$, then $\rho$ represents a pure quantum state; but if it is not the case then $\rho$ is called a mixed quantum state. This is the first way in which mixed states can arise and we will shortly discuss the second way, which is the partial trace of a pure state, that can be used to obtain a mixed state.

Pure quantum state is described by a single vector $|k\rangle$ or by a superposition of basis states $|\Psi\rangle=\sum_{i} c_{i}\left|k_{i}\right\rangle$ where $c_{i}$ are the probabilities corresponding to $\left|k_{i}\right\rangle$ and $\sum_{i}\left|c_{i}\right|^{2}=1$.

$$
\text { If } \operatorname{tr}\left(\rho^{2}\right)\left\{\begin{array}{l}
=1, \text { then } \rho \text { is a pure state } \\
<1, \text { then } \rho \text { is a mixed state. }
\end{array}\right.
$$

## The Qubit

The classic binary bit has a quantum version which is called quantum bit or qubit. Qubit is in a two dimensional Hilbert space and it can be in states labelled 0, 1, or a superposition of these. The superposition stated that any two or more quantum states can be added together, we end up with another quantum state, for example, $|\Psi\rangle=\alpha|0\rangle+\beta|1\rangle$, where $\alpha$ and $\beta$ are two complex numbers and $|0\rangle$ and $|1\rangle$ are the ket- 0 and ket-1 respectively. So, The qubit is represented by a column vector of two elements, where these elements should satisfy the normalisation condition, from the previous example, $|\alpha|^{2}+|\beta|^{2}=1$. By using the Dirac notation, where the ket is $|0\rangle$, $|1\rangle$ and the bra is the conjugate transpose, we have: $|0\rangle:=\binom{1}{0} ;|1\rangle:=\binom{0}{1}$ and these two are the orthonormal basis states. $|0\rangle,|1\rangle$ together is called computational basis.

For instance, the qubit state $|+\rangle$ is a superposition of the basis state which means that it can be described by a linear combination of $|0\rangle$ and $|1\rangle$ :

## Multi-particle system

A multiparticle system can be defined by two axioms:

1. Tensor product is crucial in the understanding of multiparticle systems. If there are two Hilbert space $H_{A}$ and $H_{B}$ on two different particles $A$ and $B$ respectively, then a Hilbert space $H_{A B}$ is described by a tensor product between these two spaces. $H_{A B}$ has two different types of states:

- a product state which is given by,

$$
\left|\phi_{1}\right\rangle \otimes\left|\phi_{2}\right\rangle
$$

where $\left|\phi_{1}\right\rangle$ and $\left|\phi_{2}\right\rangle$ are two independent states.

- a joint state that has the possibility of being an entangled state. This specific state cannot be written as a product state, as we will shortly explain.

2. If a measurement will be applied on one particle, then in the joint system, the measurement is extended by the identity operator which is applied on the second particle.

For instance, if a measurement $P_{j}$ is applied on $A$, then we need to apply a tensor product between $P_{j_{A}}$ and the identity operator $I_{B}$ which is applied on the second particle $B$. The probability of getting the measurement outcome $j$ is given by

$$
\begin{equation*}
\operatorname{Pr}(j)=\operatorname{tr}\left(\left(P_{j_{A}} \otimes I_{B}\right) \rho_{A B}\right)=\operatorname{tr}\left(P_{j_{A}} \operatorname{tr}_{B}\left(\rho_{A B}\right)\right)=\operatorname{tr}\left(P_{j_{A}} \rho_{A}\right) \tag{1.4}
\end{equation*}
$$

where $\rho_{A}$ is defined as

$$
\rho_{A}:=\operatorname{tr}_{B}\left(\rho_{A B}\right):=\sum_{k}^{d_{B}}\left(I_{A} \otimes\left\langle k_{B}\right|\right) \rho_{A B}\left(I_{A} \otimes\left|k_{B}\right\rangle\right)
$$

$d_{B}$ is the dimension of $H_{B},\left|k_{B}\right\rangle$ is any orthonormal basis on $B$ and $I_{A}$ is the identity operator in $H_{A} . \rho_{A}$ is called The reduced state or the reduced density matrix of $A$ which is computed by taking the partial trace over $B$, denoted by $t r_{B}$.

## POVM

As mentioned above, the Born's rule states that the probability of the outcome for measurement, is represented by

$$
\operatorname{prob}(i)=\sum_{j} p_{j}\left\langle\psi_{j}\right| P_{i}\left|\psi_{j}\right\rangle=\operatorname{tr}\left(\rho P_{i}\right)
$$

where $\rho=\sum_{j} p_{j}\left|\psi_{j}\right\rangle\left\langle\psi_{j}\right|$ and $P_{i}=\sum_{k}\left|\lambda_{k}\right\rangle\left\langle\lambda_{k}\right|$. The projective measurements $P_{i}$ have the axioms that we mentioned in the beginning of this part:

- $P_{i}^{\dagger}=P_{i}$,
- $\sum_{i} P_{i}=I$ where $I$ is the identity operator because the sum of all probabilities must be equal to 1 .
- $P_{i} P_{i^{\prime}}=\delta_{i, i^{\prime}} P_{i}$. This property together with the above two axioms, imply that the operator $P_{i}$ is positive and therefore the probability will be non-negative.

If we add ancilla particle which is used in the measurement, to the probability equation, these axioms can be modified. Let's try to interact the quantum system with an external one, so we are going to add an ancilla $|0\rangle\langle 0|$. Then, we have

$$
\begin{equation*}
\operatorname{Prob}(\mathrm{i})=\operatorname{tr}\left[(\rho \otimes|0\rangle\langle 0|) P_{i}\right] \equiv \operatorname{tr}(\rho M) \tag{1.5}
\end{equation*}
$$

where $P_{i}=\left|\lambda_{k}\right\rangle\left\langle\lambda_{k}\right|$ and $M=\sum_{a, b} \lambda_{b 0} \lambda_{a 0}^{*}|b 0\rangle\langle a 0|$.
Proof. By taking $\left|\lambda_{k}\right\rangle=\sum_{k, j} \lambda_{k j}|k j\rangle$, then the probability of getting $i$ is

$$
\begin{aligned}
\operatorname{Prob}(\mathrm{i}) & =\operatorname{tr}\left[(\rho \otimes|0\rangle\langle 0|) P_{i}\right] \\
& =\operatorname{tr}\left[(\rho \otimes|0\rangle\langle 0|)\left|\lambda_{k}\right\rangle\left\langle\lambda_{k}\right|\right] \\
& =\operatorname{tr}\left[\sum_{a, b, k, j, m, l}\left(\rho_{a b} \otimes|a 0\rangle\langle b 0|\right) \lambda_{k j} \lambda_{m l}^{*}|k j\rangle\langle m l|\right] \\
& =\operatorname{tr}\left[\sum_{a, b}\left(\rho_{a b} \otimes|a 0\rangle\langle b 0|\right) \lambda_{b 0} \lambda_{a 0}^{*}|b 0\rangle\langle a 0|\right] \text { by taking } k=b, m=a, j=l=0 ; \\
& =\operatorname{tr}\left[\sum_{a, b} \lambda_{b 0} \lambda_{a 0}^{*}|b 0\rangle\langle a 0| \rho_{a b}\right] \\
& =\operatorname{tr}(\rho M) ; \text { where } M=\sum_{a, b} \lambda_{b 0} \lambda_{a 0}^{*}|b 0\rangle\langle a 0|
\end{aligned}
$$

$M$ is a Positive Operator Valued Measure element because when a projective measurement $P_{i}$ is applied on a particle and an ancilla, we end up with a statistic defined by POVM where only the third axiom has been modified, which means that the orthogonality property is removed and the positivity is retained.

Therefore Positive Operator Valued Measure or Generalized measurement is described by the set of positive operators $P_{i}$ that sum to the identity matrix: $\sum_{i=1}^{n} P_{i}=I$.

## Bloch sphere

Let's consider a system of a single qubit that has a $2 \times 2$ density matrix $\rho$. This matrix is Hermitian, has trace one and can be expanded in the basis $I, \sigma_{1}, \sigma_{2}, \sigma_{3}$ where $I$ is the identity matrix, $\sigma_{1}=X, \sigma_{2}=Y$ and $\sigma_{3}=Z$ are the Pauli matrices (as will be defined on page 19). The coefficient of $I$ in the expansion of the $\rho$ is $\frac{1}{2}$ to end up with $\operatorname{tr}(\rho)=1$. The expansion is given by:

$$
\rho=\frac{1}{2}(I+x X+y Y+z Z)=\frac{1}{2}\left(\begin{array}{cc}
1+z & x-i y \\
x+i y & 1-z
\end{array}\right)
$$

The coefficients $x, y$ and $z$ are found using the Hilbert Schmidt inner product which is presented by $\operatorname{tr}\left(A^{\dagger} B\right)$. Hence, they are computed as follow

$$
\frac{1}{2} \operatorname{tr}\left(\rho \sigma_{i}\right) .
$$

Then the determinant of the $\rho$ must be computed to find the eigenvalues $\lambda_{i}$. Hence, the eigenvalues are given by

$$
\lambda_{i}=\frac{1}{2} \pm \frac{\sqrt{x^{2}+y^{2}+z^{2}}}{2} .
$$

The non-negativity of the eigenvalues will be holding if $x^{2}+y^{2}+z^{2} \leq 1$. Thus a quantum state can be represented by a vector $(x, y, z)$, where these values are the coefficients of the above decomposition. These coefficients correspond to a point on the so-called a Bloch sphere which is considered as a unit 2-sphere and the standard basis vectors $|0\rangle$ and $|1\rangle$ are represented by the north and south poles of the Bloch sphere respectively [50].

The pure quantum states of the system are the points on the surface of the sphere where the norm of the Bloch vector must be equal to 1 , while the internal points correspond to the mixed states where the norm of the vector must be less than 1 .


Figure 1.1: This image represents the Bloch sphere. It is taken from "Wikipedia", https://commons.wikimedia.org/wiki/File:Bloch_sphere.svg, February 2022.

## Entangled and separable quantum states

Let's consider two quantum systems $A$ and $B$ in two different Hilbert spaces $H_{A}$ and $H_{B}$. If the pure state $\left|\chi_{A B}\right\rangle \in H_{A} \otimes H_{B}$ can be written in the form $\left|\chi_{A B}\right\rangle=\left|\chi_{A}\right\rangle \otimes\left|\chi_{B}\right\rangle$ where $\left|\chi_{j}\right\rangle$ is a pure state of the $j^{\text {th }}$ subsystem, then the state is separable, but if it can not be written in the showed form then the state is entangled. For instance, let's take a state

$$
|\psi\rangle=\frac{1}{\sqrt{2}}(|00\rangle+|11\rangle) .
$$

This state cannot be written as a product state because the state should have the following form

$$
(a|0\rangle+b|1\rangle) \otimes(A|0\rangle+B|1\rangle)
$$

where $a A$ and $b B$ should be equal to $\frac{1}{\sqrt{2}}, a B$ and $b A$ should be equal to 0 , which means that two of these coefficients must be equal to 0 but this cannot be done, for instance, by taking a product between $a B=0$ and $b A=0$ and a second product between $a A=\frac{1}{\sqrt{2}}$ and $b B=\frac{1}{\sqrt{2}}$, then we have:

$$
\left.\left.\begin{array}{l}
a B=0 \\
b A=0
\end{array}\right\} a b A B=0 ; \quad \begin{array}{l}
a A=\frac{1}{\sqrt{2}} \\
\\
b B=\frac{1}{\sqrt{2}}
\end{array}\right\} a b A B=\frac{1}{2}
$$

we find that these products are inconsistent. Then

$$
\frac{1}{\sqrt{2}}(|00\rangle+|11\rangle) \neq(a|0\rangle+b|1\rangle) \otimes(A|0\rangle+B|1\rangle)
$$

Therefore, $|\psi\rangle$ is an entangled quantum state.
Let's consider the density matrix $\rho$, which corresponds to a mixed state, is separable if there is $p_{k}$ is greater or equal to $0, \rho_{1}^{k}$ and $\rho_{2}^{k}$ are mixed states of the subsystems:

$$
\rho=\sum_{k} p_{k} \rho_{1}^{k} \otimes \rho_{2}^{k}
$$

where

$$
\sum_{k} p_{k}=1,
$$

The state is called simply separable or product state when the state can be expressed as $\rho=\rho_{1} \otimes \rho_{2}$. If the two conditions, $p_{k} \geq 0$ and $\rho=\sum_{k} p_{k} \rho_{1}^{k} \otimes \rho_{2}^{k}$, are not attained then $\rho$ is entangled mixed state. In the next chapter, we will define and use the notion of generalised separability.

## Local hidden variable model

The original motivation for investigating separability [63] was to understand local hidden variables $[26,8]$. Let's consider a source that creates particles $\lambda$. These particles are sent to two distant detectors $D_{1}$ and $D_{2}$ controlled by two people. By assuming that each person can choose his measurement, $M_{1}$ and $M_{2}$, the two detectors generate two outcomes $A_{1}$ and $A_{2}$ respectively. Hence, $\lambda$ holds local hidden variables that have a probability distribution $p_{\lambda}$.


Figure 1.2: Local hidden variables

Local hidden variable models can be presented by the following equation,

$$
\sum_{\lambda} p_{\lambda} p^{\lambda}\left(A_{1} \mid M_{1}\right) p^{\lambda}\left(A_{2} \mid M_{2}\right)
$$

where $p^{\lambda}\left(A_{i} \mid M_{i}\right)$ is the probability of getting an outcome $A_{i}$ given a measurement $M_{i}$. Bell [15] showed that local hidden variables are inconsistent with the statistical predictions of quantum mechanics. Experiments have been performed to demonstrate
non-locality $[24,21]$.

## Positive partial transpose

In general, it is difficult to test whether a state $\rho_{A B}$ is separable or not [25], in our case the so-called PPT test will be sufficient for this research. The positive partial transpose test or PPT test is one of the useful tests that works to find whether the state is separable [47]. For instance, if we have

$$
X_{A B}=\sum_{i, j, k, l} c_{i, j, k, l}|i\rangle\left\langle\left. j\right|_{A} \otimes \mid k\right\rangle\left\langle\left. l\right|_{B}\right.
$$

then the partial transpose of $X_{A B}$ is

$$
X_{A B}^{T_{A}}=\sum_{i, j, k, l} c_{i, j, k, l}|i\rangle\left\langle\left. j\right|_{A} ^{T} \otimes \mid k\right\rangle\left\langle\left. l\right|_{B}=\sum_{i, j, k, l} c_{i, j, k, l} \mid j\right\rangle\left\langle\left. i\right|_{A} \otimes \mid k\right\rangle\left\langle\left. l\right|_{B}\right.
$$

If we consider $\rho$ is a separable state

$$
\rho=\sum_{i} p_{i}\left|\psi_{i}\right\rangle\left\langle\left.\psi_{i}\right|_{A} \otimes \mid \phi_{i}\right\rangle\left\langle\left.\phi_{i}\right|_{B},\right.
$$

then

$$
\rho^{T_{A}}=\sum_{i} p_{i}\left|\psi_{i}^{*}\right\rangle\left\langle\left.\psi_{i}^{*}\right|_{A} \otimes \mid \phi_{i}\right\rangle\left\langle\left.\phi_{i}\right|_{B}\right.
$$

where $\left|\psi_{i}^{*}\right\rangle$ are pure states and their coefficients, that are in the computational basis, are the complex conjugates of $\left|\psi_{i}\right\rangle$. After finding the partial transpose, if $\rho^{T_{A}}$ is positive semi-definite then $\rho$ is PPT. Therefore, $\rho$ is separable which implies that $\rho \in$ PPT. If $\rho \notin P P T$ then $\rho$ is an entangled quantum state but if $\rho \in P P T$ then $\rho$ can be separable. However, $2 \times 2$ and $2 \times 3$ systems is known to be equivalent to PPT separability [34]. But in the case of $3 \times 3$ and $2 \times 4$ systems, the eigenvalues of the partial transpose of entangled mixed states are found to be positive [35].

### 1.1.5 Controlled gates

The manipulation of qubits can be done with the help of quantum logic gates. So quantum logic gates are basic quantum circuits that operate on a small number of qubits. Quantum logic gates are represented by unitary operators,

$$
U^{\dagger} U=U U^{\dagger}=I
$$

where $U^{\dagger}$ is the conjugate transpose matrix of $U$.
A gate acting on $n$ qubits is represented by a $2^{n} \times 2^{n}$ unitary matrix. There exists an uncountable infinite number of gates but the most used gates are the Identity and the Pauli gates. The Identity gate is the identity matrix, and it is defined for a single qubit:

$$
I=\left(\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right)
$$

The Pauli gates are the Pauli matrices and act on a single qubit. The Pauli "X", "Y" and " $Z$ " correspond to rotations around the $x, y$ and $z$ axes of the Bloch sphere.

1. The Pauli- $X$ gate is sometimes called the bit flip as it maps $|0\rangle$ to $|1\rangle$ and $|1\rangle$ to $|0\rangle$.

$$
X=\sigma_{x}=\left(\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right)
$$

2. The Pauli- $Y$ gate maps $|0\rangle$ to $i|1\rangle$ and $|1\rangle$ to $-i|0\rangle$.

$$
Y=\sigma_{y}=\left(\begin{array}{cc}
0 & -i \\
i & 0
\end{array}\right)
$$

3. The Pauli- $Z$ gate leaves the basis state $|0\rangle$ unchanged and maps $|1\rangle$ to $-|1\rangle$.

$$
Z=\sigma_{z}=\left(\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right)
$$

The square of a Pauli matrix is the identity matrix,

$$
I^{2}=X^{2}=Y^{2}=Z^{2}=-i X Y Z=I
$$

Also, they anti-commute, for example, $Z X=i Y=-X Z$.
There is also the Hadamard gate that maps the basis state $|0\rangle$ to $\frac{|0\rangle+|1\rangle}{\sqrt{2}}$ and $|1\rangle$ to $\frac{|0\rangle-|1\rangle}{\sqrt{2}}$,

$$
H=\frac{1}{\sqrt{2}}\left(\begin{array}{cc}
1 & 1 \\
1 & -1
\end{array}\right)
$$

Controlled gates are the gates that act on two or more qubits, where one or more qubits act as a control. In general, let's consider $U$ is a gate that acts on one single
qubit. The matrix of this gate is represented by

$$
U=\left(\begin{array}{cc}
U_{00} & U_{01}  \tag{1.6}\\
U_{10} & U_{11}
\end{array}\right)
$$

Then the controlled gate operates on two qubits where the first qubit acts as a control. So the matrix that represents $C U$ is

$$
C U=\left(\begin{array}{cccc}
1 & 0 & 0 & 0  \tag{1.7}\\
0 & 1 & 0 & 0 \\
0 & 0 & U_{00} & U_{01} \\
0 & 0 & U_{10} & U_{11}
\end{array}\right)
$$

For example, the controlled- $N O T$ gate or $C X$ acts on two qubits and performs the $X$ on the second qubit only when the first qubit is $|1\rangle$ and when the first qubit is $|0\rangle$ leaves it unchanged, which means

$$
\begin{aligned}
|00\rangle & \rightarrow|00\rangle \\
|01\rangle & \rightarrow|01\rangle \\
|10\rangle & \rightarrow|11\rangle \\
|11\rangle & \rightarrow|10\rangle
\end{aligned}
$$

It is represented by the following matrix:

$$
C N O T=\left(\begin{array}{llll}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 0 & 1 \\
0 & 0 & 1 & 0
\end{array}\right)
$$

The Controlled- $Z$ gate or $C Z$ acts on two qubits and performs the $Z$ operation on the second qubit only when the first qubit is $|1\rangle$ and when the first qubit is $|0\rangle$ leaves it unchanged, which means

$$
\begin{aligned}
|00\rangle & \rightarrow|00\rangle \\
|01\rangle & \rightarrow|01\rangle \\
|10\rangle & \rightarrow|10\rangle
\end{aligned}
$$

$$
|11\rangle \rightarrow-|11\rangle
$$

The matrix is given by:

$$
C Z=\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & -1
\end{array}\right)
$$

### 1.1.6 Affine and convex sets

In this section, we are going to define the convex hull that we are going to use in following chapters.

Let's take two different points $x_{1}$ and $x_{2}$ in $\mathbb{R}^{n}$ and these points are of the form

$$
y=\lambda x_{1}+(1-\lambda) x_{2}
$$

where $\lambda \in \mathbb{R}$. The previous equation forms the line passing through $x_{1}$ and $x_{2}$. If the parameter $\lambda=0$ then the equation $y=x_{2}$, and if $\lambda=1$, then $y=x_{1}$. So the values of $\lambda$ between 0 and 1 represent the closed line segment between the two points $x_{1}$ and $x_{2}$.

Now, let's consider a set $C \subseteq \mathbb{R}^{n}$. $C$ is affine if the line between any two different points in $C$ lies in $C$, which means that if for any $x_{1}, x_{2} \in C$ and $\lambda \in \mathbb{R}$, there is $\lambda x_{1}+(1-\lambda) x_{2} \in C . C$ has the linear combination of any two points in $C$, and all the coefficients in the linear combination sum to one. In general, by taking more than two points in the following form

$$
\lambda_{1} x_{1}+\lambda_{2} x_{2}+\ldots+\lambda_{k} x_{k}
$$

where

$$
\lambda_{1}+\lambda_{2}+\ldots+\lambda_{k}=1
$$

the equation is defined as an affine combination of the points $x_{1}, x_{2}, \ldots, x_{k}$.
Moreover, if $C$ is an affine set, $x_{1}, x_{2}, \ldots, x_{k} \in C$ and $\lambda_{1}+\lambda_{2}+\ldots+\lambda_{k}=1$, then the point

$$
\lambda_{1} x_{1}+\lambda_{2} x_{2}+\ldots+\lambda_{k} x_{k} \in C
$$

The affine hull (aff) is the set of all affine combinations of points in the set $C$.

$$
\text { aff } C=\left\{\lambda_{1} x_{1}+\lambda 2 x_{2}+\ldots+\lambda_{k} x_{k} \mid x_{1}, x_{2}, \ldots, x_{k} \in C, \lambda_{1}+\lambda_{2}+\ldots+\lambda_{k}=1\right\}
$$

By taking a line segment between two points in $C, C$ is convex set if this line segment lies in $C$, it means that for any $x_{1}, x_{2} \in C$, any $\lambda$ between 0 and 1 , it hold that

$$
\lambda x_{1}+(1-\lambda) x_{2} \in C
$$

Every affine set is a convex set because the set contains the whole line between two different points in the set.

A convex combination of the points $x_{1}, x_{2}, \ldots, x_{k}$ has a form

$$
\lambda_{1} x_{1}+\lambda_{2} x_{2}+\ldots+\lambda_{k} x_{k}
$$

where $\lambda_{1}+\lambda_{2}+\ldots+\lambda_{k}=1$ and $\lambda_{i} \geq 0$ for $i=1,2, \ldots, k$.
The set of all convex combinations of points in a set $C$ is called the convex hull (conv) of $C$,
$\operatorname{conv} C=\left\{\lambda_{1} x_{1}+\lambda_{2} x_{2}+\ldots+\lambda_{k} x_{k} \mid x_{i} \in C, \lambda_{i} \geq 0, i=1,2, \ldots, k, \lambda_{1}+\lambda_{2}+\ldots+\lambda_{k}=1\right\}$

### 1.1.7 Cluster States

In quantum information and quantum conputation, a cluster state is not a single quantum state, it is a family of quantum state that can support quantum computation. It can be explained as follows: any graph that contains a number of vertices, $n$, we can set $n$-qubit cluster state, it means that each vertex corresponds to a qubit in the state $\frac{|0\rangle+|1\rangle}{\sqrt{2}}$. Each edge that represents the $C Z$ gate, connect two vertices or two qubits. For two dimensions, the cluster state is considered as rectangular lattice and for $n$ dimensions, it corresponds to a graph with $n$-dimensional lattice. Initially, the "cluster state" [53, 54, 45] is a type of quantum computation and it is made of one qubit measurements.

Cluster states can be defined as graph states which is a particular sort of multiqubit state that can be shown as a graph.

In this thesis, we are going to consider what happens if we change the inputs and


Figure 1.3: 9-qubit cluster state
the gates of the cluster state quantum computation and we will show that there are regions that can be efficiently simulated classically.

### 1.2 Classical simulation

In general, it is not believed that quantum system can be efficiently simulated classically [57]. But for several years great effort has been devoted to the study of efficient simulation on a classical computer. These studies have offered rich contributions to the field of research. In this section, we will review a non-exhaustive selection of quantum systems that have been shown to be classical simulatable.

## 1. Stabilizer circuit

A stabilizer formalism can be described by an example. Let's consider a state of two qubits,

$$
|\psi\rangle=\frac{|00\rangle+|11\rangle}{\sqrt{2}} .
$$

By applying two operators $X \otimes X$ and $Z \otimes Z$ on $|\psi\rangle$, we find that the state remains the same,

$$
(X \otimes X)|\psi\rangle=|\psi\rangle \text { and }(Z \otimes Z)|\psi\rangle=|\psi\rangle
$$

then the state $|\psi\rangle$ is stabilized by these two operators [46]. The stabilizer formalism has an important concept which states that a quantum state can be presented by using Pauli operators that stabilize it. By applying unitary gates on the quantum state, it means that these gates act on the string of the Pauli operators. But if these unitary gates are made of $C N O T$, Hadamard and phase gates, these Pauli operators are transformed into new set of stabilizer Pauli operators. Then a circuit that consists only of these gates is called a stabilizer circuit [1].

Gottesman-Knill theorem shows that a quantum computation, that consists of state preparations in the computational basis, gates in the so-called Clifford group [27], Pauli gates and measurements in the Pauli basis, may be efficiently simulated on a classical computer.
This theorem has been extended in a variety of ways [64, 13, 14, 40]. It also has been extended to quantum circuits that are composed of Clifford and a single type of non-Clifford gates [19]. For instance, in [19], the classical algorithm that simulate the system is polynomial in the number of qubits and the number of Clifford gates but it is exponential in the limited number of the non-Clifford gates.

## 2. Matrix product states and tensor networks

The product state $|\psi\rangle_{1 \ldots n}$ can be written in the following form,

$$
\begin{equation*}
|\psi\rangle_{1 \ldots n}=|\alpha\rangle_{1}|\beta\rangle_{2} \ldots|\kappa\rangle_{n} \tag{1.8}
\end{equation*}
$$

where this state consists of $n$ qubits. A matrix product state or $M P S$ is a one dimension quantum state of $n$ qubits and it is a generalisation of the above equation. Each state in equation (1.8) is replaced by a matrix of state; for instance, the state $|\alpha\rangle$ is changed into $\left|\alpha_{i j}\right\rangle$. Therefore, the matrix product of $n$ matrices is given by the following form:

$$
\begin{equation*}
|\psi\rangle=\sum_{i, j, k, \ldots, \ldots}\left|\alpha_{i j}\right\rangle\left|\beta_{j k}\right\rangle \ldots\left|\kappa_{m i}\right\rangle \tag{1.9}
\end{equation*}
$$

If the matrix of states, $\left|\alpha_{i j}\right\rangle, \ldots,\left|\kappa_{m i}\right\rangle$, are written as $A, B \ldots, K$, then the state $|\psi\rangle=\operatorname{tr}(A B \ldots K)$. In the standard basis, states can be written as,

$$
\left|\alpha_{i j}\right\rangle=\sum_{i_{1}} A_{i j}^{\left(i_{1}\right)}\left|i_{1}\right\rangle, \ldots,\left|\kappa_{m i}\right\rangle=\sum_{i_{n}} K_{m i}^{\left(i_{n}\right)}\left|i_{n}\right\rangle,
$$

then the matrix product state $|\psi\rangle$ will be given by

$$
\begin{equation*}
|\psi\rangle=\sum_{i, \ldots m, i_{1} \ldots i_{n}} A_{i j}^{\left(i_{1}\right)} B_{j k}^{\left(i_{2}\right)} \ldots K_{m i}^{\left(i_{n}\right)}\left|i_{1}\right\rangle\left|i_{2}\right\rangle \ldots\left|i_{n}\right\rangle . \tag{1.10}
\end{equation*}
$$

As mentioned above, the MPS is a one dimension state [48] but there is another description of MPS, which is called Projected entangled pairs state or PEPS. $P E P S$ extends the $M P S$ in two and higher dimensions.

Using MPS, Jozsa [37] considered any poly sized qauntum circuit of $n$ qubits and has shown that this circuit can be efficiently simulated classically only if it is logarithmic in the number of qubit gates. Using the matrix product state representation, a classical methodology is presented to simulate cluster state quantum computation that has a polynomial number of qubits and an exponential number in the width of the cluster [66]; then it is shown that any logarithmic depth quantum computation with gates that relates only near qubits can be efficiently simulated classically.

The tensor network which is a generalisation of matrix product states, can be described by a graphical notation [16]. The matrix product state $|\psi\rangle$ which is presented by equation (1.10) can be illustrated by the following graph:


Figure 1.4: This image presents the tensor network. The MPS can also be written as: $\sum_{i_{1}, \ldots, i_{n}} c_{i_{1}, \ldots, i_{n}}\left|i_{1}, \ldots, i_{n}\right\rangle$.

Markov and Shi [41] proved that a quantum circuit with a non-Clifford gates whose graph is close to a tree graph as measured by the treewidth, can be simulated in an exponential time. In addition, any quantum circuit can be efficiently simulated using the tree tensor network contraction can only produce a limited amount of entanglement[65].

## 3. Matchgates

A matchgate can be defined as a gate $G$ of two qubits that takes the following form in the computational basis,

$$
G=\left(\begin{array}{cccc}
m & 0 & 0 & n \\
0 & r & s & 0 \\
0 & t & w & 0 \\
p & 0 & 0 & q
\end{array}\right)
$$

where

$$
A=\left(\begin{array}{cc}
m & n \\
p & q
\end{array}\right) \text { and } B=\left(\begin{array}{cc}
r & s \\
t & w
\end{array}\right)
$$

are the inner and the outer block of $G$ respectively and they are in the special unitary group or $S U(2)$ which is a $2 \times 2$ unitary matrix with a determinant equal to 1 . Valiant $[60,59]$ defines the notion of matchgate and demonstrates that matchgate circuits can be classically simulated in polynomial time under specific conditions. Then Terhal and Divincenzo [58] analyze the class of quantum computations demonstrated by Valiant [60] and show that this class is related to a fermions model in one dimension. Matchgate circuits in which these gates are only applied to neighbour qubit lines can also be efficiently simulated classically [39]. Instead of taking a circuit that has product input and the output consists
of a single qubit measurement [39], Brod [20] displays that even if the input is product states and the output has measurements of many qubits, matchgates can be simulated on a classical computer.

### 1.3 Quantum supremacy

The term quantum supremacy describes the goal of building controllable quantum devices with well understood dynamics that cannot be efficiently simulated on classical computers [32, 17]. The Martinis-Google group claimed a demonstration of quantum supremacy by taking a random quantum circuit and seeing whether the outcome of measurements can be sampled on a classical computer [6]. However, there have been number of attempts to replicate their work on classical computers [9, 2], and so the debate is far from being conclusive.

### 1.4 Preview of the thesis results

In this thesis, we will begin by changing the inputs of cluster state quantum computation to end up with some inputs that allow the system to be efficiently simulated classically. If different pure state choices for $|\psi\rangle=|+\rangle=\frac{|0\rangle+|1\rangle}{\sqrt{2}}$ are made, two facts are well known already:

1. If $|\psi\rangle=|0\rangle$ or $|\psi\rangle=|1\rangle$ then the system is classically simulatable efficiently.
2. If $|\psi\rangle$ is a state from the $X Y$ plane of the Bloch sphere then the system has the power of quantum computation.

We explore questions related to this in the various chapters of the thesis:

- In chapter 2, we will show that if $|\psi\rangle$ has the (unormalised) form $|0\rangle+\epsilon|1\rangle$ where $\epsilon$ is a small non-zero number, the system can be efficiently simulated classically. We are going to use new state spaces that involve cylinders rather than Bloch sphere and these cylinders are parameterized by radius $r$. We must grow the radius in order to get separable decomposition. However, the radius cannot grow too much. Then instead of $C Z$ gates, we will try to apply diagonal gates $V(\theta)=|0\rangle\langle 0| \otimes I+|1\rangle\langle 1| \otimes Z_{\theta}$ on two cylindrical state spaces and we will find that almost all inputs can be simulated classically when $\theta$ is small. In addition, it turns out that $C Z$ gate requires the most growth factor $\lambda$, which is
the ratio of the output radius to the input radius, to maintain separability. My contribution was the computation of the growth factors for the $V(\theta)$ gates.
- In chapter 3, by taking some examples of alternative local state spaces and by growing the size of these specific states, we will find that the setup can also be efficiently simulated on a classical computer. Then, we tried optimize our classical simulation over these state spaces that we tried and it turns out that the cylinder requires the least growth. My contribution was to perform all these calculations using linear programming.
- In chapter 4 , we attempt to improve the radius of a system that can be efficiently simulated classically. By using the coarse graining, we will find that the value of the radius of four particles in a block is greater than the radius of two particles. We will discover that by adding more vertices in a block, the value of the radius, that can be efficiently simulated classically, increases slightly [7]. My contribution was to calculate the radius threshold for four particles, assuming a conjecture about which measurements determine when certain invalid negative probabilities arise.


## Chapter 2

## Efficient Classical Simulation of quantum circuits with alternative inputs

A typical quantum computer is a sequence of unitary quantum gates acting on qubits. Quantum gates are capable of generating entanglement between these qubits. But real quantum gates suffer from noise which can destroy the creation of entanglement.

This means that very noisy systems cannot be used for better than classical computation, because it is known that gate model quantum computers with gates that don't generate entanglement can be efficiently simulated classically [33] (this is a part of variety of works that investigate when classical simulation is possible if there is a limited number of entanglement $[38,62])$. Here, by efficiently classical simulation, we mean sampling in polynomial time from a probability distribution that approximates the probability distribution of measurement outcomes to arbitrary accuracy.

When we say that a gate doesn't generate entanglement, we mean that a quantum gate $\mathcal{E}$ takes a product input, $\rho \otimes \sigma$, to a separable output state. We then say that $\mathcal{E}$ is a separable gate. So, it preserves the separability.

One of the motivations of the above classical algorithm of [33] was to obtain upper bounds on the value of the fault-tolerance threshold, which is the noise level below which quantum computation can be achieved, because once the noise level is high enough to make the gate separable the device can be efficiently simulated classically. If we assume that quantum computers cannot be efficiently simulated classically, the fault-tolerance threshold must be lower than this.

This algorithm cannot be used to simulate pure entangled quantum systems. However, we will see later in this chapter that by modifying the notion of separability, the algorithm in Harrow and Nielsen [33] can be used to simulate some entangled quantum systems. Specifically, we will consider particular entangled quantum systems that are variant of cluster state quantum computation. As discussed in chapter 1, in cluster state, we initialize in the state $|\psi\rangle=|+\rangle=\frac{|0\rangle+\epsilon|1\rangle}{\sqrt{2}}$; but we are going
to show that if we initialize in the state $|\psi\rangle=|0\rangle+\epsilon|1\rangle$ where $\epsilon$ is a small enough non-zero number, instead of the $|+\rangle$ state, the system can be efficiently simulated on a classical computer. It means that by providing a small enough $\epsilon$, it turns out that we can give a classical simulation algorithm that uses cylindrical separability rather than normal separability. Then we will replace the $C Z$ gates by other diagonal gates $V(\theta)=|0\rangle\langle 0| \otimes I+|1\rangle\langle 1| \otimes Z_{\theta}$ where $Z_{\theta}=|0\rangle\langle 0|+e^{i \theta}|1\rangle\langle 1|$ and we will find similar results that hold for these diagonal gates.

First, we are going to describe the algorithm of Harrow and Nielsen [33] which we use in the research. Then we are going to modify the pure states of the cluster state quantum computation to end up with a system that can be efficiently simulated classically using the new notion of generalised separability. Finally, we will apply diagonal gates instead of $C Z$ gates on the two inputs and it turns out that almost all inputs can be simulated classically.

### 2.1 Description of the algorithm in Harrow and Nielsen[33]

The algorithm of [33] starts with an assumption that the input is a product state, $\rho_{A} \otimes$ $\rho_{B} \otimes \rho_{C} \ldots$, acted upon by separable quantum gates, $\mathcal{E}_{1}, \ldots, \mathcal{E}_{p(n)}$ where the polynomial $p(n)$ is the number of gates and $n$ is the measure of problem size. For instance, if the first gate, $\mathcal{E}_{1}$, which is assumed to be separable, acts on two qubits $A$ and $B$,

$$
\mathcal{E}_{1}\left(\rho_{A} \otimes \rho_{B}\right)=\sum_{k} p_{k} \rho_{A}^{k} \otimes \rho_{B}^{k}
$$

The input is taken to be a separable output state. Then the algorithm samples the probability distribution, $p_{k}$, gets one outcome and stores the product state which corresponds to the obtained outcome, $\rho_{A}^{k} \otimes \rho_{B}^{k}$. This product state is considered as an input for the next separable gate, for instance, $\mathcal{E}_{2}$. These gates may act not only on the stored product state but also on different qubits. This process is repeated until it reaches the final outcome which is produced from the final separable gate $\mathcal{E}_{p(n)}$.

At the end of the process, qubits are measured at the output of the quantum computation, where the output is a string of product states. Using the trace, the probability distribution of each one of the products is computed. Because these computations involve a linear number of $2 \times 2$ matrices, they can be computed classically. The algorithm simulates the final measurement when it samples from those
probability distributions and produces a classical outcome for measurements of each qubit. It ends up with an equivalence between the classical simulation and certain measurement on quantum states.

Now let's describe in detail how the classical simulation is performed in the paper of [33]. The initial state of the computer is assumed to be in a computational basis state. The algorithm uses variables in the classical simulation, that are three dimensional real vectors $\vec{s}_{j}$ for each $j=1, \ldots, q(n)$ where the polynomial $q(n)$ is the number of qubits. These vectors represent the bloch vectors of qubits in the quantum computers which the algorithm simulates. Each $\vec{s}_{j}$ is valid which means that it satisfies many properties. Two of these properties are: each component of the vector is in $[-1,1]$ and the norm of the vector is less than or equal to $1 . \rho(\vec{s})$ is defined to be the density matrix of $q(n)$ qubits when $\vec{s}$ is valid.
W.l.o.g. a single qubit gate can be considered as a two qubit product gate. Then the algorithm assumes that $\mathcal{E}_{1}, \ldots, \mathcal{E}_{p(n)}$ are two-qubit separable gates. It starts to simulate the first separate gate $\mathcal{E}_{1}$ that acts on two qubits $A$ and $B$, using the input $\vec{s}$ in the gate simulation procedure.

The gate simulation procedure is described as followed. The input was represented by the three dimensional vector $\vec{s}$ as mentioned above. The body of the procedure consisted of finding valid vectors, ${\overrightarrow{s_{A}}}^{j}$ and ${\overrightarrow{s_{B}}}^{j}$ on the output. The equation can be written as:

$$
\mathcal{E}_{1}\left(\rho\left(\overrightarrow{s_{A}}\right) \otimes \rho\left(\overrightarrow{s_{B}}\right)\right)=\sum_{k} p_{k} \rho\left({\overrightarrow{s_{A}}}^{j}\right) \otimes \rho\left({\overrightarrow{s_{B}}}^{j}\right)
$$

These separable decompositions are computed to a certain accuracy, which is introduced to be less than or equal to $2^{-l}$ where $l$ is the number of bits of precision. To find the probabilities and the vectors, it requires operations that are poly $\left(2^{l}\right)$ which turns out to be sufficient (as explained in [33] in section 3.B) and $l$ is picked to be logarithmic in the circuit; because if the circuit is large, more gates, that introduce more errors, are required; so $l$ needs to be large.

By Caratheodory's theorem, there are at most 16 terms in the sum; we mean that if a point of $\mathbb{R}^{d}$ is in a convex hull of a set, this point can be written as the convex combination of at most $d+1$ in the set. Then the algorithm samples the probability distribution, $p_{k}$, and stores the obtained product state, $\rho\left({\overrightarrow{s_{A}}}^{j}\right) \otimes \rho\left({\overrightarrow{s_{B}}}^{j}\right)$, and it is used as an input for the following gate. This process is repeated many times until the gate $\mathcal{E}_{p(n)}$ is reached and a final output is produced. At the end, this procedure produced a valid vector as an output. And this vector is the set of all three dimensional real vectors of $n$ qubits.

Finally, the final measurement is simulated in the computational basis. The algorithm measures all qubits at the output of the quantum computation. The probability distribution is computed of each one of the products. The algorithm samples from those probability distributions and produces a classical outcome for measurements of each qubit.

As a conclusion, [33] prove that a quantum device built from noisy quantum gates, that don't generate entanglement, can be efficiently simulated classically.

### 2.2 Cylinder separability

Explicitly, as we mentioned in the introduction, we will be studying a variant of cluster state quantum computation and we will show that it could be efficiently simulated on a classical computer. The algorithm of [33] cannot be used because these cluster state circuits lead to pure entangled quantum states; and these quantum states cannot be approximated by a separable decomposition. But we develop a new method; we are going to use a different notion of separability to show that these states can be efficiently simulated classically.

Separable states are quantum states that can be decomposed into individual states which belong to separate particles [63]. In the multipartite case, a pure state $|\phi\rangle$ is separable if it has the form: $|\phi\rangle=\left|\phi_{1}\right\rangle \otimes \ldots \otimes\left|\phi_{n}\right\rangle$.

And a mixed state $\rho$ is separable if it has the following form

$$
\begin{equation*}
\rho=\sum_{i} p_{i} \rho_{A}^{i} \otimes \rho_{B}^{i} \otimes \rho_{C}^{i} \otimes \ldots \tag{2.1}
\end{equation*}
$$

where $\rho_{A}^{i}, \rho_{B}^{i}, \rho_{C}^{i}, \ldots$, are positive density matrices. In normal quantum state, these density matrices contain bloch vectors that come from the bloch sphere.

In generalised separability [52, 51],

$$
\rho=\sum_{j} p_{j} \sigma_{A}^{j} \otimes \sigma_{B}^{j} \otimes \sigma_{C}^{j} \otimes \ldots
$$

where $\sigma_{A}^{j}, \sigma_{B}^{j}, \sigma_{C}^{j}, \ldots$, are positive density matrices that are drawn from a set $Q$ of quantum states. We relax this condition by allowing this set to be a new state space. For instance, instead of the Bloch sphere, we will consider Bloch vectors drawn from a cylindrical state space with a specific radius $r \leq 1$ and a height $h=1$. This is the new notion of generalised separability, the so-called "cylinder separability", which is
going to lead to efficient classical simulation. Generalised separability has also been previously used to construct efficient simulation on a classical computer of certain type of PEPS [3] where these states are in general, hard to simulate[55].

As we will see later in this chapter, by taking two of these cylindrical state spaces acting with the controlled- $Z$ gate, we will end up with a separable output state with respect to two cylindrical state spaces with larger radii " $R$ ". Then the system could be efficiently simulated on a classical computer only when the output radii are large enough for the input radii, which means that the output radii should be less than or equal to 1 because of the negativity that always appears [5, 4].

Note that, even though we have systems that are close to product state, it does not mean that the setup can be efficiently simulated classically [61, 29, 28]. It has also been shown that variant of the systems that we are considering can enable quantum computation under the assumption that repeated measurements on any given qubit are allowed $[43,12]$, however no previous work has shown the efficient classical simulation results that we present. In addition, our framework does not fit into Generalized probabilistic theories or GPTs because if $C Z$ gates are applied many times, the radius will grow and we will end up with negative probabilities. By GPT, we mean that it is a framework for generalizing quantum theory to explore the cause of its unusual features [49, 36, 31, 10, 11].

### 2.3 Cylinders and the restricted measurements

For a specific set of measurement operators, denoted by $\{M\}$, the dual is a set of operators $\rho$, that give positive numbers under the Born rule, $\operatorname{tr}(M \rho)$. For instance, motivated by the cluster state quantum computation schemes, if we measure in the $X-Y$ plane and $Z$ directions, the dual is a cylinder. To demonstrate how a cylinder is obtained, we start by measuring in the $Z$ direction, we need to compute the trace, that should be positive,

$$
\operatorname{tr}\left(\frac{I \pm Z}{2} \rho\right)=\operatorname{tr}\left\{\left(\frac{I \pm Z}{2}\right)\left(\frac{I+x X+y Y+z Z}{2}\right)\right\}=\frac{1 \pm z}{2}
$$

By adding a constraint that the operators are unit trace, we end up with two parallel horizontal planes that should be orthogonal to the $Z$ axis where $Z$ is in the range $[-1,1]$, and all the operators of the bloch space between these planes give
positive probabilities. Then we add the $X-Y$ measurement to the $Z$ measurement,

$$
\operatorname{tr}\left(\frac{I \pm X}{2} \rho\right)=\operatorname{tr}\left\{\left(\frac{I \pm X}{2}\right)\left(\frac{I+x X+y Y+z Z}{2}\right)\right\}=\frac{1 \pm x}{2}
$$

and

$$
\operatorname{tr}\left(\frac{I \pm Y}{2} \rho\right)=\operatorname{tr}\left\{\left(\frac{I \pm Y}{2}\right)\left(\frac{I+x X+y Y+z Z}{2}\right)\right\}=\frac{1 \pm y}{2}
$$

$X$ and $Y$ are in the range $[-1,1]$. Therefore, the intersection of all these planes makes a cylindrical state space. We are using this specific non-physical input state space because we are only considering the restricted measurements, $Z$ and $X Y$ measurements. So, we are interested in quantum and non-quantum operators that give valid and positive probabilities for these measurements under the born rule.

Thus, the same process of [33] algorithm is applied with this new type of state space, that allows us to express entangled pure states as separable states. Instead of the bloch sphere, we can use the cylinder because of the measurement restrictions, and those cylinder states still give us a valid probability distribution.

### 2.4 Bloch vectors of the cylinder

In this section, we discuss two ways in which we characterize the cylinder.
The Bloch vectors of the cylinder can be characterized by

$$
\begin{equation*}
\left\{(x, y, z) \mid z \in[-1,1], x^{2}+y^{2} \leq 1\right\} \tag{2.2}
\end{equation*}
$$

Also, the cylinder can be defined as:

$$
\begin{equation*}
C y l(r):=\left\{\rho \mid \rho=\rho^{\dagger}, \operatorname{tr} \rho=1,\left\|\rho-D_{Z}(\rho)\right\| \leq r\right\} \tag{2.3}
\end{equation*}
$$

where $D_{Z}(\rho):=\frac{1}{2}\left(\rho+Z \rho Z^{\dagger}\right)$ is the dephasing of $\rho$.

Equation (2.3) is a description of a cylinder with radius $r$; it doesn't consist of bloch vectors, it's about density matrices, that are a $2 \times 2$ matrices. Hence, equation (2.3) is equivalent to the equation (2.2), which is a description in terms of bloch vector $X, Y$ and $Z$ coordinates.

If we compute $\rho-D_{Z}(\rho)$ where

$$
\rho=\left(\begin{array}{ll}
\rho_{00} & \rho_{01} \\
\rho_{10} & \rho_{11}
\end{array}\right)
$$

so

$$
\begin{gathered}
\rho-D_{Z}(\rho)=\left(\begin{array}{ll}
\rho_{00} & \rho_{01} \\
\rho_{10} & \rho_{11}
\end{array}\right)-\frac{1}{2}\left\{\left(\begin{array}{ll}
\rho_{00} & \rho_{01} \\
\rho_{10} & \rho_{11}
\end{array}\right)+\left(\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right)\left(\begin{array}{ll}
\rho_{00} & \rho_{01} \\
\rho_{10} & \rho_{11}
\end{array}\right)\left(\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right)\right\} \\
=\left(\begin{array}{ll}
\rho_{00} & \rho_{01} \\
\rho_{10} & \rho_{11}
\end{array}\right)-\frac{1}{2}\left\{\left(\begin{array}{ll}
\rho_{00} & \rho_{01} \\
\rho_{10} & \rho_{11}
\end{array}\right)+\left(\begin{array}{cc}
\rho_{00} & -\rho_{01} \\
-\rho_{10} & \rho_{11}
\end{array}\right)\right\} \\
\rho-D_{Z}(\rho)=\left(\begin{array}{cc}
0 & \rho_{01} \\
\rho_{10} & 0
\end{array}\right)
\end{gathered}
$$

By taking $A=\rho-D_{Z}(\rho)$, we have to compute the determinant of $A-\lambda I$ to find the eigenvalues:

$$
A-\lambda I=\left(\begin{array}{cc}
0 & \rho_{01} \\
\rho_{10} & 0
\end{array}\right)-\lambda\left(\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right)=\left(\begin{array}{cc}
-\lambda & \rho_{01} \\
\rho_{10} & -\lambda
\end{array}\right)
$$

Determinant $(A-\lambda I)=\lambda^{2}-\rho_{01} \rho_{10}$, then

$$
\begin{gathered}
\lambda^{2}-\rho_{01} \rho_{10}=0 \\
\lambda^{2}=\rho_{01} \rho_{10} \\
\lambda= \pm \sqrt{\left|\rho_{01}\right|\left|\rho_{10}\right|}
\end{gathered}
$$

where, $\rho_{01}=\frac{x+i y}{2}, \rho_{10}=\frac{x-i y}{2}$ and $\left|\rho_{01}\right|=\left|\rho_{10}\right|$ because they are the complex conjugate of each other, so we have $\left|\rho_{01}\right|^{2}=\left|\rho_{10}\right|^{2}=\frac{x^{2}+y^{2}}{4}$,

$$
\lambda= \pm \sqrt{\left|\rho_{01}\right|^{2}}= \pm \sqrt{\frac{x^{2}+y^{2}}{4}}= \pm \frac{1}{2} \sqrt{x^{2}+y^{2}}
$$

so,

$$
\lambda_{1}=\frac{1}{2} \sqrt{x^{2}+y^{2}} ; \lambda_{2}=-\frac{1}{2} \sqrt{x^{2}+y^{2}}
$$

The trace norm is the sum of the absolute values of these eigenvalues where $|\lambda|=$

$$
\begin{aligned}
& \left|\lambda_{1}\right|=\left|\lambda_{2}\right|=\frac{1}{2} \sqrt{x^{2}+y^{2}} \\
& \quad \text { Trace } \operatorname{norm}(A-\lambda I)=\left|\lambda_{1}\right|+\left|\lambda_{2}\right|=2|\lambda|=2 \times \frac{1}{2} \times \sqrt{x^{2}+y^{2}}=\sqrt{x^{2}+y^{2}}
\end{aligned}
$$

By replacing $r$ in equation (2.3) with 1 , so we end up with equation (2.2).

### 2.5 Cylinder state space with different radii

To develop efficient classical simulation, we have to demonstrate how the generalised separable decomposition can be obtained in terms of cylinder state spaces.

### 2.5.1 $C Z$ gate acting on two input qubits

The important idea in this section is to show that if there are two inputs represented by two cylindrical state spaces with radii $r$ to which the $C Z$ gate has been applied, the output is separable with respect to two different cylindrical state spaces with larger radii $R$. It can be illustrated as follow:


Figure 2.1: A $C Z$ gate acts on two cylindrical state spaces with two radii $r_{A}$ and $r_{B}$, the output is separable with respect to two cylindrical state spaces with larger radii $R_{A}$ and $R_{B}$

Also, it can be demonstrated by the following Lemma:
Lemma 1. Consider a $C Z$ gate acting on two cylindrical state spaces with two different radii. Let's us state this more precisely: let's take $C Z\left(c y l\left(r_{A}\right) \otimes \operatorname{cyl}\left(r_{B}\right)\right)$ and any operator, as being $\operatorname{cyl}\left(r_{A}\right)$ and $\operatorname{cyl}\left(r_{B}\right)$ separable, can be written in the generalised separable form if and only if (where we define $f_{A}:=\frac{r_{A}}{R_{A}}$ and $f_{B}:=\frac{r_{B}}{R_{B}}$ )

$$
\begin{equation*}
\left(f_{A}+f_{B}\right)^{2}+f_{A}^{2} f_{B}^{2} \leq 1 \tag{2.4}
\end{equation*}
$$

Before starting with the proof, a question can be asked: when the $C Z$ can be considered as a separable operation?
Let's define the growth factor as:

$$
\begin{equation*}
g_{i}:=\frac{R_{i}}{r_{i}} \tag{2.5}
\end{equation*}
$$

As we mentioned previously, if $C Z$ gate acting on two cylindrical state spaces with different radii $r$, then the output that consists of two new cylindrical state spaces with larger radii $R$ is separable. It means that $g_{i}$ is large enough.

Let's consider the growth factors of the two states are equivalent. By applying the ratio to the determinant equation of the outer block we get:

$$
1-\frac{4}{g^{2}}-\frac{1}{g^{4}} \geq 0
$$

The value of $g$ is

$$
g \geq \lambda:=\sqrt{2+\sqrt{5}} \approx 2.05817
$$

so,

$$
R \geq 2.05817 r
$$

This means that the radii of the output should be minimum twice the radii of the input, then the $C Z$ can be considered a separable operation.

Proof. Let's start by writing the Pauli decomposition of two qubit operators that is represented by:

$$
\begin{equation*}
\rho_{A B}=\frac{1}{4} \sum_{i, j} \rho_{i, j} \sigma_{i} \otimes \sigma_{j} \tag{2.6}
\end{equation*}
$$

where $\sigma_{0}, \sigma_{1}, \sigma_{2}$ and $\sigma_{3}$ are the $2 \times 2$ matrices representing $I, X, Y$ and $Z$ operators respectively. The above equation can also be represented as a $4 \times 4$ matrix with rows and columns numbered from $0, \ldots, 3$ :

$$
\left(\begin{array}{llll}
\rho_{00} & \rho_{01} & \rho_{02} & \rho_{03} \\
\rho_{10} & \rho_{11} & \rho_{12} & \rho_{13} \\
\rho_{20} & \rho_{21} & \rho_{22} & \rho_{23} \\
\rho_{30} & \rho_{31} & \rho_{32} & \rho_{33}
\end{array}\right)
$$

By applying the $C Z$ gate on the matrix we get:

$$
\left(\begin{array}{cccc}
\rho_{00} & \rho_{01} & \rho_{02} & \rho_{03} \\
\rho_{10} & \rho_{11} & \rho_{12} & \rho_{13} \\
\rho_{20} & \rho_{21} & \rho_{22} & \rho_{23} \\
\rho_{30} & \rho_{31} & \rho_{32} & \rho_{33}
\end{array}\right) \rightarrow\left(\begin{array}{cccc}
\rho_{00} & \rho_{31} & \rho_{32} & \rho_{03} \\
\rho_{13} & \rho_{22} & -\rho_{21} & \rho_{10} \\
\rho_{23} & -\rho_{12} & \rho_{11} & \rho_{20} \\
\rho_{30} & \rho_{01} & \rho_{02} & \rho_{33}
\end{array}\right)
$$

The product operator can be written as:

$$
\frac{1}{2}\left(\sigma_{0}+x_{A} \sigma_{1}+y_{A} \sigma_{2}+z_{A} \sigma_{3}\right) \otimes \frac{1}{2}\left(\sigma_{0}+x_{B} \sigma_{1}+y_{B} \sigma_{2}+z_{B} \sigma_{3}\right),
$$

for conciseness, we will represent it by the following vector notation

$$
\left(1, x_{A}, y_{A}, z_{A}\right) \otimes\left(1, x_{B}, y_{B}, z_{B}\right)
$$

To determine if the output of a $C Z$ gate lead to cylindrical separable states that have different radii, we need to consider only the extremal input points from the top and the bottom faces of the cylinder. The reason for this consideration is that $C Z$ gate is a linear transformation which is a mapping between two vector spaces that preserves the operations of vector addition; so a linear transformation commutes with addition. Furthermore, when $C Z$ is applied on the the extremal input, the output is separable; then if a mixture of the extremal points is taken and $C Z$ has been applied, the output is also separable because $C Z$ is linear into a mixture of the original parts and those are separable. So the output for all the inputs will be separable.

We will take advantage of the symmetry around the $Z$ axis. Let's consider a
cylinder separable decomposition:

$$
C Z\left(\rho_{A} \otimes \rho_{B}\right)=\sum_{i} p_{i} \omega_{A}^{i} \otimes \omega_{B}^{i}
$$

where $\omega_{A}^{i} \in \operatorname{Cyl}\left(R_{A}\right)$ and $\omega_{B}^{i} \in \operatorname{Cyl}\left(R_{B}\right)$. We know that $C Z$ gate commutes with $Z$ rotation $U_{z}$ and cylinders are invariant under $Z$ rotations, because of these properties, we have the following separable decomposition:

$$
C Z\left(U_{z}^{A}\left(\rho_{A}\right) \otimes U_{z}^{B}\left(\rho_{B}\right)\right)=\sum_{i} p_{i} U_{z}^{A}\left(\omega_{A}\right) \otimes U_{z}^{B}\left(\omega_{B}\right)
$$

This is used to reduce the number of inputs that we need to test. First, we set the Pauli operator $y=0$ because if we take two inputs $\left(1, r_{A}, 0, \pm 1\right)$ and $\left(1, r_{B}, 0, \pm 1\right)$ and we apply the $Z$ rotation, we notice that these inputs are similar to other inputs, such as $\left(1, r_{A} \cos \phi X, r_{A} \sin \phi Y, \pm 1\right)$ and $\left(1, r_{B} \cos \phi X, r_{B} \sin \phi Y, \pm 1\right)$ where $\phi \neq 0$. Then, if we obtain a separable decomposition for the first two inputs, we will end up with separable decompositions for the other inputs by applying the $Z$ rotation.

So now, we can concentrate on the two inputs that have the form $\left(1, r_{A}, 0, \pm 1\right)$ and $\left(1, r_{B}, 0, \pm 1\right)$.

Let's see which value of $z$ we need to take for the both inputs particles. If we take $z=1$ for particle $A$, we will have:

$$
C Z\left(1, r_{A}, 0,1\right) \otimes\left(1, r_{B}, 0, \pm 1\right)=\sum_{i} p_{i} \omega_{A}^{i} \otimes \omega_{B}^{i}
$$

But if we take $z=-1$ for particle $A$, the separable decomposition will be

$$
C Z\left(1, r_{A}, 0,-1\right) \otimes\left(1, r_{B}, 0, \pm 1\right)=\sum_{i} p_{i} X \omega_{A}^{i} X^{\dagger} \otimes Z \omega_{B}^{i} Z^{\dagger}
$$

So, we need to take $z=1$ because changing $z$ to -1 when $y=0$, is the same as doing a $Z$ rotation on the second particle. And because $C Z$ is symmetric, the value of $z$ in the second input will be equal to 1 as well.

The two input state particles are given by:

$$
\begin{equation*}
\frac{1}{2}\left(I+r_{A} X+0 Y+1\right) \otimes \frac{1}{2}\left(I+r_{B} X+0 Y+1\right) \tag{2.7}
\end{equation*}
$$

After applying $C Z$ on the input, the matrix will be represented by:

$$
\left(\begin{array}{cccc}
1 & r_{B} & 0 & 1  \tag{2.8}\\
r_{A} & 0 & 0 & r_{A} \\
0 & 0 & r_{A} r_{B} & 0 \\
1 & r_{B} & 0 & 1
\end{array}\right)
$$

We need to show that the above matrix can be written as the outer product:

$$
\sum_{i} p_{i}\left(\begin{array}{c}
1 \\
R_{A} \cos \left(\theta_{i}\right) \\
R_{A} \sin \left(\theta_{i}\right) \\
1
\end{array}\right)\left(\begin{array}{llll}
1 & R_{B} \cos \left(\theta_{i}\right) & R_{B} \sin \left(\theta_{i}\right) & 1
\end{array}\right)
$$

If we take $\theta_{i}=0$ and if we multiply the previous matrix by two $4 \times 4$ diagonal matrices:

$$
\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & \frac{1}{R_{A}} & 0 & 0 \\
0 & 0 & \frac{1}{R_{A}} & 0 \\
0 & 0 & 0 & 1
\end{array}\right) \text { and }\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & \frac{1}{R_{B}} & 0 & 0 \\
0 & 0 & \frac{1}{R_{B}} & 0 \\
0 & 0 & 0 & 1
\end{array}\right)
$$

one on the left and the other on the right respectively, we end up with two cylinders with radii equal to $1, \operatorname{Cyl}(1)$.

By multiplying these two matrices by matrix (2.8), one on the right and the other on the left, we get the following matrix:

$$
\left(\begin{array}{cccc}
1 & \frac{r_{B}}{R_{B}} & 0 & 1  \tag{2.9}\\
\frac{r_{A}}{R_{A}} & 0 & 0 & \frac{r_{A}}{R_{A}} \\
0 & 0 & \frac{r_{A} r_{B}}{R_{A} R_{B}} & 0 \\
1 & \frac{r_{B}}{R_{B}} & 0 & 1
\end{array}\right)
$$

So, The matrix (2.8) is $\operatorname{Cyl}\left(r_{A}\right), \operatorname{Cyl}\left(r_{B}\right)$ separable if and only if matrix (2.9) is $C y l(1), C y l(1)$ separable.

If a cylinder decomposition exists for matrix (2.9), then all the $z$ components in the matrix are replaced by 0 to obtain a quantum separable decomposition of a two qubit quantum operator, which means that by taking a vector in the $X-Y$ plane where $z=0$, we end up with a bloch vector decomposition. On the contrary, if a
quantum separable decomposition exists,

$$
\sum_{i} p_{i}\left(1, x_{A}^{i}, y_{A}^{i}, z_{A}^{i}\right) \otimes\left(1, x_{B}^{i}, y_{B}^{i}, z_{B}^{i}\right)
$$

then all the $z$ components are replaced by 1 to get a cylinder decomposition. So we know that determining the operator presented by matrix (2.9) is equivalent to determining a two qubit quantum operator.

When we replace all $z$ components by 0 , the above matrix can be seen as follows:

$$
\left(\begin{array}{cccc}
1 & \frac{r_{B}}{R_{B}} & 0 & 0  \tag{2.10}\\
\frac{r_{A}}{R_{A}} & 0 & 0 & 0 \\
0 & 0 & \frac{r_{A} r_{B}}{R_{R} R_{B}} & 0 \\
0 & 0 & 0 & 0
\end{array}\right)
$$

This means that if matrix (2.10) has a quantum separable decomposition then (2.9) has cylinders separable decomposition with radii " 1 ":

$$
\sum_{i} p_{i}\left(1, x_{A}^{i}, y_{A}^{i}, 1\right) \otimes\left(1, x_{B}^{i}, y_{B}^{i}, 1\right)
$$

Then (2.9) is $C y l(1)$ separable which makes (2.8) is a $C y l\left(r_{A}\right), C y l\left(r_{B}\right)$ separable if and only if $(2.10)$ corresponds to a positive and $P P T$ operator.

Now let's calculate and check the minimal eigenvalues of the operator represented by the matrix (2.10) are non-negative. In the form of Pauli operators, the equation is displayed as follows:

$$
I+\left(\frac{r_{A}}{R_{A}} X \otimes I+\frac{r_{B}}{R_{B}} I \otimes X\right)+\frac{r_{A} r_{B}}{R_{A} R_{B}} Y \otimes Y
$$

And its partial transpose

$$
I+\left(\frac{r_{A}}{R_{A}} X \otimes I+\frac{r_{B}}{R_{B}} I \otimes X\right)-\frac{r_{A} r_{B}}{R_{A} R_{B}} Y \otimes Y
$$

When the partial transpose has been applied, we noticed a changed in the third column which is the same as doing an $X$ rotation. Now, we will work out the eigenvalues of one equation because we found that all the eigenvalues of these two operators are equivalent. Before computing the eigenvalues, we used the Hadamard gate on both qubits to get the following equation:

$$
\begin{equation*}
I+\left(\frac{r_{A}}{R_{A}} Z \otimes I+\frac{r_{B}}{R_{B}} I \otimes Z\right)-\frac{r_{A} r_{B}}{R_{A} R_{B}} Y \otimes Y \tag{2.11}
\end{equation*}
$$

In computational basis, this is equal to

$$
\left(\begin{array}{cccc}
1+f_{A}+f_{B} & 0 & 0 & f_{A} f_{B}  \tag{2.12}\\
0 & 1+f_{A}-f_{B} & -f_{A} f_{B} & 0 \\
0 & -f_{A} f_{B} & 1-f_{A}+f_{B} & 0 \\
f_{A} f_{B} & 0 & 0 & 1-f_{A}-f_{B}
\end{array}\right)
$$

Now, we can find the eigenvalues by dividing the matrix into two blocks: the inner block and the outer block. The determinant of the inner and the outer block are:

$$
\begin{aligned}
& 1-\left(f_{A}-f_{B}\right)^{2}-f_{A}^{2} f_{B}^{2} \\
& 1-\left(f_{A}+f_{B}\right)^{2}-f_{A}^{2} f_{B}^{2}
\end{aligned}
$$

Between these two blocks, the determinant of the outer block is the lowest as $f_{A}$ and $f_{B}$ are positive. As we conclude, if the determinant of the outer block is non-negative then the output of the two cylindrical state spaces will be separable.

### 2.5.2 The diagonal gates acting on two input qubits

In this section, instead of applying the $C Z$ gate, we will consider a diagonal gate that acts on two qubits. This specific gate, $V(\theta)$, is represented by the following matrix:

$$
V(\theta)=\left(\begin{array}{cccc}
1 & 0 & 0 & 0  \tag{2.13}\\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{i \theta}
\end{array}\right)
$$

This matrix can be written in the form:

$$
V(\theta)=|0\rangle\langle 0| \otimes I+|1\rangle\langle 1| \otimes Z_{\theta}
$$

where $\theta$ goes from 0 to $2 \pi$ and $Z_{\theta}=\left(\begin{array}{cc}1 & 0 \\ 0 & e^{i \theta}\end{array}\right)$.

As mentioned previously, to show that the output of this diagonal gate lead to cylindrical separable states that have different radii, we need to consider only extremal input points of the cylinder because if the output from the extremal inputs is separable then the output from all the inputs can be separable because the $V(\theta)$ gate is linear transformation.

To reduce the number of inputs, we take advantage of the symmetry around the $z$ axis. By considering a cylinder separable decomposition,

$$
V(\theta)\left(\rho_{A} \otimes \rho_{B}\right)=\sum_{j} p_{j} m_{A}^{j} \otimes m_{B}^{j}
$$

where $m_{A}^{j} \in \operatorname{Cyl}\left(R_{A}\right)$ and $m_{B}^{j} \in \operatorname{cyl}\left(R_{B}\right)$, we know that $V(\theta)$ gate commutes with $Z$ rotation $U_{z}$ and cylinders are invariant under $Z$ rotations; we have the decomposition:

$$
V(\theta)\left(U_{z}^{A}\left(\rho_{A}\right) \otimes U_{z}^{B}\left(\rho_{B}\right)\right)=\sum_{j} p_{j} U_{z}^{A}\left(m_{A}\right) \otimes U_{z}^{A}\left(m_{B}\right)
$$

If we take two inputs with $y=0$ and we apply the $Z$ rotation, we notice that these inputs can be taken to other inputs. Then, if we have a separable decomposition for the first two inputs, we can obtain separable decompositions for all other inputs by applying the $Z$ rotation.

So, we will concentrate on the two inputs and we will consider all possibilities of $Z,\left(1, r_{A}, 0, \pm 1\right)$ and $\left(1, r_{B}, 0, \pm 1\right)$.

By taking the first possibility where $z=1$ for both inputs,

$$
V(\theta)\left(\left(1, r_{A}, 0,1\right) \otimes\left(1, r_{B}, 0,1\right)\right)=\sum_{j} p_{j} m_{A}^{j} \otimes m_{B}^{j}
$$

In pauli basis, these two inputs are given by

$$
\frac{1}{2}\left(I+r_{A} X+0 Y+Z\right) \otimes \frac{1}{2}\left(I+r_{B} X+0 Y+Z\right)=\left(\begin{array}{cccc}
1 & r_{B} & 0 & 1 \\
r_{A} & r_{A} r_{B} & 0 & r_{A} \\
0 & 0 & 0 & 0 \\
1 & r_{B} & 0 & 1
\end{array}\right)
$$

In the previous part, we could reduce the number of inputs because when $z=-1$ in the first particle is the same as applying an $X$ rotation and a $Z$ rotation to the first and the second particle respectively. Then, the $P P T$ criterion is used when $C Z$ acted on the product term in which $z=1$ in both particles.

Lemma 2. After applying the $V(\theta)$ gate on the matrix $\left(\rho_{A} \otimes \rho_{B}\right)$ i.e. this state is
represented by the following matrix in the pauli basis,

$$
\left(\begin{array}{cccc}
1 & r_{B} & 0 & 1 \\
r_{A} & r_{A} r_{B} & 0 & r_{A} \\
0 & 0 & 0 & 0 \\
1 & r_{B} & 0 & 1
\end{array}\right)
$$

we end up with the following:

- the output top row is identical to the bottom row,
- the left column and the right column of the output are the same.

Hence, $P P T$ criterion can be applied in the possibility where $z=1$ for both inputs.

Proof. We consider $V(\theta)\left(\rho_{A} \otimes \rho_{B}\right) V(\theta)^{\dagger}=\rho_{\text {out }}$. We need to find if the top and the bottom row are the same. So, we subtract these two rows to see if we will end up with 0 . The top row is given by $\operatorname{tr}_{A}\left[\left(I_{A} \otimes I_{B}\right) \rho_{\text {out }}\right]$ and the bottom row is given by $\operatorname{tr}_{A}\left[\left(Z_{A} \otimes I_{B}\right) \rho_{\text {out }}\right]$.

$$
\begin{aligned}
\text { Top - bottom row } & =\operatorname{tr}_{A}\left[\left(I_{A} \otimes I_{B}\right) \rho_{o u t}-\left(Z_{A} \otimes I_{B}\right) \rho_{\text {out }}\right] \\
& =\operatorname{tr}_{A}\left[\left(I_{A}-Z_{A}\right) \otimes I_{B} \rho_{\text {out }}\right] \\
& =\operatorname{tr}_{A}\left[(I-Z)_{A} \otimes I_{B} V(\theta)\left(\rho_{A} \otimes \rho_{B}\right) V(\theta)^{\dagger}\right] \\
& =\operatorname{tr}_{A}\left[2|1\rangle\langle 1| \otimes I_{B}\left(|0\rangle\langle 0| \otimes I+|1\rangle\langle 1| \otimes Z_{\theta}\right)\left(\rho_{A} \otimes \rho_{B}\right) V(\theta)^{\dagger}\right] \\
& =\operatorname{tr}_{A}\left[2|1\rangle\langle 1| \otimes I_{B}\left(|1\rangle\langle 1| \otimes Z_{\theta}\right)\left(\rho_{A} \otimes \rho_{B}\right) V(\theta)^{\dagger}\right] \\
& =\operatorname{tr}_{A}\left[2|1\rangle\langle 1| \otimes I_{B}\left(\rho_{A} \otimes Z_{\theta} \rho_{B}\right) V(\theta)^{\dagger}\right] \\
& =\operatorname{tr}_{A}\left[2|1\rangle\langle 1| \rho_{A} \otimes Z_{\theta} \rho_{B} Z_{\theta}^{\dagger}\right] \\
& =2\langle 1| \rho_{A}|1\rangle Z_{\theta} \rho_{B} Z_{\theta}^{\dagger} ; \text { where }\langle 1| \rho_{A}|1\rangle=0 \text { because it is on the surface of the cylinder. } \\
& =0
\end{aligned}
$$

By obtaining that the top and the bottom row are the same, this argument is the main factor for the use of $P P T$ criterion when $z=1$ for both inputs.

Explicitly, we can go through the same argument for all other possibilities $\left(1, r_{A}, 0, \pm 1\right) \otimes$ $\left(1, r_{B}, 0, \pm 1\right)$, and we find that we end up with the same result. Therefore we know that $P P T$ can be applied in this section.

## The growth factors of different thetas

In this part, we will attempt to find the growth factors for each gate $V(\theta)$ where $\theta$ will change from 0 to $2 \pi$. We will first work out the action of $V(\theta)$ gate on all product of pauli. So, we have to apply a tensor product between each two pauli matrices from the first and the second inputs respectively. For example, we choose $I$ from the first input and then we do a tensor product between the chosen pauli matrix and the $I$, $X, Y$, and $Z$ from the second input. The same is applied for $X, Y$, and $Z$ from the first input (i.e. $I \otimes I, I \otimes X, I \otimes Y, I \otimes Z, X \otimes I, X \otimes X$, etc). Then, we obtained a $4 \times 4$ matrix. Then, we multiply the matrix $V(\theta)$ with each of the calculated $4 \times 4$ pauli matrix and with the conjugate matrix of $V(\theta)$. That's how We end up with matrices of $4 \times 4$. For more details, see Appendix $A$.

After written out these matrices in terms of pauli operators, we note down all the coefficients in this $16 \times 16$ matrix, that we called " $V V(\theta)$ "

$$
\left(\begin{array}{llllllllllllllll}
1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & \mathrm{a} & \mathrm{~b} & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & \mathrm{c} & \mathrm{~d} & 0 \\
0 & \mathrm{~d} & \mathrm{a} & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & \mathrm{~b} & \mathrm{c} & 0 \\
0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & \mathrm{e} & 0 & 0 & \mathrm{j} & \mathrm{p} & 0 & 0 & \mathrm{q} & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & \mathrm{e} & \mathrm{p} & 0 & 0 & \mathrm{p} & \mathrm{j} & 0 & 0 & 0 & 0 & 0 \\
(I \otimes X) \\
0 & 0 & 0 & 0 & 0 & \mathrm{q} & \mathrm{e} & 0 & 0 & \mathrm{k} & \mathrm{p} & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & \mathrm{j} & 0 & 0 & \mathrm{e} & \mathrm{q} & 0 & 0 & \mathrm{p} & 0 & 0 & 0 & 0 \\
(I \otimes Z) \\
(X \otimes I) \\
0 & 0 & 0 & 0 & \mathrm{q} & 0 & 0 & \mathrm{p} & \mathrm{e} & 0 & 0 & \mathrm{j} & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & \mathrm{q} & \mathrm{k} & 0 & 0 & \mathrm{e} & \mathrm{p} & 0 & 0 & 0 & 0 & 0 \\
(X \otimes Y) \\
0 & 0 & 0 & 0 & 0 & \mathrm{j} & \mathrm{q} & 0 & 0 & \mathrm{q} & \mathrm{e} & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & \mathrm{p} & 0 & 0 & \mathrm{q} & \mathrm{j} & 0 & 0 & \mathrm{e} & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 1 & 0 & 0 & 0 \\
(Y \otimes I) \\
0 & \mathrm{c} & \mathrm{~d} & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & \mathrm{a} & \mathrm{~b} & 0 \\
0 & \mathrm{~b} & \mathrm{c} & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & \mathrm{~d} & \mathrm{a} & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 1
\end{array}\right)(Y \otimes Z)
$$

where

$$
a=\frac{1}{2}+\frac{1}{2} \cos \theta ; b=\frac{1}{2} \sin \theta ; c=\frac{1}{2}-\frac{1}{2} \cos \theta ; d=-\frac{1}{2} \sin \theta ;
$$

$e=\frac{1}{2}+\frac{1}{4} \exp (-i \theta)+\frac{1}{4} \exp (i \theta) ; j=\frac{1}{2}-\frac{1}{4} \exp (-i \theta)-\frac{1}{4} \exp (i \theta) ; k=-\frac{1}{2}+\frac{1}{4} \exp (-i \theta)+\frac{1}{4} \exp (i \theta) ;$

$$
p=\frac{1}{4} i \exp (-i \theta)-\frac{1}{4} i \exp (i \theta) ; q=-\frac{1}{4} i \exp (-i \theta)+\frac{1}{4} i \exp (i \theta)
$$

Our goal is to check whether the output is within the output state space which is the convex hull of all the points on the cylinder. So, we are going to theoretically and numerically work out the required growth factor $g$ for different values of $\theta$, where $g$ is a function of theta, to make the output separable. We know two values of $g$ :

- When $\theta$ is 0 , the growth factor $g$ is 1 . Because when $\theta=0, V(1)$ is the identity gate and there is no entangling.
- By taking $z=1$ for both inputs $\left(1, r_{A / B}, 0,1\right)$, we know from the previous section we got $g \geq 2.05817$ when $\theta$ is equal to $\pi$.


## Product input with component $\mathrm{z}=1$

We will start by computing the growth factor theoretically. As we showed previously, to determine if the output of $V(\theta)$ gate leads to cylindrical separable state, we only need to consider extremal input points of the cylinder. We know that $V(\theta)$ gate commutes with $Z$ rotation $U_{z}$ and cylinders are invariant under $Z$ rotations; which means that we are dealing with the two inputs used in the previous section:

$$
\frac{1}{2}\left(I+r_{A} X+0 Y+1\right) \otimes \frac{1}{2}\left(I+r_{B} X+0 Y+1\right)
$$

When we apply $V(\theta)$ to the product input, the matrix is given by:

$$
\left(\begin{array}{c}
1 \\
r_{B} \\
0 \\
1 \\
r_{A} \\
\left(\frac{1}{2}+\frac{1}{4} e^{-i \theta}+\frac{1}{4} e^{i \theta}\right) r_{A} r_{B} \\
\left(-\frac{1}{4} i e^{-i \theta}+\frac{1}{4} i e^{i \theta}\right) r_{A} r_{B} \\
r_{A} \\
0 \\
\left(-\frac{1}{4} i e^{-i \theta}+\frac{1}{4} i e^{i \theta}\right) r_{A} r_{B} \\
\left(\frac{1}{2}-\frac{1}{4} e^{-i \theta}-\frac{1}{4} e^{i \theta}\right) r_{A} r_{B} \\
0 \\
1 \\
r_{B} \\
0 \\
1
\end{array}\right)\left(\begin{array}{c}
I \otimes I \\
I \otimes X \\
X \otimes I \\
X \otimes X \\
X \otimes Y \\
Y \otimes I \\
Y \otimes X \\
Y \otimes Y \\
Y \otimes Z \\
Z \otimes I \\
Z \otimes X \\
Z \otimes Y \\
Z \otimes Z
\end{array}\right)
$$

If we represent the above $16 \times 1$ matrix by a $4 \times 4$ matrix, we have:

$$
\left(\begin{array}{cccc}
1 & r_{B} & 0 & 1  \tag{2.14}\\
r_{A} & \left(\frac{1}{2}+\frac{1}{4} e^{-i \theta}+\frac{1}{4} e^{i \theta}\right) r_{A} r_{B} & \left(-\frac{1}{4} i e^{-i \theta}+\frac{1}{4} i e^{i \theta}\right) r_{A} r_{B} & r_{A} \\
0 & \left(-\frac{1}{4} i e^{-i \theta}+\frac{1}{4} i e^{i \theta}\right) r_{A} r_{B} & \left(\frac{1}{2}-\frac{1}{4} e^{-i \theta}-\frac{1}{4} e^{i \theta}\right) r_{A} r_{B} & 0 \\
1 & r_{B} & 0 & 1
\end{array}\right)
$$

As we mentioned in the previous section, we will show that the matrix can be written as the outer product:

$$
\sum_{i} p_{i}\left(\begin{array}{c}
1 \\
R_{A} \cos \left(\theta_{i}\right) \\
R_{A} \sin \left(\theta_{i}\right) \\
1
\end{array}\right)\left(\begin{array}{llll}
1 & R_{B} \cos \left(\theta_{i}\right) & R_{B} \sin \left(\theta_{i}\right) & 1
\end{array}\right)
$$

So, we multiply the previous matrix by two $4 \times 4$ diagonal matrices:

$$
\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & \frac{1}{R_{A}} & 0 & 0 \\
0 & 0 & \frac{1}{R_{A}} & 0 \\
0 & 0 & 0 & 1
\end{array}\right) \text { and }\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & \frac{1}{R_{B}} & 0 & 0 \\
0 & 0 & \frac{1}{R_{B}} & 0 \\
0 & 0 & 0 & 1
\end{array}\right)
$$

one on the left and the other on the right respectively, we end up with two cylinders with radii equal to $1, \operatorname{Cyl}(1)$.

We get the following matrix:

$$
\left(\begin{array}{cccc}
1 & \frac{r_{B}}{R_{B}} & 0 & 1  \tag{2.15}\\
\frac{r_{A}}{R_{A}} & \left(\frac{1}{2}+\frac{1}{4} e^{-i \theta}+\frac{1}{4} e^{i \theta}\right) \frac{r_{A} r_{B}}{R_{A} R_{B}} & \left(-\frac{1}{4} i e^{-i \theta}+\frac{1}{4} i e^{i \theta}\right) \frac{r_{A} r_{B}}{R_{A} R_{B}} & \frac{r_{A}}{R_{A}} \\
0 & \left(-\frac{1}{4} i e^{-i \theta}+\frac{1}{4} i e^{i \theta}\right) \frac{r_{A} r_{B}}{R_{A} R_{B}} & \left(\frac{1}{2}-\frac{1}{4} e^{-i \theta}-\frac{1}{4} e^{i \theta}\right) \frac{r_{A} R_{B} R_{A}}{R_{A} R_{B}} & 0 \\
1 & \frac{r}{R_{B}} & 0 & 1
\end{array}\right)
$$

So, The matrix (2.14) is $\operatorname{Cyl}\left(r_{A}\right), \operatorname{Cyl}\left(r_{B}\right)$ separable if and only if matrix (2.15) is $C y l(1), C y l(1)$ separable.

Then, we replace all $z$ components by 0 to obtain a quantum separable decomposition of a two qubit quantum operators, the above matrix can be seen as follows:

$$
\left(\begin{array}{cccc}
1 & \frac{r_{B}}{R_{B}} & 0 & 0  \tag{2.16}\\
\frac{r_{A}}{R_{A}} & \left(\frac{1}{2}+\frac{1}{4} e^{-i \theta}+\frac{1}{4} e^{i \theta}\right) \frac{r_{A} r_{B}}{R_{A} R_{B}} & \left(-\frac{1}{4} i e^{-i \theta}+\frac{1}{4} i e^{i \theta}\right) \frac{r_{A} r_{B}}{R_{A} R_{B}} & 0 \\
0 & \left(-\frac{1}{4} i e^{-i \theta}+\frac{1}{4} i e^{i \theta}\right) \frac{r_{A} r_{B}}{R_{A} R_{B}} & \left(\frac{1}{2}-\frac{1}{4} e^{-i \theta}-\frac{1}{4} e^{i \theta}\right) \frac{r_{A} r_{B}}{R_{A} R_{B}} & 0 \\
0 & 0 & 0 & 0
\end{array}\right)
$$

The partial transpose of the above matrix will be equal to:

$$
\left(\begin{array}{cccc}
1 & \frac{r_{B}}{R_{B}} & 0 & 0  \tag{2.17}\\
\frac{r_{A}}{R_{A}} & \left(\frac{1}{2}+\frac{1}{4} e^{-i \theta}+\frac{1}{4} e^{i \theta}\right) \frac{r_{A} r_{B}}{R_{A} R_{B}} & -\left(-\frac{1}{4} i e^{-i \theta}+\frac{1}{4} i e^{i \theta}\right) \frac{r_{A} r_{B}}{R_{A} R_{B}} & 0 \\
0 & \left(-\frac{1}{4} i e^{-i \theta}+\frac{1}{4} i e^{i \theta}\right) \frac{r_{A} r_{B}}{R_{A} R_{B}} & -\left(\frac{1}{2}-\frac{1}{4} e^{-i \theta}-\frac{1}{4} e^{i \theta}\right) \frac{r_{A} r_{B}}{R_{A} R_{B}} & 0 \\
0 & 0 & 0 & 0
\end{array}\right)
$$

Using the the pauli operators, the matrix (2.16) can be written as :

$$
\begin{aligned}
& \quad I+\frac{r_{B}}{R_{B}} I \otimes X+\frac{r_{A}}{R_{A}} X \otimes I+\left[\left(\frac{1}{2}+\frac{1}{4} e^{-i \theta}+\frac{1}{4} e^{i \theta}\right) X \otimes X+\left(\frac{1}{2}-\frac{1}{4} e^{-i \theta}-\frac{1}{4} e^{i \theta}\right) Y \otimes\right. \\
& Y] \frac{r_{A} r_{B} R_{B}}{R_{A} R_{B}}+\left(-\frac{1}{4} i e^{-i \theta}+\frac{1}{4} i e^{i \theta}\right)(X \otimes Y+Y \otimes X) \frac{r_{A} r_{B}}{R_{A} R_{B}}
\end{aligned}
$$

And the partial transpose is given by:

$$
\begin{aligned}
& \quad I+\frac{r_{B}}{R_{B}} I \otimes X+\frac{r_{A}}{R_{A}} X \otimes I+\left[\left(\frac{1}{2}+\frac{1}{4} e^{-i \theta}+\frac{1}{4} e^{i \theta}\right) X \otimes X-\left(\frac{1}{2}-\frac{1}{4} e^{-i \theta}-\frac{1}{4} e^{i \theta}\right) Y \otimes\right. \\
& Y] \frac{r_{A} r_{A} R_{B}}{R_{A}}+\left(-\frac{1}{4} i e^{-i \theta}+\frac{1}{4} i e^{i \theta}\right)(-X \otimes Y+Y \otimes X) \frac{r_{A} r_{B}}{R_{A} R_{B}}
\end{aligned}
$$

When the partial transpose has been applied, we noticed only a change in the third column which is the same as doing an $X$ rotation. When the $X$ rotation is applied, the eigenvalues of the two matrices (2.16) and (2.17) are the same. We checked numerically if the minimum eigenvalues don't change and we found that they are equivalent for different values of theta starting from 0 ending at $2 \pi$.

Now, by taking the symmetric case $R_{A}=R_{B}=R$ and for each value of $R$ going from 1 to 2.1 where the input radius is 1 , we store the value of $R$ only when the minimum eigenvalue for each value of theta becomes negative, it means that when the minimum eigenvalue of a specific theta is negative, it requires a the growth $R$ to be separable. We end up with a graph that represents the storing value of $R$ for each theta.


Figure 2.2: The required growth for each $\theta$ to end up with a separable gate.

It turns out that the $C Z$ gate, where $\theta=\pi$, requires the most growth in $R$.

In these following parts, the same procedure will be repeated but for different product inputs.

## Product input with component $z=1$ and $z=-1$

We will start by applying the diagonal $V(\theta)$ gate on the product input that can be written in the form:

$$
\frac{1}{2}\left(I+r_{A} X+0 Y+1\right) \otimes \frac{1}{2}\left(I+r_{B}+0 Y-1\right)
$$

So $V(\theta)\left(1, r_{A}, 0,1\right) \otimes\left(1, r_{B}, 0,-1\right) V(\theta)^{\dagger}$ is given by a $4 \times 4$ matrix:

$$
\left(\begin{array}{cccc}
1 & r_{B} & 0 & -1 \\
\left(\frac{1}{2} e^{-i \theta}+\frac{1}{2} e^{i \theta}\right) r_{A} & \left(\frac{1}{2}+\frac{1}{4} e^{-i \theta}+\frac{1}{4} e^{i \theta}\right) r_{A} r_{B} & \left(-\frac{1}{4} i e^{-i \theta}+\frac{1}{4} i e^{i \theta}\right) r_{A} r_{B} & \left(-\frac{1}{2} e^{-i \theta}-\frac{1}{2} e^{i \theta}\right) r_{A} \\
\left(-\frac{1}{2} i e^{-i \theta}+\frac{1}{2} i e^{i \theta}\right) r_{A} & \left(-\frac{1}{4} i e^{-i \theta}+\frac{1}{4} i e^{i \theta}\right) r_{A} r_{B} & \left(\frac{1}{2}-\frac{1}{4} e^{-i \theta}-\frac{1}{4} e^{i \theta}\right) r_{A} r_{B} & \left(\frac{1}{2} i e^{-i \theta}-\frac{1}{2} i e^{i \theta}\right) r_{A} \\
1 & r_{B} & 0 & -1
\end{array}\right)
$$

This matrix is separable if and only if the following matrix is $\operatorname{cyl}(1), \operatorname{cyl}(1)$ separable

$$
\left(\begin{array}{cccc}
1 & \frac{r_{B}}{R_{B}} & 0 & -1 \\
\left(\frac{1}{2} e^{-i \theta}+\frac{1}{2} e^{i \theta}\right) \frac{r_{A}}{R_{A}} & \left(\frac{1}{2}+\frac{1}{4} e^{-i \theta}+\frac{1}{4} e^{i \theta}\right) \frac{r_{A} r_{B}}{R_{A} R_{B}} & \left(-\frac{1}{4} i e^{-i \theta}+\frac{1}{4} i e^{i \theta}\right) \frac{r_{A} r_{B}}{R_{A} R_{B}} & \left(-\frac{1}{2} e^{-i \theta}-\frac{1}{2} e^{i \theta}\right) \frac{r_{A}}{R_{A}} \\
\left(-\frac{1}{2} i e^{-i \theta}+\frac{1}{2} i e^{i \theta}\right) \frac{r_{A}}{R_{A}} & \left(-\frac{1}{4} i e^{-i \theta}+\frac{1}{4} i e^{i \theta}\right) \frac{r_{A} r_{A}}{R_{B} R_{B}} & \left(\frac{1}{2}-\frac{1}{4} e^{-i \theta}-\frac{1}{4} e^{i \theta}\right) \frac{r_{A} r_{B}}{R_{A} R_{B}} & \left(\frac{1}{2} i e^{-i \theta}-\frac{1}{2} i e^{i \theta}\right) \frac{r_{A}}{R_{A}} \\
1 & \frac{r_{B}}{R_{B}} & 0 & -1
\end{array}\right)
$$

By replacing all $z$ component by 0 , a quantum separable decomposition is obtained,

$$
\left(\begin{array}{cccc}
1 & \frac{r_{B}}{R_{B}} & 0 & 0 \\
\left(\frac{1}{2} e^{-i \theta}+\frac{1}{2} e^{i \theta}\right) \frac{r_{A}}{R_{A}} & \left(\frac{1}{2}+\frac{1}{4} e^{-i \theta}+\frac{1}{4} e^{i \theta}\right) \frac{r_{A} r_{B}}{R_{A} R_{B}} & \left(-\frac{1}{4} i e^{-i \theta}+\frac{1}{4} i e^{i \theta}\right) \frac{r_{A} r_{B}}{R_{A} R_{B}} & 0 \\
\left(-\frac{1}{2} i e^{-i \theta}+\frac{1}{2} i e^{i \theta}\right) \frac{r_{A}}{R_{A}} & \left(-\frac{1}{4} i e^{-i \theta}+\frac{1}{4} i e^{i \theta}\right) \frac{r_{A} r_{A}}{R_{B} R_{B}} & \left(\frac{1}{2}-\frac{1}{4} e^{-i \theta}-\frac{1}{4} e^{i \theta}\right) \frac{r_{A} r_{B}}{R_{A} R_{B}} & 0 \\
0 & 0 & 0 & 0
\end{array}\right)
$$

The matrix can be written as:

$$
\begin{gathered}
I+\frac{r_{B}}{R_{B}} I \otimes X+\left[\left(\frac{1}{2} e^{-i \theta}+\frac{1}{2} e^{i \theta}\right) X \otimes I+\left(-\frac{1}{2} i e^{-i \theta}+\frac{1}{2} i e^{i \theta}\right) Y \otimes I\right] \frac{r_{A}}{R_{A}}+\left[\left(\frac{1}{2}+\frac{1}{4} e^{-i \theta}+\right.\right. \\
\left.\left.\frac{1}{4} e^{i \theta}\right) X \otimes X+\left(-\frac{1}{4} i e^{-i \theta}+\frac{1}{4} i e^{i \theta}\right)(X \otimes Y+Y \otimes X)+\left(\frac{1}{2}-\frac{1}{4} e^{-i \theta}-\frac{1}{4} e^{i \theta}\right) Y \otimes Y\right] \frac{r_{A} r_{B}}{R_{A} R_{B}}
\end{gathered}
$$

The partial transpose of the matrix is given by

$$
\left(\begin{array}{cccc}
1 & \frac{r_{B}}{R_{B}} & 0 & 0 \\
\left(\frac{1}{2} e^{-i \theta}+\frac{1}{2} e^{i \theta}\right) \frac{r_{A}}{R_{A}} & \left(\frac{1}{2}+\frac{1}{4} e^{-i \theta}+\frac{1}{4} e^{i \theta}\right) \frac{r_{A} r_{B}}{R_{A} R_{B}} & \left(-\frac{1}{4} i e^{-i \theta}+\frac{1}{4} i e^{i \theta}\right) \frac{r_{A} r_{B}}{R_{A} R_{B}} & 0 \\
-\left(-\frac{1}{2} i e^{-i \theta}+\frac{1}{2} i e^{i \theta}\right) \frac{r_{A}}{R_{A}} & -\left(-\frac{1}{4} i e^{-i \theta}+\frac{1}{4} i e^{i \theta}\right) \frac{r_{A} r_{A}}{R_{B} R_{B}} & -\left(\frac{1}{2}-\frac{1}{4} e^{-i \theta}-\frac{1}{4} e^{i \theta}\right) \frac{r_{A} r_{B} R_{B}}{R_{A} R_{B}} & 0 \\
0 & 0 & 0 & 0
\end{array}\right)
$$

it is same as:
$I+\frac{r_{B}}{R_{B}} I \otimes X+\left[\left(\frac{1}{2} e^{-i \theta}+\frac{1}{2} e^{i \theta}\right) X \otimes I-\left(-\frac{1}{2} i e^{-i \theta}+\frac{1}{2} i e^{i \theta}\right) Y \otimes I\right] \frac{r_{A}}{R_{A}}+\left[\left(\frac{1}{2}+\frac{1}{4} e^{-i \theta}+\right.\right.$ $\left.\left.\frac{1}{4} e^{i \theta}\right) X \otimes X+\left(-\frac{1}{4} i e^{-i \theta}+\frac{1}{4} i e^{i \theta}\right)(X \otimes Y-Y \otimes X)-\left(\frac{1}{2}-\frac{1}{4} e^{-i \theta}-\frac{1}{4} e^{i \theta}\right) Y \otimes Y\right] \frac{r_{A} r_{B}}{R_{A} R_{B}}$

As before, when we apply $P P T$ criterion to the matrix, we end up with negative signs in the third row or in the third column, which also can be the same as applying the $X$ rotation. Then, we will check numerically the eigenvalues of these last two matrices if they are equal. We notice that the minimum eigenvalues are the same. In addition, we notice that the minimum eigenvalues of these matrices are identical to those of matrices (2.16) and (2.17).

As mentioned in the previous part, we are going to find the necessary value of growth for each theta to end up with separable gates.

By taking the symmetric case, $R_{A}=R_{B}=R$ and by changing the $z$ component of the second particle in the product input, it turns out that the same graph (2.2) is obtained and the $C Z$ gate requires the most growth in $R$.

For the two remaining cases where the product input has

- component $z=-1$ and $z=1$ for the first and second particle respectively,
- and $z=-1$ for both particles,
the same procedure has been applied. A $4 \times 4$ matrix is obtained after $V(\theta)$ acting on the product input. This matrix is $\operatorname{Cyl}\left(r_{A}\right), C y l\left(r_{B}\right)$ separable if and only if $C y l(1), C y l(1)$ is separable. Then all the components of $z$ in the matrix are replaced by 0 , to have a quantum separable decomposition. $P P T$ is applied, as before, we notice that all minimum eigenvalues are identical. We compare all the minimum eigenvalues with the previous two cases, where $z=1$ for both particles and $z=-1$ for the second particle, we notice that these eigenvalues remain the same.

By considering the symmetric case where $R_{A}=R_{B}=R$, we produce the graph that represents the required value of growth or $R$ for each $\theta$ going from 0 to $2 \pi$ to have separable gates. The obtained graphs of these two cases are exactly the same as the graph (2.2).

### 2.6 Summary of chapter 2

We started by relaxing the notion of separability to get a new generalised notion that we called cylinder separability, where, instead of Bloch sphere, we had cylinders.

Then we applied this notion to systems where we changed the inputs $|\psi\rangle$ of the cluster state quantum computation, i.e. the modified input state has the unormalised form $|\psi\rangle=|0\rangle+\epsilon|1\rangle$. Then we find that, after applying $C Z$ gates on two cylindrical state spaces with radii $r$, the output is separable with respect to two different cylindrical state spaces with larger radii $R=\lambda r$. Separability can be maintained by growing the cylinder by $\lambda$ but we must not grow beyond radius 1 , because this state space should be in the dual of the cylinder measurements. Then, instead of the $C Z$ gates, we applied diagonal gates $V(\theta)$ on two cylindrical state spaces. We tried to find the minimum growth required for separability for each value of $\theta$ and it turned out that the $C Z$ gate is the gate that requires the most growth factor $\lambda$ to maintain separability. In the following chapter, we will find that, by growing the size of new state spaces, the system can be efficiently simulated classically.

## Chapter 3

## Different state spaces

To efficiently simulate a complex quantum system on a classical computer, we approximate a pure entangled quantum state by a separable decomposition. As we mentioned in the previous chapter, we use a state space as the cylinder instead of Bloch sphere to get a separable decomposition.

In this chapter, we will try to find a non-quantum state space that is better than the cylinder. So, we are going to try to find a state space that starts with greater quantum inputs than the cylinder, but grows more slowly each time $C Z$ gate has been applied to maintain the separable decomposition, so that we do not exit the dual of the allowed measurement. However, it turns out that the cylinder is the most optimal state space among all the state spaces that we searched through, which means that the cylinder requires the lowest growth of $\sqrt{2+\sqrt{5}} \approx 2.05817$.

### 3.1 State spaces

We consider state spaces that are convex hull of $N(5-2)+2$ vertices where $N$ is an integer number and these vertices are described by Bloch vectors. The considered state spaces are chosen because of the symmetry about the $z$ axis and they are among the simplest we could try for relatively small number of extremal points. This study requires small number of extremal points because we are going to use Linear Programming on Matlab on a regular laptop to work out numerically the growth factors to see if there is a state space grows less than the cylinder. If our trial shapes have too many points, the algorithms run prohibitively slowly.

We will begin by describing the boundary of these state spaces in the positive $x$ part of the $X-Z$ plane only. Then this cross-section will be rotated around the $z$ axis $N$ times to construct the whole shape.

The boundary of these state spaces in the positive $x$ part of the $X-Z$ plane is defined by five points. These points are $Q=(0,1), S=(0,-1)$, which represent
the top and the bottom of the state space respectively, $T=(1,0)$ which represents a point at the "equator", and two other points $K=(w, h)$ and $L=(w,-h)$ where " $h$ " represents a chosen height and " $w$ " represents a related width. For instance, the following two dimensional graph illustrates the boundary we get in the $X-Z$ plane for $h=0.5$ and $w=\sqrt{\frac{3}{4}}$.


Figure 3.1: The boundary of the state space in the $X-Z$ plane for $h=0.5$ and $w=\sqrt{\frac{3}{4}}$

We will repeat this boundary by rotating it $N$ times about the $z$ axis, at regular azimuthal angle intervals $\theta=n \frac{2 \pi}{N}$, where $n=0, \ldots, N-1$, so overall, the vertices of the state spaces will be given by:

- $(0,0,1)$ represents the top of the state space
- $(0,0,-1)$ represents the bottom of the state space
- $(w \cos (\theta), w \sin (\theta), \pm h)$
- $(\cos \theta, \sin (\theta), 0)$ where these points are located on the horizontal plane

The following image illustrates the state space with $h=0.5, w=\sqrt{\frac{3}{4}}$ and $N=6$ which is the total number of angles.


Figure 3.2: The state space for $N=6, h=0.5$ and $w=\sqrt{\frac{3}{4}}$

This means that the shape of the state space that we will consider, is parameterized by three numbers: $N, h$ and $w$. We will change these values to find the growth factor needed to maintain separability.

But, to see the relationship between these shapes and the bloch sphere, we are going to define a new parameter $\epsilon$ to replace the parameter $w$, where we have:

$$
\begin{equation*}
w=\left(\sqrt{1-h^{2}}\right)(1+\epsilon) \tag{3.1}
\end{equation*}
$$

When $\epsilon=0$ all the vertices are on the bloch sphere surface. When $\epsilon$ is bigger than 0 , some of the vertices protrude outside the Bloch sphere. Then we will vary the value of $\epsilon$, which does not need to be small, to see what happens to the growth factor when we step away from the Bloch sphere. Note however that when the value of $\epsilon$ is too high, where $w>1$,

$$
\begin{equation*}
\epsilon>\frac{1}{\sqrt{1-h^{2}}}-1 \tag{3.2}
\end{equation*}
$$

the shape of these state spaces changes, and we end up with a non-convex shape. For instance, the following graph represents the obtained boundary in the $X-Z$ plane for $h=0.5$ and $w=\sqrt{\frac{3}{4}}$ when $\epsilon=1.1$; in theses cases, we are going to take the convex hull which reduce the number of extremal points and change the outermost boundary to a flat vertical face as illustrated in the figure,


Figure 3.3: The boundary of the convex hull in the $X-Z$ plane for $h=0.5, \epsilon=1.1$ and $w=\sqrt{\frac{3}{4}}$

As a conclusion, we are going to use three parameters: $N, h$ and $\epsilon$. When $\epsilon \leq$ $\frac{1}{\sqrt{1-h^{2}}}-1$, we end up with cross-section which is illustrated by the following graph, and when $\epsilon>\frac{1}{\sqrt{1-h^{2}}}-1$, we end up with cross-section with a flat boundary as


Figure 3.4: The boundary of the convex hull in the $X-Z$ plane for $h=0.75, \epsilon=0.1$ and $w=\sqrt{\frac{7}{16}}$
illustrated, for example, in the following figure,


Figure 3.5: The boundary of the convex hull in the $X-Z$ plane for $h=0.75, \epsilon=1.1$ and $w=\sqrt{\frac{7}{16}}$

Our goal in this chapter is to use Linear Programming to work out numerically the growth factors for different values of $N, h$ and $\epsilon$ to see whether any of these shapes may grow slower than the cylinder.

### 3.2 Using Linear Programming to decide the convex hull membership

The problem that we are working on is whether we end up with separability, it means that we need to work out whether a point is in the convex hull of a set of points, in our case we are going to use the Linear Programming [18].

In this section, we will describe a general method of using Linear Programming to decide if a point is inside the convex hull. We start by considering a set $G$ of vectors where $G=\left\{g_{1}, \ldots, g_{d}\right\}$ where the $g_{i} \in \mathbb{R}^{n}$. The state space is the convex hull of $G$ and we attempt to work out whether a vector $v$ is in the convex hull of $G$. So we use the Linear Programming to check if the vector is inside the convex hull and if the problem is feasible.

The Linear Programming consists of minimizing $C^{T} x$ where $C$ is Linear objective function. The minimisation is subject to a linear equality and linear inequality constraints where both of them consist of a matrix and a vector.

The equality constraint states that $\sum_{i} x_{i} g_{i}=v$. This constraint can be presented as:

$$
\left(\begin{array}{llll} 
& & & \\
g_{1} & g_{2} & \ldots & \\
& & &
\end{array}\right)\left(\begin{array}{c}
x_{1} \\
x_{2} \\
\vdots \\
\\
x_{d}
\end{array}\right)=v
$$

If $x$ satisfies this equality constraint, it means that $v$ can be written as a linear combination of $g_{i}$ vectors.
In order to have $v$ inside the convex hull of $G, x$ should be positive numbers and $\sum_{i} x_{i}=1$. Then in our case, the linear inequality constraint consists of an identity matrix $I$ and a zero vector $b$, where we have,

$$
\left(\begin{array}{ccccccc}
1 & 0 & 0 & 0 & \ldots & 0 \\
0 & 1 & 0 & 0 & \ldots & & 0 \\
& & & & & \\
0 & 0 & 0 & 0 & \ldots & & 1
\end{array}\right)\left(\begin{array}{c}
x_{1} \\
x_{2} \\
\vdots \\
x_{d}
\end{array}\right) \geq b
$$

So $v$ is not just a linear combination but it is a convex combination of $g_{i}$ vectors.
We are not interested in optimization, we only care whether these constraints can
be satisfied to get $v$ in the convex hull of $G$. So we can choose any vector for the linear objective function but we pick the zero vector.

After applying the Linear Programming, if the problem is feasible then the vector $v$ is inside the convex hull but if the problem is infeasible, then $v \notin G$.

### 3.3 Linear Programming

In this section, we will see how to cast our problem in the form that we discussed in the previous section. For a particular shape, the set $G$ contains vectors $g_{i} . g_{1}$ and $g_{2}$ are the vectors that represent the top and the bottom of the state space,

$$
g_{1}=\left(\begin{array}{l}
1 \\
0 \\
0 \\
1
\end{array}\right) \text { and } g_{2}=\left(\begin{array}{c}
1 \\
0 \\
0 \\
-1
\end{array}\right)
$$

For $g_{j}$ where $j \geq 3$, these vectors have two forms. The first form is given by

$$
\left(\begin{array}{c}
1 \\
R w(1+\epsilon) \cos (\theta) \\
R w(1+\epsilon) \sin (\theta) \\
\pm h
\end{array}\right)
$$

where $R$ is the output radius, $\theta=n \frac{2 \pi}{N}$ and $n=0, \ldots, N-1$, and the second form represents the points located on the horizontal plane,

$$
\left(\begin{array}{c}
1 \\
\cos (\theta) \\
\sin (\theta) \\
0
\end{array}\right)
$$

Note that $g_{i}$ begin always with 1 which is the coefficient of the pauli operator $I$ because of normalisation.

After putting all these $g_{i}$ together, we end up with a $4 \times(N(5-2)+2)$ output matrix, $m_{\text {output state space. This matrix }}$ describes the shape, for instance by using Matlab, we take $R=2.468, h=0.5, w=\sqrt{\frac{3}{4}}, \epsilon=0$ and $N=6$, we obtain the following
output matrix,

$$
m_{\text {output state space }}=\left(\begin{array}{cccccc}
1 & 1 & 1 & 1 & \ldots & 1 \\
0 & 0 & 2.14 & 1.07 & \ldots & 1.07 \\
0 & 0 & 0 & 1.85 & \ldots & -1.85 \\
1 & -1 & 0.5 & 0.5 & \ldots & -0.5
\end{array}\right)
$$

Each of these columns represents a vertex in the state space.
In our case, $G$ is going to be the product of $m_{\text {output state space }}$. So, the set $G$ which is represented by a $\left(16 \times(N(5-2)+2)^{2}\right)$ matrix, is given by,

$$
G=m_{\text {output state space }} \otimes m_{\text {output state space }}
$$

The matrix $m_{\text {input state space }}$ has the same coefficients as $m_{\text {output state space }}$ but the $X$ and $Y$ coordinates are not multiplied by $R$. Each column in $m_{\text {inputstatespace }}$ represents a vector $t_{i}$ where $i=1, \ldots, N(5-2)+2$. Then $v$ is constructed from a particular pair of inputs $t_{i} \otimes t_{j}$ acting with the $C Z$ gate. Explicitly, each $v$ is is going to be in the following form,

$$
v=C Z\left(t_{i} \otimes t_{j}\right)
$$

We need to work out whether all the $t_{i} \otimes t_{j}$ inputs are taken to separable outputs. But we are not going to test the $(N \times(5-2)+2)^{2}$ inputs, we only need to take $5^{2}$ inputs because the $C Z$ commutes with the $Z$ rotation; because of this property, we have the following separable decomposition:

$$
C Z\left(U_{z i}\left(t_{i}\right) \otimes U_{z j}\left(t_{j}\right)\right)=\sum_{k} p_{k} U_{z i}\left(\lambda_{i}\right) \otimes U_{z j}\left(\lambda_{j}\right)
$$

If $C Z$ is applied on the measurement operators, $X$ and $Y$ measurements will be rotated. Then if we apply the $Z$ rotation on a separable decomposition for inputs with $y=0$, we find that the other inputs with $y=\sin (\theta)$ as a coefficient of $Y$ will have separable decomposition.

So the inputs $t_{i}$ are the following vectors where $y$ has the coefficient 0 ,

$$
t_{1}=\left(\begin{array}{l}
1 \\
0 \\
0 \\
1
\end{array}\right), t_{2}=\left(\begin{array}{c}
1 \\
0 \\
0 \\
-1
\end{array}\right), t_{3}=\left(\begin{array}{l}
1 \\
1 \\
0 \\
0
\end{array}\right) \text { and for } t_{j} \geq 4 t_{j}=\left(\begin{array}{c}
1 \\
w(1+\epsilon) \\
0 \\
\pm h
\end{array}\right)
$$

where $t_{1}, t_{2}$ represent the top and the bottom of the state space respectively, $t_{3}$
represents the point on the horizontal line and the other vectors have the form $t_{j}$. For example, by using Matlab, we take $R=2.468, h=0.5, w=\sqrt{\frac{3}{4}}, \epsilon=0$ and $N=6$, we present one pair of input $m_{\text {input state space }}$,

$$
m_{\text {input state space }}=\left(\begin{array}{cccccc}
1 & 1 & 1 & 1 & \ldots & 1 \\
0 & 0 & 0.866 & 0.433 & \ldots & 0.433 \\
0 & 0 & 0 & 0.75 & \ldots & -0.75 \\
1 & -1 & 0.5 & 0.5 & \ldots & -0.5
\end{array}\right)
$$

We will vary the output radius $R$ across a range until the problem becomes feasible, then this specific $R$ is taken as the minimum growth required to have separability.

We start by fixing the output radius on the first value of the chosen range for a specific value of $N$. We will fix the first input $t_{1}$ and apply the $C Z$ gate on a tensor product between the fixed input and itself. Then we use the Linear Programming to check if it ends up with a separable decomposition. So we have two cases:

1. If it has such decomposition, we do the same process but we change the tensor product, which will occur between the fixed and the second input $t_{2}$. If we get a separable decomposition, the same procedure will be repeated until we end up with a separable decomposition for all the 5 inputs. Then, we fix the second input instead of the first and re-do the same technique until we obtain a separable decomposition for a tensor product between the final input and itself. If all these tensor products have separable decomposition then the output radius is considered as the optimal one. But if in the middle of this procedure, we don't reach a separable decomposition, the second case will be applied.
2. If we don't obtain a separable decomposition, then the same procedure will be replicated with the second value of output radius until it reaches the last required output radius in the range.

The same procedure is repeated but for different values of $N$. If the procedure ends up by giving us a minimum value of the growth factor (or a minimum output radius by taking input radius equal to 1 ), this value will be compared with the growth (or the output radius) of the cylinder to choose which state space is optimal. But if the process has gone through all the radii and we don't have a separable decomposition, we know that the cylinder is the optimal state space.

### 3.4 Results

Our state spaces are characterised by $h, N$ and $\epsilon$. We vary these parameters to get a minimum growth factor. The following table shows, for each state space, the minimum radius growth required for all the pairs of inputs to end up with a separable output.

| height | $\epsilon$ | $\mathbf{N}$ | Growth factor or R |
| :---: | :---: | :---: | :---: |
| 0.5 | 0 | 6 | 2.468 |
|  |  | 10 | 2.442 |
|  |  | 15 | 2.442 |
|  | 0.1 | 20 | 2.415 |
|  |  | 10 | 2.507 |
|  |  | 15 | 2.434 |
|  |  | 20 | 2.443 |
|  | 1.1 | 6 | 2.415 |
|  |  | 10 | 2.528 |
|  |  | 15 | 2.434 |
|  |  | 20 | 2.443 |
| 0.75 | 0 | 6 | 2.415 |
|  |  | 10 | 2.514 |
|  |  | 15 | 2.439 |
|  |  | 20 | 2.443 |
|  | 0.1 | 6 | 2.415 |
|  |  | 10 | 2.528 |
|  |  | 15 | 2.452 |
|  | 1.1 | 6 | 2.443 |
|  |  | 10 | 2.415 |
|  |  | 15 | 2.528 |
|  |  | 20 | 2.452 |
|  |  | 2.415 |  |

Table 3.1: The radius growth for different state spaces

| h | $\epsilon$ | N | Growth factor or R |
| :---: | :---: | :---: | :---: |
| 0.8 | 0 | 6 | 2.528 |
|  |  | 10 | 2.434 |
|  |  | 15 | 2.443 |
|  |  | 20 | 2.415 |
|  | 0.1 | 6 | 2.528 |
|  |  | 10 | 2.452 |
|  |  | 15 | 2.443 |
|  |  | 20 | 2.415 |
|  | 1.1 | 6 | 2.528 |
|  |  | 10 | 2.452 |
|  |  | 15 | 2.443 |
|  |  | 20 | 2.415 |
| 0.6 | 0 | 6 | 2.484 |
|  |  | 10 | 2.447 |
|  |  | 15 | 2.439 |
|  |  | 20 | 2.415 |
|  | 0.1 | 6 | 2.515 |
|  |  | 10 | 2.447 |
|  |  | 15 | 2.441 |
|  |  | 20 | 2.415 |
|  | 1.1 | 6 | 2.528 |
|  |  | 10 | 2.447 |
|  |  | 15 | 2.441 |
|  |  | 20 | 2.415 |
| 0.4 | 0 | 6 | 2.465 |
|  |  | 10 | 2.436 |
|  |  | 15 | 2.441 |
|  |  | 20 | 2.415 |
|  | 0.1 | 6 | 2.500 |
|  |  | 10 | 2.448 |
|  |  | 15 | 2.441 |
|  |  | 20 | 2.415 |
|  | 1.1 | 6 | 2.500 |
|  |  | 10 | 2.448 |
|  |  | 15 | 2.441 |
|  |  | 20 | 2.415 |

Table 3.2: The radius growth for different state spaces

| $\mathbf{h}$ | $\epsilon$ | $\mathbf{N}$ | Growth factor or $\mathbf{R}$ |
| :---: | :---: | :---: | :---: |
| 0.1 | 0 | 6 | 2.488 |
|  |  | 10 | 2.432 |
|  |  | 15 | 2.432 |
|  | 0.1 | 20 | 2.415 |
|  |  | 10 | 2.484 |
|  |  | 15 | 2.429 |
|  |  | 20 | 2.429 |
|  | 1.1 | 6 | 2.415 |
|  |  | 10 | 2.484 |
|  |  | 15 | 2.429 |
|  |  | 20 | 2.429 |
| 0.25 | 0 | 6 | 2.415 |
|  |  | 10 | 2.458 |
|  |  | 15 | 2.431 |
|  |  | 20 | 2.438 |
|  | 0.1 | 6 | 2.415 |
|  |  | 10 | 2.468 |
|  |  | 15 | 2.444 |
|  |  | 20 | 2.443 |
|  | 1.1 | 6 | 2.415 |
|  |  | 10 | 2.468 |
|  |  | 15 | 2.444 |
|  |  | 20 | 2.443 |
|  |  | 2.415 |  |

Table 3.3: The radius growth for different state spaces

As a conclusion of this section, when we change the height for $N=20$, which is the total number of points around the equator, we notice that the minimum value of growth factor is 2.145 for all the considered state spaces in both cases: within and outside the Bloch sphere. This means that the cylinder with a growth factor 2.05817 is the most optimal state space among all the state spaces that we studied.

### 3.5 Summary of chapter 3

In this chapter, we started by considering new state spaces that are the convex hull of $N(5-2)+2$ vertices. Then, we grew the size of these spaces to maintain the non-entangled representation. We also used the Linear programming to work out numerically the growth factors to see whether any of these state spaces may grow slower than the cylinder. It turned out that the minimum value of growth factor for the considered state spaces is 2.145 which is bigger than the cylinder growth factor.

Therefore, the cylinder that requires the least growth, is the most optimal state space among all the state spaces that we studied. In the following chapter, we will try to increase the region that can be efficiently simulated classically, using the method of the coarse graining.

## Chapter 4

## Coarse Graining

When a $C Z$ gate is applied on two cylindrical state spaces, we must grow the radius $r$ in order to describe the system with separable decomposition. However, $r$ cannot grow too much, i.e. $r$ must not grow beyond 1, because the separable decomposition must use state spaces that are in the dual of the cylinder measurements. In this chapter, we are going to introduce a method where we put particles into blocks and these blocks are treated as single systems. This method is called coarse graining. So, we are going to try to use coarse graining to increase the size of the region that we can simulate efficiently, and we will find that this enables us to improve the value of the radius $r$ by putting two particles and then four particles in a block.

### 4.1 Description of coarse graining

We start by taking a lattice similar to the cluster state and we divide it into blocks. Each block contains vertices that are connected by several $C Z$ gates. The $C Z$ gates that are within the block, are defined as "internal $C Z$ gates". We also connect vertices from different blocks by $C Z$ gates that we called "external $C Z$ gate". For instance, we are considering a $4 \times 4$ lattice in each block:

$$
\begin{array}{|cccc|cccc|}
\hline \cdot & \cdot & \cdot & \circ & \circ & \cdot & \cdot & \cdot \\
\cdot & \cdot & \cdot & 0 & \circ & \cdot & \cdot & \cdot \\
\cdot & \cdot & \cdot & \circ & \circ & \cdot & \cdot & \cdot \\
\cdot & \cdot & \cdot & \circ & \circ & \cdot & \cdot & \cdot \\
\hline
\end{array}
$$

In this diagram, the dots and the circles within each block are related by internal $C Z$ gates. The circles in the first block are connected to those in the second block by the external $C Z$ gates.

In each block that is considered as a single particle, we are going to grow the radii
of the cylindrical state spaces according to the applied external $C Z$ gates to maintain separable decomposition. By taking the previous example, where we considered a $4 \times 4$ lattice in each block, $C y l(r)$ are represented by the dots and $C y l(r \lambda)$ are represented by the circles where $\lambda=\sqrt{2+\sqrt{5}} \simeq 2.05817$, which was found in the previous chapter, then $r \leq \frac{1}{\lambda}$.

Then the operators, that are within the block, are created from these cylindrical state spaces $C y l\left(r \lambda^{i}\right)$, where $i$ is the number of the external $C Z$ gates, and the internal $C Z$ gates that are applied on these state spaces. The set of these operators is defined as "block state space".

By growing the radii of the state spaces, we construct separable decomposition before applying the internal $C Z$ gates. Now, we work out whether the separable decomposition has positive outcomes for allowed measurements.

We start by taking one particle in the block state space, it means that the internal $C Z$ gate cannot be applied, only four external $C Z$ gates are applied to this cylindrical state space.

$C y l\left(r \lambda^{4}\right)$ are given by the circle and the dots may grow only by $\lambda$ or they don't grow. But we won't consider cylinder presented by dots because if $C y l\left(r \lambda^{4}\right)$ satisfies the dual constraint then the other cylinders that have $r$ or $r \lambda$ will be in the dual. By taking the allowed measurements, we compute the trace that should be positive. Then $r \lambda^{4}$ must be less than 1 which means that $r \leq \frac{1}{r \lambda^{4}} \approx 0.0557$.

Now we consider the case where a block consists of two particles. Then, we have only one internal gate and three external gates.

$C y l\left(r \lambda^{3}\right)$ are given by the big circles and the small circles may grow only by $\lambda$ or $\lambda^{2}$ and the dots represent $C y l(r)$. As mentioned above, we won't consider cylinders that are presented by dots and small circles.

All the extremal inputs of a cylindrical state space have the following form:

$$
\left(1, r_{2 \text { cylinders }} \cos \theta, r_{2 \text { cylinders }} \sin \theta, \pm 1\right)
$$

where $r_{2 c y l i n d e r s}$ is defined by $r_{2 c y l i n d e r s}:=r \lambda^{3}$ and $\theta$ varies between 0 and $2 \pi$. We are going to work out the maximum value $r_{2 c y l i n d e r s}$ to see whether the separable decomposition has positive measurement outcomes.

Instead of taking all the possible inputs, we use the following to reduce the number of inputs:

- If we apply $C Z$ on two $Z$ measurements, there is no negativity because we end up with the same operators. However, if the $C Z$ gate acts on a $Z$ and an $X-Y$ measurements, we end up with a rotation of the $X-Y$ measurement. Then the product measurement on the two particles does not affect the positivity of the outcomes. Explicitly, we assume that $r_{2 c y l i n d e r s}$ should not grow beyond 1 in order to satisfy the dual constraint. Therefore we do not need to consider the $Z$ measurement.
- We consider that $y$ components of $C y l\left(r_{2 c y l i n d e r s}\right)$ are equal to 0 . Hence, if we consider $C y l\left(r_{2 \text { cylinders }}\right)$ with $z=-1$, we find that it is the same as applying an $X$ rotation on the first particle with $z=1$ and $Z$ rotation on the second particle.
- $C Z$ commutes with the $Z$ rotation. If we take the projectors of $X-Y$ plane measurements and we apply the $Z$ rotation, we end up with projectors that are in a different direction. Instead of taking $C y l\left(r_{2 c y l i n d e r s}\right)$ with component $y=0$, we are going to consider the measurement projector $\left(\frac{I-X}{2}\right)$.

Now, we need to compute the $\operatorname{tr}(\rho \mathrm{P}) \geq 0$, where

$$
\rho=C Z\left(1, r_{2 \text { cylinders }} \cos \theta_{A} X, r_{2 \text { cylinders }} \sin \theta_{A} Y, Z\right) \otimes\left(1, r_{2 \text { cylinders }} \cos \theta_{B} X, r_{2 \text { cylinders }} \sin \theta_{B} Y, Z\right),
$$

where

$$
z=1 \text { and } P=\left(\frac{I-X}{2}\right)^{\otimes 2}
$$

So we end up with the following equation,

$$
\operatorname{tr}(\rho \mathrm{P})=\mathrm{r}_{2 \text { cylinders }}^{2} \sin \theta_{\mathrm{A}} \sin \theta_{\mathrm{B}}-\mathrm{r}_{2 \text { cylinders }}\left(\cos \theta_{\mathrm{A}}+\cos \theta_{\mathrm{B}}\right)+1 \geq 0 .
$$

Also, it can be written as

$$
-r_{2 \text { cylinders }}^{2} \cos \left(\theta_{A}-\theta_{B}\right)+\left(1-r_{2 \text { cylinders }} \cos \theta_{A}\right)\left(1-r_{2 \text { cylinders }} \cos \theta_{B}\right) \geq 0
$$

By taking $\theta_{A}=\theta_{B}=0$, we have $r_{2 c y l i n d e r s} \leq \frac{1}{2}$. Therefore, all the allowed measurements and all the inputs have positive outcomes if and only if $r_{2 c y l i n d e r s} \leq \frac{1}{2}$.

By taking $r_{2 c y l i n d e r s}=r \lambda^{3}$, then $r$ will be less than or equal to $\frac{1}{2 \lambda^{3}} \simeq 0.0573$.
We find that $r$ of the block that contains only one particle is less than $r$ of the block of two particles, $\frac{1}{\lambda^{4}} \simeq 0.0557 \leq \frac{1}{2 \lambda^{3}} \simeq 0.0573$. Therefore, the value of $r$ increases as long as it satisfies the dual constraints.

We will show under certain assumptions which we believe to be true how we compute the maximum value of $r$ for a block state space that contains only four state spaces. These assumptions are:

1. the measurement projector is $\left(\frac{I-X}{2}\right)$,
2. instead of considering all possible inputs, we take the input ( $I, r X, 0 Y, 1 Z$ )

### 4.2 Block state space of four cylindrical state spaces

In this section, we are going to combine four cylindrical state spaces into the block state space and apply four internal $C Z$ gates, to work out the maximum radius of the block $r_{\text {pos }}$, such that the block state space belongs to the dual of cylinder measurement. We will find this maximum value by computing the $\operatorname{trace}(\rho P)$, where $\rho$ is the set of state spaces and $P$ is the set of the allowed measurements.

To work out $r_{p o s}$, we assume that the Bloch vector for each particle is $(r, 0,1)$ and all the angles are equal to zero. Then, without applying $C Z, \rho^{(1,2,3,4)}$ is equal to:

$$
\begin{equation*}
\rho^{(1,2,3,4)}=\frac{1}{2}(I+r X+Z) \otimes \frac{1}{2}(I+r X+Z) \otimes \frac{1}{2}(I+r X+Z) \otimes \frac{1}{2}(I+r X+Z) \tag{4.1}
\end{equation*}
$$

After applying the $C Z_{1}$ on the first and the second qubits, $C Z_{2}$ on the second and third qubits, $C Z_{3}$ on the third and the fourth qubits and finally $C Z_{4}$ on the fourth and first qubits, we obtained the following $\rho^{(1,2,3,4)^{\prime}}$ :

$$
\rho^{(1,2,3,4)^{\prime}}=\frac{1}{16}\left[I \otimes I \otimes I \otimes I+r Z \otimes I \otimes Z \otimes X+I \otimes I \otimes I \otimes Z+r I \otimes Z \otimes X \otimes Z+r^{2} Z \otimes Z \otimes\right.
$$



Figure 4.1: The block state space that consists of four cylinder state spaces with four $C Z$ gates
$Y \otimes Y+r I \otimes Z \otimes X \otimes I+I \otimes I \otimes Z \otimes I+r Z \otimes I \otimes I \otimes X+I \otimes I \otimes Z \otimes Z+r Z \otimes X \otimes Z \otimes I+r^{2} I \otimes$ $X \otimes I \otimes X+r Z \otimes X \otimes Z \otimes Z+r^{2} Z \otimes Y \otimes Y \otimes Z+r^{3} I \otimes Y \otimes X \otimes Y+r^{2} Z \otimes Y \otimes Y \otimes I+r Z \otimes X \otimes$ $I \otimes I+r^{2} I \otimes X \otimes Z \otimes X+r Z \otimes X \otimes I \otimes Z+I \otimes Z \otimes I \otimes I+r Z \otimes Z \otimes Z \otimes X+I \otimes Z \otimes I \otimes Z+r I \otimes$ $I \otimes X \otimes Z+r^{2} Z \otimes I \otimes Y \otimes Y+r I \otimes I \otimes X \otimes I+I \otimes Z \otimes Z \otimes I+r Z \otimes Z \otimes I \otimes X+I \otimes Z \otimes Z \otimes Z+$ $r X \otimes Z \otimes I \otimes Z+r^{2} Y \otimes Z \otimes Z \otimes Y+r X \otimes Z \otimes I \otimes I+r^{2} X \otimes I \otimes X \otimes I+r^{3} Y \otimes I \otimes Y \otimes X+r^{2} X \otimes$ $I \otimes X \otimes Z+r X \otimes Z \otimes Z \otimes Z+r^{2} Y \otimes Z \otimes I \otimes Y+r X \otimes Z \otimes Z \otimes I+r^{2} Y \otimes Y \otimes Z \otimes Z+r^{3} X \otimes Y \otimes$ $I \otimes Y+r^{2} Y \otimes Y \otimes Z \otimes I+r^{3} Y \otimes X \otimes Y \otimes I+r^{4} X \otimes X \otimes X \otimes X+r^{3} Y \otimes X \otimes Y \otimes Z+r^{2} Y \otimes Y \otimes$ $I \otimes Z+r^{3} X \otimes Y \otimes Z \otimes Y+r^{2} Y \otimes Y \otimes I \otimes I+r X \otimes I \otimes I \otimes Z+r^{2} Y \otimes I \otimes Z \otimes Y+r X \otimes I \otimes I \otimes I+$ $r^{2} X \otimes Z \otimes X \otimes I+r^{3} Y \otimes Z \otimes Y \otimes X+r^{2} X \otimes Z \otimes X \otimes Z+r X \otimes I \otimes Z \otimes Z+r^{2} Y \otimes I \otimes I \otimes Y+$ $r X \otimes I \otimes Z \otimes I+Z \otimes I \otimes I \otimes I+r I \otimes I \otimes Z \otimes X+Z \otimes I \otimes I \otimes Z+r Z \otimes Z \otimes X \otimes Z+r^{2} I \otimes Z \otimes Y \otimes$ $Y+r Z \otimes Z \otimes X \otimes I+Z \otimes I \otimes Z \otimes I+r I \otimes I \otimes I \otimes X+Z \otimes I \otimes Z \otimes Z+r I \otimes X \otimes Z \otimes I+r^{2} Z \otimes X \otimes$ $I \otimes X+r I \otimes X \otimes Z \otimes Z+r^{2} I \otimes Y \otimes Y \otimes Z+r^{3} Z \otimes Y \otimes X \otimes Y+r^{2} I \otimes Y \otimes Y \otimes I+r I \otimes X \otimes I \otimes$ $I+r^{2} Z \otimes X \otimes Z \otimes X+r I \otimes X \otimes I \otimes Z+Z \otimes Z \otimes I \otimes I+r I \otimes Z \otimes Z \otimes X+Z \otimes Z \otimes I \otimes Z+r Z \otimes$ $\left.I \otimes X \otimes Z+r^{2} I \otimes I \otimes Y \otimes Y+r Z \otimes I \otimes X \otimes I+Z \otimes Z \otimes Z \otimes I+r I \otimes Z \otimes I \otimes X+Z \otimes Z \otimes Z \otimes Z\right]$.
$r_{\text {pos }}$ can be calculated by studying the trace, that should be positive, of $\rho^{(1,2,3,4)^{\prime}}$
and the measurement projector $P$ :

$$
\begin{equation*}
\operatorname{tr}\left(\rho^{(1,2,3,4)^{\prime}} P\right) \geq 0 \tag{4.2}
\end{equation*}
$$

where

$$
\begin{equation*}
P=\left[\frac{1}{2}(I-X)\right]^{\otimes 4} \tag{4.3}
\end{equation*}
$$

the computation of trace is given by

$$
\begin{equation*}
\operatorname{tr}\left(\rho^{(1,2,3,4)^{\prime}} P\right)=\frac{1}{16}\left(r_{p o s}^{4}+2 r_{p o s}^{2}-4 r_{p o s}+1\right) \geq 0 \tag{4.4}
\end{equation*}
$$

The maximum value of $r_{p o s}$ satisfying the equation of trace is:

$$
r_{p o s} \leq 0.295597743 \simeq 0.2956
$$



Figure 4.2: maximum radius required in the block state space to be separable.

In this block, we can only apply two external $C Z$ gates which means that $r$ will be taken to $r \lambda^{2}$ to maintain the separable decomposition. Hence, $r \lambda^{2} \leq 0.2956$, then $r$ is given by

$$
r \leq \frac{0.2956}{\lambda^{2}}=\frac{0.2956}{2.05817^{2}} \approx 0.0698
$$

The value 0.0698 is computed based on the assumptions that the negativity will first appear for measurement $\frac{(I-X)}{2}$.

It turns out that when we add more vertices in the block state space, the value of $r$ increases as long as it satisfies the dual constraints to maintain separability [7].

### 4.3 Block state space of many cylindrical state spaces

This section is a review of further progress which was made by my co-authors and described in article [7]. We are going to increase the size of the block state space and we will obtain two sequences (upper and lower) of optimization problems. The lower sequence is non-decreasing while the upper sequence is non-increasing. The lower sequence represents the values of the input radius where the system can be efficiently simulated classically and the upper sequence represents the input radius where the system gives negative probabilities. In the case of $2 D$ lattice, these two sequences converge to limits that are close but we do not know if both sequences have the same limits, i.e. if they have the same limit then below this specific value the system can be efficiently simulated classically and above it the system has negative probabilities.

We consider a block B that has $H \times W$ qubits, where $H$ and $W$ are greater than 2. To initialise the qubits, we consider two ways:

1. All cylinders have the same state $(1, r, 0,1)$. After applying internal $C Z$, the resulting operator is denoted by $\rho(B, r)$.
2. All cylinders that are in the corner of the block, are prepared in a state with $\lambda^{2} r$ and the boundary cylinders are in a state with $\lambda r$, while the remaining cylinders are in a state with input radius $r$. All these states have $z=1$. After applying the internal $C Z$, the resulting operator is denoted by $\rho_{\lambda}(B, r)$.

We will also define the following:

$$
\begin{aligned}
s(B) & :=\max \left\{r \mid \rho(B, r) \geq_{M} 0\right\} \\
s_{\lambda}(B) & :=\max \left\{r \mid \rho_{\lambda}(B, r) \geq_{M} 0\right\}
\end{aligned}
$$

where $\rho(B, r) \geq_{M} 0$ and $\rho_{\lambda}(B, r) \geq_{M} 0$ means that these operators are in the dual of the set $M$ of the allowed measurements.

We are interested in $s_{\lambda}(B)$ to implement the coarse graining scheme. If $r \leq s_{\lambda}(B)$ in a block $B$, then cylinders with radius $r$ can be efficiently simulated classically because they lead to positive probabilities. But if $r \geq s(B)$ in a block $B$, then cylinders with radius $r$ lead to negative probabilities which mean that we won't end up with separable decomposition. So, we will try to find blocks where $s_{\lambda}(B)$ is large.

Lemma 3. Let's consider a region $K L$ of qubits in a lattice and we divide this region into two disjoint subregions $K$ and $L$. By disjoints we mean that we remove all the $C Z$ gates that relate these two regions. For any region $F$, we have

$$
\begin{equation*}
s(F) \geq s_{\lambda}(F) \tag{4.5}
\end{equation*}
$$

Then for the region $K L$, we have:

$$
\begin{gather*}
s(K L) \leq \min \{s(K), s(L)\}  \tag{4.6}\\
s_{\lambda}(K L) \geq \min \left\{s_{\lambda}(K), s_{\lambda}(L)\right\} \tag{4.7}
\end{gather*}
$$

If we consider that the two regions are identical, $K=L$, then

$$
\begin{gather*}
s(L L) \leq s(L)  \tag{4.8}\\
s_{\lambda}(L L) \geq s_{\lambda}(L) \tag{4.9}
\end{gather*}
$$

This tells us that every time we join the subregions together, we increase the size of a region, then $s$ will only decrease and $s_{\lambda}$ will only increase.

Proof. To end up with equation (4.5), we start by obtaining $\rho(F, r)$ from $\rho_{\lambda}(F, r)$ by dephasing the edges of the corner of qubits, i.e. we need to remove $\lambda^{2}$ and $\lambda$. Hence, if $\rho_{\lambda}(F, r) \geq 0$ then $\rho(F, r) \geq 0$. When $\rho(F, r)$ becomes negative, the value of $r$ has to be greater than the value of $r$ when $\rho_{\lambda}(F, r)$ becomes negative. So, as dephasing maintains positivity, $s(F) \geq s_{\lambda}(F)$.

To end up with equation (4.6), we begin by considering a cylinder with $z=1$ that has the form $|0\rangle\langle 0|+a|0\rangle\langle 1|+b|1\rangle\langle 0|$ where $a, b \in \mathbb{C}$. Now suppose we interact this cylinder with two other particles in an arbitrary state $T$ using two $C Z$. Then we have

$$
|0\rangle\langle 0| \otimes T+a|0\rangle\langle 1| \otimes\left(T Z^{\otimes 2}\right)+b|1\rangle\langle 0| \otimes\left(Z^{\otimes 2} T\right)
$$

If we trace out the first particle, then the remaining particles are in their original marginal state $T$. So, cutting the external $C Z$ gates that join the two regions $K$ and $L$, does not change the marginal operators on $K$ or $L, \operatorname{tr}_{K}\{\rho(K L, r)\}=\rho(L, r)$. If... $\rho(K L, r) \geq_{M} 0$ then we have $\operatorname{tr}_{K}\{\rho(K L, r)\} \geq_{M} 0$, but as $\operatorname{tr}_{K}\{\rho(K L, r)\}=$ $\rho(L, r)$ this means that $\rho(K L, r) \geq_{M} 0$ implies $\rho(L, r) \geq_{M} 0$ and $\rho(K, r) \geq_{M} 0$. So by taking the minimum between $s(L)$ and $s(K), s(K L)$ must be less than or equal to the minimum to have $\rho(K L, r) \geq_{M} 0$.

To end up with equation (4.7), we have $\rho_{\text {lambda }}(K L, r)$ is in the convex hull of
$\rho_{\lambda}(K, r) \otimes \rho_{\lambda}(L, r)$. Hence $\rho_{\lambda}(K L, r)$ on the block must be positive if $\rho_{\lambda}(K, r) \geq_{M} 0$ and $\rho_{\lambda}(L, r) \geq_{M} 0$ and so the equation $s_{\lambda}(K L) \geq \min \left\{s_{\lambda}(K), s_{\lambda}(L)\right\}$ must hold.

This lemma allows us to define sequences that capture when $r$ is classically simulated classically using the coarse graining approach. Let's consider that we construct a sequence of blocks and we start with a $2 \times 2$ block $B_{1}$ and then construct larger blocks by joining two copies of $B_{n-1}$ to make $B_{n}$. We define the sequences as follows

$$
\begin{align*}
& u_{n}:=s\left(B_{n}\right)  \tag{4.10}\\
& l_{n}:=s_{\lambda}\left(B_{n}\right) \tag{4.11}
\end{align*}
$$

From equation (4.5), we have $u_{n} \geq l_{n}$ where $l_{n}$ is non-decreasing and $u_{n}$ is nonincreasing and hence both sequences converge. We denote the limits as:

$$
\begin{aligned}
u & :=\lim u_{n} \\
l & :=\lim l_{n}
\end{aligned}
$$

If $r<l$ then $r$ is classically simulatable efficiently but if $r>u$ then we end up with negative probabilities. If it turns out that $u=l$ then for $r<l$ the system is classically simulated and for $r>l$ the system has negative probabilities. But we do not know whether or not $u$ is equal to $l$ for any system. For a $2 D$ lattice with the interaction of $C Z$ gates, we calculated the lower and the upper bounds using trial measurements and inputs on rectangles of size $6 \times 7$. These numbers indicates that $0.0698 \leq l \leq u \leq 0.139$ for a $2 D$ square lattice. For this specific lattice, the numerical experiments seem to suggest these two assumptions: the upper and lower sequences are determined by taking a measurement projector that has the form $\frac{I-X}{2}$ and the input that has the form $(1, \alpha, 0,1)$ in the Pauli basis where $\alpha$ includes $r$ and any growth factors $\lambda$ applied in the coarse graining approach. When these inputs and measurements give positive probabilities, the maximum $r$ appears to be the maximum in the following equations $s(B)=\max \left\{r \mid \rho(B, r) \geq_{M} 0\right\}$.

### 4.4 Summary of chapter 4

By using the method of coarse graining, we tried to expand the size of the region of the system that can be efficiently simulated classically. It turns out that the value of the radius increases slightly for two particles. But for four particles in a block, the system can be efficiently classically simulated under the following two assumptions:
the measurement projector is $\left(\frac{I-X}{2}\right)^{\otimes 4}$ and the input is $(1, r, 0,1)$ in the Pauli basis. We use this measurement projector because it seems for numerical experiments that the negativity appears first for this measurement, i.e. if, by applying this measurement, we end up with positivity then the remaining allowed measurements will give positive values. In [7], we found that by adding more vertices in a block, the value of the radius increases, however it doesn't exceed 1.

## Chapter 5

## Summary

It is unknown if quantum computers are better than classical computers. We can proceed one of two ways to understand this problem. The first way is that we can try to find that quantum computers cannot be efficiently simulated classically and the second is that we can try to develop an algorithm that can classically simulate a quantum system.

It is believed that quantum computers, built with quantum system, can efficiently solve hard problems that are unknown whether they can be efficiently solved on classical computers. If it is found that simulating a quantum computing can be done classically, then any algorithm that can be done in quantum computers can be efficiently simulated on a classical computer.

In this thesis, our goal is to develop a new method of classically simulating certain types of quantum system that are previously unknown to be efficiently simulatable on classical computers.

We started by defining cylinder separability which is a new generalised notion of separability. Then we applied this notion to systems where we vary the inputs of cluster state quantum computation and we end up with some inputs that allow the system to be efficiently simulated classically. As we mentioned previously, the cluster state has input states $|\psi\rangle=|+\rangle=\frac{|0\rangle+|1\rangle}{\sqrt{2}}$, but we have shown that if we change $|\psi\rangle$ into the unormalised form $|0\rangle+\epsilon|1\rangle$ where $\epsilon$ is a small non-zero number, the system can be efficiently simulated classically.

The classical simulation is achieved by demonstrating that if we take $C Z$ gate acting on two inputs that are presented by two cylindrical state spaces with radii $r$, the output is separable with respect to two different cylindrical state spaces with larger radii $R=\lambda r$ where $\lambda=\sqrt{2+\sqrt{5}}$. We can maintain a separable decomposition by growing the cylinder by $\lambda$ each time but we must not grow beyond radius 1 , because this state space should be in the dual of the cylinder measurements. This work forms the basis of the classical simulation.

We also considered, instead of $C Z$ gates, diagonal gates $V(\theta)$ that acts on two cylindrical state spaces. Our goal is to try to find the minimum growth required for separability. After working out the required growth factor as a function of $\theta$ for all the diagonal gates, we discover that the $C Z$ gate is the gate that requires the most growth factor $\lambda$ to maintain separability. This work can be taken further. Instead of using $V(\theta)$ gates of two qubits, we could try to use different diagonal gates such as $|0\rangle\langle 0| \otimes z_{\theta}+|1\rangle\langle 1| \otimes z_{\theta}$ where $z_{\theta}=|0\rangle\langle 0|+e^{i \theta}|1\rangle\langle 1|$, or we could try other gates with number of qudits to see whether or not it is possible to efficiently simulated on a classical computer.

Then, we tried to find if there are other state spaces require a lower growth factor. So, we considered new state spaces that are the convex hull of $N(5-2)+2$ vertices. We chose these state spaces because of the symmetry about the $z$ axis and they are among the simplest we could try for small number of extremal points. Then we grew the size of these state spaces to maintain separability and we used Linear programming to compute the required growth factors to see if any of the state spaces that we studied can grow slower than the cylinder. But it turns out that, among all these state spaces that we tried, the cylinder is the most optimal state space that needs the least growth. In further work, we could try different state spaces rather than the ones we studied but these state spaces have to be also symmetric about the $z$ axis.

Finally, we introduce the method of coarse graining and we try to improve the value of the radius of a system that can be efficiently simulated on a classical computer using this specific method. Then we find that the radius that can be simulated efficiently, increases slightly for two particles and then even further subject to the following two assumptions for four particles in a block: the measurement projector has the form $\frac{I-X}{2}$ and instead of taking all the possible inputs, we consider the input $(I, r X, 0 Y, Z)$. If we end up with positivity using this measurement projector then the remaining allowed measurements will give positive values. We will also find that by adding more vertices in a block, the value of the radius will increase slightly [7]. Two open problems are whether the upper and lower sequences converge and whether there are other ways of coarse graining that can improve the classical simulation. These open problems need further investigations.

## Appendix A

## Representation of $\mathbf{V}(\theta)$ in the Pauli basis

In this part, we will compute the representation of $V(\theta)$ in Pauli basis that we called $V V(\theta)$. It is a $16 \times 16$ matrix and is used in chapter 2 .

First, we take two inputs in the following Pauli basis:

$$
(I, X, Y, Z)
$$

where we have:

$$
I=\left(\begin{array}{cc}
1 & 0 \\
0 & 1
\end{array}\right) ; X=\left(\begin{array}{cc}
0 & 1 \\
1 & 0
\end{array}\right) ; Y=\left(\begin{array}{cc}
0 & -i \\
i & 0
\end{array}\right) I \text { and } Z=\left(\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right)
$$

Then, we act the Controlled gate $V(\theta)$ on these two qubits by taking one Pauli operator from the first and second input; where $V(\theta)$ and $V^{\dagger}(\theta)$ are represented as the following:
$V(\theta)=\left(\begin{array}{cccc}1 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & e^{i \theta}\end{array}\right)$ and $V^{\dagger}(\theta)=\left(\begin{array}{cccc}1 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & e^{-i \theta}\end{array}\right)$

- Let's start by taking $I$ from the first input and changing the Pauli operators of the second input.

$$
\begin{aligned}
V(\theta)(I \otimes I) V^{\dagger}(\theta) & =\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(i \theta)}
\end{array}\right)\left(\begin{array}{llll}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & 1
\end{array}\right)\left(\begin{array}{lllc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(-i \theta)}
\end{array}\right) \\
& =\left(\begin{array}{llll}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & 1
\end{array}\right) \\
& =I \otimes I .
\end{aligned}
$$

- $V(\theta)$ is acting on $I$ and $X$ where $X$ is the second input considered.

$$
\begin{aligned}
V(\theta)(I \otimes X) V^{\dagger}(\theta) & =\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(i \theta)}
\end{array}\right)\left(\begin{array}{llll}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 0 & 1 \\
0 & 0 & 1 & 0
\end{array}\right)\left(\begin{array}{lllc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(-i \theta)}
\end{array}\right) \\
& =\left(\begin{array}{cccc}
0 & 1 & 0 & 0 \\
1 & 0 & 0 & 0 \\
0 & 0 & 0 & e^{-i \theta} \\
0 & 0 & e^{i \theta} & 0
\end{array}\right) \\
& =\left(\frac{I+Z}{2}\right) \otimes X+\left(\frac{I-Z}{2}\right) \otimes(\cos \theta X+\sin \theta Y) \\
& =a I \otimes X+b I \otimes Y+c Z \otimes X+d Z \otimes Y
\end{aligned}
$$

- Let's take $Y$ the second input.

$$
\begin{aligned}
V(\theta)(I \otimes Y) V^{\dagger}(\theta) & =\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(i \theta)}
\end{array}\right)\left(\begin{array}{cccc}
0 & -i & 0 & 0 \\
i & 0 & 0 & 0 \\
0 & 0 & 0 & -i \\
0 & 0 & i & 0
\end{array}\right)\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(-i \theta)}
\end{array}\right) \\
& =\left(\begin{array}{cccc}
0 & -i & 0 & 0 \\
i & 0 & 0 & 0 \\
0 & 0 & 0 & -i e^{-i \theta} \\
0 & 0 & i e^{i \theta} & 0
\end{array}\right) \\
& =\left(\frac{I+Z}{2}\right) \otimes Y+\left(\frac{I-Z}{2}\right) \otimes(\cos \theta Y-\sin \theta X) \\
& =a I \otimes Y+d I \otimes X+b Z \otimes X+c Z \otimes Y
\end{aligned}
$$

- The second input considered is $Z$.

$$
\begin{aligned}
V(\theta)(I \otimes Z) V^{\dagger}(\theta) & =\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(i \theta)}
\end{array}\right)\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & -1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & -1
\end{array}\right)\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(-i \theta)}
\end{array}\right) \\
& =\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & -1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & -e^{-i \theta} e^{i \theta}
\end{array}\right) \\
& =I \otimes Z
\end{aligned}
$$

Instead of $I$, we consider $X$ as the first input and change the second input.

- The second input is $I$.

$$
\begin{aligned}
V(\theta)(X \otimes I) V^{\dagger}(\theta) & =\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(i \theta)}
\end{array}\right)\left(\begin{array}{llll}
0 & 0 & 1 & 0 \\
0 & 0 & 0 & 1 \\
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0
\end{array}\right)\left(\begin{array}{lllc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(-i \theta)}
\end{array}\right) \\
& =\left(\begin{array}{cccc}
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{-i \theta} \\
1 & 0 & 0 & 0 \\
0 & e^{i \theta} & 0 & 0
\end{array}\right) \\
& =\left(\frac{X+i Y}{2}\right) \otimes\left(\frac{I+Z}{2}+\frac{I-Z}{2} e^{-i \theta}\right)+\left(\frac{X-i Y}{2}\right) \otimes\left(\frac{I+Z}{2}+\frac{I-Z}{2} e^{i \theta}\right) \\
& =e X \otimes I+j X \otimes Z+p Y \otimes I+q Y \otimes Z
\end{aligned}
$$

- $X$ is considered as the second input.

$$
\begin{aligned}
V(\theta)(X \otimes X) V^{\dagger}(\theta) & =\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(i \theta)}
\end{array}\right)\left(\begin{array}{llll}
0 & 0 & 0 & 1 \\
0 & 0 & 1 & 0 \\
0 & 1 & 0 & 0 \\
1 & 0 & 0 & 0
\end{array}\right)\left(\begin{array}{llll}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(-i \theta)}
\end{array}\right) \\
& =\left(\begin{array}{cccc}
0 & 0 & 0 & e^{-i \theta} \\
0 & 0 & 1 & 0 \\
0 & 1 & 0 & 0 \\
e^{i \theta} & 0 & 0 & 0
\end{array}\right) \\
& =\frac{X+i Y}{2} \otimes\left(\frac{X-i Y}{2}+\frac{X+i Y}{2} e^{-i \theta}\right)+\frac{X-i Y}{2} \otimes\left(\frac{X+i Y}{2}+\frac{X-i Y}{2} e^{i \theta}\right) \\
& =e X \otimes X+p X \otimes Y+p Y \otimes X+j Y \otimes Y
\end{aligned}
$$

- The second input taken is $Y$.

$$
\begin{aligned}
V(\theta)(X \otimes Y) V^{\dagger}(\theta) & =\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(i \theta)}
\end{array}\right)\left(\begin{array}{cccc}
0 & 0 & 0 & -i \\
0 & 0 & i & 0 \\
0 & -i & 0 & 0 \\
i & 0 & 0 & 0
\end{array}\right)\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(-i \theta)}
\end{array}\right) \\
& =\left(\begin{array}{cccc}
0 & 0 & 0 & -i e^{-i \theta} \\
0 & 0 & i & 0 \\
0 & -i & 0 & 0 \\
i e^{i \theta} & 0 & 0 & 0
\end{array}\right) \\
& =\frac{X+i Y}{2} \otimes\left(\frac{i X+Y}{2}+\frac{-i X+Y}{2} e^{-i \theta}\right)+\frac{X-i Y}{2} \otimes\left(\frac{-i X+Y}{2}+\frac{i X+Y}{2} e^{i}\right. \\
& =q X \otimes X+e X \otimes Y+k Y \otimes X+p Y \otimes Y
\end{aligned}
$$

- The last operator is $Z$ which is considered for the second input.

$$
\begin{aligned}
V(\theta)(X \otimes Z) V^{\dagger}(\theta) & =\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(i \theta)}
\end{array}\right)\left(\begin{array}{cccc}
0 & 0 & 1 & 0 \\
0 & 0 & 0 & -1 \\
1 & 0 & 0 & 0 \\
0 & -1 & 0 & 0
\end{array}\right)\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(-i \theta)}
\end{array}\right) \\
& =\left(\begin{array}{cccc}
0 & 0 & 1 & 0 \\
0 & 0 & 0 & -e^{-i \theta} \\
1 & 0 & 0 & 0 \\
0 & -e^{i \theta} & 0 & 0
\end{array}\right) \\
& =\left(\frac{X+i Y}{2}\right) \otimes\left(\frac{I+Z}{2}-\frac{I-Z}{2} e^{-i \theta}\right)+\left(\frac{X-i Y}{2}\right) \otimes\left(\frac{I+Z}{2}-\frac{I-Z}{2} e^{i \theta}\right) \\
& =j X \otimes I+q Y \otimes I+e X \otimes Z+p Y \otimes Z
\end{aligned}
$$

The third operator is $Y$ which is taken as the first unchangeable input.

- Let's start by taking $I$ the second input.

$$
\begin{aligned}
V(\theta)(Y \otimes I) V^{\dagger}(\theta) & =\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(i \theta)}
\end{array}\right)\left(\begin{array}{cccc}
0 & 0 & -i & 0 \\
0 & 0 & 0 & -i \\
i & 0 & 0 & 0 \\
0 & i & 0 & 0
\end{array}\right)\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(-i \theta)}
\end{array}\right) \\
& =\left(\begin{array}{cccc}
0 & 0 & -i & 0 \\
0 & 0 & 0 & -i e^{-i \theta} \\
i & 0 & 0 & 0 \\
0 & i e^{i \theta} & 0 & 0
\end{array}\right) \\
& =q X \otimes I+p X \otimes Z+e Y \otimes I+j Y \otimes Z
\end{aligned}
$$

- The second input is $X$.

$$
\begin{aligned}
V(\theta)(Y \otimes X) V^{\dagger}(\theta) & =\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(i \theta)}
\end{array}\right)\left(\begin{array}{cccc}
0 & 0 & 0 & -i \\
0 & 0 & -i & 0 \\
0 & i & 0 & 0 \\
i & 0 & 0 & 0
\end{array}\right)\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(-i \theta)}
\end{array}\right) \\
& =\left(\begin{array}{cccc}
0 & 0 & 0 & -i e^{-i \theta} \\
0 & 0 & -i & 0 \\
0 & i & 0 & 0 \\
i e^{i \theta} & 0 & 0 & 0
\end{array}\right) \\
& =q X \otimes X+k X \otimes Y+e Y \otimes X+p Y \otimes Y
\end{aligned}
$$

- Let's take $Y$ the second input.

$$
\begin{aligned}
V(\theta)(Y \otimes Y) V^{\dagger}(\theta) & =\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(i \theta)}
\end{array}\right)\left(\begin{array}{cccc}
0 & 0 & 0 & -1 \\
0 & 0 & 1 & 0 \\
0 & 1 & 0 & 0 \\
-1 & 0 & 0 & 0
\end{array}\right)\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(-i \theta)}
\end{array}\right) \\
& =\left(\begin{array}{cccc}
0 & 0 & 0 & -e^{-i \theta} \\
0 & 0 & 1 & 0 \\
0 & 1 & 0 & 0 \\
-e^{i \theta} & 0 & 0 & 0
\end{array}\right) \\
& =j X \otimes X+q X \otimes Y+q Y \otimes X+e Y \otimes Y
\end{aligned}
$$

- $Z$ is considered as the second input.

$$
\begin{aligned}
V(\theta)(Y \otimes Z) V^{\dagger}(\theta) & =\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(i \theta)}
\end{array}\right)\left(\begin{array}{cccc}
0 & 0 & -i & 0 \\
0 & 0 & 0 & i \\
i & 0 & 0 & 0 \\
0 & -i & 0 & 0
\end{array}\right)\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(-i \theta)}
\end{array}\right) \\
& =\left(\begin{array}{cccc}
0 & 0 & -i & 0 \\
0 & 0 & 0 & i e^{-i \theta} \\
i & 0 & 0 & 0 \\
0 & -i e^{i \theta} & 0 & 0
\end{array}\right) \\
& =p X \otimes I+q X \otimes Z+j Y \otimes I+e Y \otimes Z
\end{aligned}
$$

And finally, $V(\theta)$ is acting on the last operator $Z$ which is taken as the first input.

- The second input considered is $I$.

$$
\begin{aligned}
V(\theta)(Z \otimes I) V^{\dagger}(\theta) & =\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(i \theta)}
\end{array}\right)\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & -1 & 0 \\
0 & 0 & 0 & -1
\end{array}\right)\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(-i \theta)}
\end{array}\right) \\
& =\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & -1 & 0 \\
0 & 0 & 0 & -e^{i \theta} e^{-i \theta}
\end{array}\right) \\
& =Z \otimes I
\end{aligned}
$$

- $X$ is taken as the second input.

$$
\begin{aligned}
V(\theta)(Z \otimes X) V^{\dagger}(\theta) & =\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(i \theta)}
\end{array}\right)\left(\begin{array}{cccc}
0 & 1 & 0 & 0 \\
1 & 0 & 0 & 0 \\
0 & 0 & 0 & -1 \\
0 & 0 & -1 & 0
\end{array}\right)\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(-i \theta)}
\end{array}\right) \\
& =\left(\begin{array}{cccc}
0 & 1 & 0 & 0 \\
1 & 0 & 0 & 0 \\
0 & 0 & 0 & -e^{-i \theta} \\
0 & 0 & -e^{i \theta} & 0
\end{array}\right) \\
& =c I \otimes X+a Z \otimes X+d I \otimes Y+b Z \otimes Y
\end{aligned}
$$

- The second changeable input is $Y$.

$$
\begin{aligned}
V(\theta)(Z \otimes Y) V^{\dagger}(\theta) & =\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(i \theta)}
\end{array}\right)\left(\begin{array}{cccc}
0 & -i & 0 & 0 \\
i & 0 & 0 & 0 \\
0 & 0 & 0 & i \\
0 & 0 & -i & 0
\end{array}\right)\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(-i \theta)}
\end{array}\right) \\
& =\left(\begin{array}{cccc}
0 & -i & 0 & 0 \\
i & 0 & 0 & 0 \\
0 & 0 & 0 & i e^{-i \theta} \\
0 & 0 & -i e^{i \theta} & 0
\end{array}\right) \\
& =c I \otimes Y+a Z \otimes Y+b I \otimes X+d Z \otimes X
\end{aligned}
$$

- The last operator $Z$ is taken as the second input.

$$
\begin{aligned}
V(\theta)(Z \otimes Z) V^{\dagger}(\theta) & =\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(i \theta)}
\end{array}\right)\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & -1 & 0 & 0 \\
0 & 0 & -1 & 0 \\
0 & 0 & 0 & 1
\end{array}\right)\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & e^{(-i \theta)}
\end{array}\right) \\
& =\left(\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & -1 & 0 & 0 \\
0 & 0 & -1 & 0 \\
0 & 0 & 0 & e^{i \theta} e^{-i \theta}
\end{array}\right) \\
& =Z \otimes Z
\end{aligned}
$$

After applying $V(\theta)$ to these inputs, we get the $16 \times 16$ matrix, that we called $V V(\theta)$.

$$
\left(\begin{array}{llllllllllllllll}
1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & \mathrm{a} & \mathrm{~b} & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & \mathrm{c} & \mathrm{~d} & 0 \\
0 & \mathrm{~d} & \mathrm{a} & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & \mathrm{~b} & \mathrm{c} & 0 \\
0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & \mathrm{e} & 0 & 0 & \mathrm{j} & \mathrm{p} & 0 & 0 & \mathrm{q} & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & \mathrm{e} & \mathrm{p} & 0 & 0 & \mathrm{p} & \mathrm{j} & 0 & 0 & 0 & 0 & 0 \\
(I \otimes X) \\
0 & 0 & 0 & 0 & 0 & \mathrm{q} & \mathrm{e} & 0 & 0 & \mathrm{k} & \mathrm{p} & 0 & 0 & 0 & 0 & 0 \\
(I \otimes Y) \\
0 & 0 & 0 & 0 & \mathrm{j} & 0 & 0 & \mathrm{e} & \mathrm{q} & 0 & 0 & \mathrm{p} & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & \mathrm{q} & 0 & 0 & \mathrm{p} & \mathrm{e} & 0 & 0 & \mathrm{j} & 0 & 0 & 0 & 0 \\
(X \otimes I) \\
0 & 0 & 0 & 0 & 0 & \mathrm{q} & \mathrm{k} & 0 & 0 & \mathrm{e} & \mathrm{p} & 0 & 0 & 0 & 0 & 0 \\
(X \otimes Y) \\
0 & 0 & 0 & 0 & 0 & \mathrm{j} & \mathrm{q} & 0 & 0 & \mathrm{q} & \mathrm{e} & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & \mathrm{p} & 0 & 0 & \mathrm{q} & \mathrm{j} & 0 & 0 & \mathrm{e} & 0 & 0 & 0 & 0 \\
(Y \otimes Z) \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 1 & 0 & 0 & 0 \\
0 & \mathrm{c} & \mathrm{~d} & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & \mathrm{a} & \mathrm{~b} & 0 \\
0 & \mathrm{~b} & \mathrm{c} & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & \mathrm{~d} & \mathrm{a} & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 1
\end{array}\right)(Y \otimes Y)
$$

where

$$
a=\frac{1}{2}+\frac{1}{2} \cos \theta ; b=\frac{1}{2} \sin \theta ; c=\frac{1}{2}-\frac{1}{2} \cos \theta ; d=-\frac{1}{2} \sin \theta ;
$$

$e=\frac{1}{2}+\frac{1}{4} \exp (-i \theta)+\frac{1}{4} \exp (i \theta) ; j=\frac{1}{2}-\frac{1}{4} \exp (-i \theta)-\frac{1}{4} \exp (i \theta) ; k=-\frac{1}{2}+\frac{1}{4} \exp (-i \theta)+\frac{1}{4} \exp (i \theta) ;$

$$
p=\frac{1}{4} i \exp (-i \theta)-\frac{1}{4} i \exp (i \theta) ; q=-\frac{1}{4} i \exp (-i \theta)+\frac{1}{4} i \exp (i \theta)
$$

## Appendix B

## Matlab code

## B. 1 Matlab code for Chapter 2

When the gate $V(\theta)$ acts on cylindrical state spaces in chapter 2, we end up with identical minimum eigenvalues. We check this result using Matlab code:

```
    clear
%for the two inputs (1,rA,0,1)tensor (1,rB,0,1)
% pauli matrices
i = complex(0,1);
I = [1 0; 0 1];
X = [0 1; 1 0];
Y = [0 -i; i 0];
Z = [1 0; 0 -1];
rA = 1;
rB = 1;
j = 1;
array_min = zeros(5,361*numel(0:0.1:1.5));
%cs = cosd(theta)+ i*sind(theta);
%cs_m = cosd(theta) - i* sind(theta);
II = kron(I,I);%II
IX = kron(I,X);%IX
IY = kron(I,Y);%IY
IZ = kron(I,Z);%IZ
XI = kron(X,I);%XI
XX = kron(X,X);%XX
XY = kron(X,Y);%XY
```

```
XZ = kron(X,Z);%XZ
YI = kron(Y,I);%YI
YX = kron(Y,X);%YX
YY = kron(Y,Y);%YY
YZ = kron(Y,Z);%YZ
ZI = kron(Z,I);%ZI
ZX = kron(Z,X);%ZX
ZY = kron(Z,Y);%ZY
ZZ = kron(Z,Z);%ZZ
for K = 0:0.1:1.5
%add 1 to RA: going from 1 to 2.5 by 0.1
RA = 1+K;
RB = 1+K;
for theta = 0:360
cs = cosd(theta) + i*sind(theta);
cs_m = cosd(theta) - i* sind(theta);
```

\%this matrix is similar to equation 8 in the paper, where all the $Z$ components are 0 :
matrix_1 $=I I+(r B / R B) * I X+(r A / R A) * X I+$
$\left(\left(0.5+0.25 * c_{\text {_ }} \quad \mathrm{m}+0.25 * \mathrm{cs}\right) *((r A * r B) /(R A * R B))\right) * X X+\left(\left(-0.25 * i * c_{\text {_ }} m\right.\right.$
$+0.25 * \mathrm{i} * \mathrm{Cs}) *((\mathrm{rA} * \mathrm{rB}) /(\mathrm{RA} * \mathrm{RB}))) * \mathrm{XY}+\left(\left(-0.25 * i * \mathrm{cs} \_\mathrm{m}\right.\right.$
$+0.25 * i * c s) *((r A * r B) /(R A * R B))) * Y X$
$+\left(\left(0.5-0.25 * c s \_m-0.25 * c s\right) *((r A * r B) /(R A * R B))\right) * Y Y$;
\%this matrix is the partial transpose:
$\% 3 r d$ column is minus
matrix_2 $=I I+(r B / R B) * I X+(r A / R A) * X I+$
$\left(\left(0.5+0.25 * c_{\text {_ }} \mathrm{m}+0.25 * \mathrm{cs}\right) *((r A * r B) /(R A * R B))\right) * X X-\left(\left(-0.25 * i * c s \_m\right.\right.$
$+0.25 * i * \operatorname{cs}) *((r A * r B) /(R A * R B))) * X Y+\left(\left(-0.25 * i * c s \_m\right.\right.$
$+0.25 * i * C s) *((r A * r B) /(R A * R B))) * Y X$
$-\left(\left(0.5-0.25 * c s \_m-0.25 * c s\right) *((r A * r B) /(R A * R B))\right) * Y Y$;
\%or 3rd row is minus
matrix_3 $=I I+(r B / R B) * I X+(r A / R A) * X I+$
$\left(\left(0.5+0.25 * c_{s} \_m+0.25 * c s\right) *((r A * r B) /(R A * R B))\right) * X X+\left(\left(-0.25 * i * c s \_m\right.\right.$
$+0.25 * \mathrm{i} * \mathrm{cs}) *((\mathrm{rA} * \mathrm{rB}) /(\mathrm{RA} * \mathrm{RB}))) * \mathrm{XY}-\left(\left(-0.25 * \mathrm{i} * \mathrm{cs} \_\mathrm{m}\right.\right.$

```
+0.25*i*cs)*((rA*rB)/(RA*RB)))*YX
-((0.5-0.25*cs_m-0.25*cs)*((rA*rB)/(RA*RB)))*YY ;
%min eigenvalues of matrix_1
eigenvalues_1 = eig(matrix_1);
min_eig_1 = min(eigenvalues_1);
eigenvalues_2 = eig(matrix_2);
min_eig_2 = min(eigenvalues_2);
eigenvalues_3 = eig(matrix_3);
min_eig_3 = min(eigenvalues_3);
array_min(1,j) = theta;
array_min(2,j) = min_eig_1;
array_min}(3,j) = min_eig_2
array_min(4,j) = min_eig_3;
array_min(5,j) = RA;
j = j+1;
end
end
```

The graph (2.2) which represents the required growth for each theta, is given by the following code:

```
clear
%for the two inputs (1,rA,0,1)tensor (1,rB,0,1)
% pauli matrices
i = complex(0,1);
I = [1 0; 0 1];
X = [0 1; 1 0];
Y = [0 -i; i 0];
Z = [1 0; 0 -1];
rA = 1;
rB = 1;
j = 1;
```

array_R = zeros(1,numel(0:0.1:1.1 ));
array_theta $=$ zeros(1, numel(0:0.1:1.1 ));
$\% c s=\operatorname{cosd}($ theta $)+i *$ sind (theta) $;$
\%cs_m = cosd(theta) - i* sind(theta);
II $=k r o n(I, I) ; \% I I$
IX = kron(I,X); \%IX
IY $=k r o n(I, Y) ; \% I Y$
IZ = kron(I,Z);\%IZ
XI = kron(X,I);\%XI
XX $=$ kron( $\mathrm{X}, \mathrm{X}$ ); \% XX
$X Y=k r o n(X, Y) ; \% X Y$
XZ $=$ kron( $\mathrm{X}, \mathrm{Z}$ );\%XZ
YI $=k r o n(Y, I) ; \% Y I$
$Y X=k r o n(Y, X) ; \% Y X$
YY = kron(Y,Y);\%YY
YZ = kron(Y,Z);\%YZ
ZI = kron(Z,I);\%ZI
ZX = kron(Z,X);\%ZX
ZY = kron(Z,Y);\%ZY
ZZ = kron(Z,Z);\%ZZ

```
for K = 0:0.00001:1.1
%add 1 to RA: going from 1 to 2.1 by 0.00001
RA = 1+K;
RB = 1+K;
for theta = 0:360
cs = cosd(theta)+ i*sind(theta);
cs_m = cosd(theta) - i* sind(theta);
%this matrix is similar to equation 8 in the paper, where all the Z components are 0:
matrix_1 = II + (rB/RB)*IX + (rA/RA)* XI +
((0.5+0.25*cs_m+0.25*cs)*((rA*rB)/(RA*RB)))* XX + ((-0.25*i*cs_m
+0.25*i*cs)*((rA*rB)/(RA*RB)))*XY + ((-0.25*i*cs_m
+0.25*i*cs)*((rA*rB)/(RA*RB)))*YX
+((0.5-0.25*cs_m-0.25*cs)*((rA*rB)/(RA*RB)))*YY ;
%min eigenvalues of matrix_1
eigenvalues_1 = eig(matrix_1);
min_eig_1 = min(eigenvalues_1);
if min_eig_1 < 0
    array_R(1,j) = RA;
    array_theta(1,j) = theta;
else
        continue
end
if min_eig_1 < 0
    break
else
    continue
end
end
j = j + 1;
end
```

```
plot(array_theta, array_R)
%title('Required growth to be separable')
xlabel('values of theta')
ylabel('values of growth')
```


## B. 2 Matlab code for Chapter 3

The Matlab code, which is shown below, is the same code applied to the section of Different state spaces. The difference is that we are changing values of " $N$ ", epsilon, heights and radii:

```
r = 1; %(radius of the input)
h = 1; %(height of the input)
M = 5; %(number of points going down the Z axis (with (1,0,0,1) and (1,0,0,-1)))
N = 15; %(number of angles going around the equator)
epsilon = 0.1;
theta = 360/N;
height = [1 0.75 0 -0.75 -1]; %(heights of the five inputs)
width = [0, (\sqrt(7)/4)*(1+epsilon), 1, (\sqrt(7)/4)*(1+epsilon), 0]; %(Radii of th
```

$C Z=\backslash$ begin\{blockarray $\}\{*\{17\}\{c\}\}$
\begin\{block\}\{(*\{16\}\{c\})c\} }
$1 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \&(\$ I \backslash o t i m e s ~ I \$)$
$0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 1 \& 0 \& 0 \&(\$ I \backslash o t i m e s ~ X \$)$
0 \& 0 \& 0 \& 0 \& $0 \& 0 \& 0$ \& $0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0$ \& $1 \& 0 \&(\$ I \backslash o t i m e s ~ Y \$)$
0 \& 0 \& 0 \& $1 \& 0$ \& 0 \& 0 \& $0 \& 0 \& 0$ \& 0 \& $0 \& 0 \& 0 \& 0 \& 0 \&(\$ I \backslash o t i m e s ~ Z \$) ~$
$0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 1 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \&(\$ X \backslash o t i m e s ~ I \$)$
0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& $0 \& 0 \& 1 \& 0 \& 0 \& 0 \& 0 \& 0 \&(\$ X \backslash o t i m e s ~ X \$)$
$0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \&-1 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \&(\$ X \backslash o t i m e s ~ Y \$)$
$0 \& 0 \& 0 \& 0 \& 1 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \&(\$ X \backslash o t i m e s ~ Z \$)$
0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& $1 \& 0$ \& 0 \& 0 \& 0 \& (\$Y
$0 \& 0 \& 0 \& 0 \& 0 \& 0 \&-1 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \&(\$ Y \backslash o t i m e s ~ X \$)$
0 \& 0 \& 0 \& $0 \& 0 \& 1 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \&(\$ Y \backslash o t i m e s ~ Y \$)$
0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& $1 \& 0$ \& $0 \& 0 \& 0 \& 0$ \& 0 \& 0 \& (\$Y\otimes $Z \$)$
$0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 1 \& 0 \& 0 \& 0 \&(\$ Z \backslash o t i m e s ~ I \$)$
0 \& $1 \& 0$ \& 0 \& $0 \& 0 \& 0$ \& $0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \&(\$ Z \backslash o t i m e s ~ X \$)$
0 \& 0 \& 1 \& $0 \& 0$ \& 0 \& 0 \& $0 \& 0 \& 0$ \& $0 \& 0 \& 0 \& 0 \& 0 \& 0 \&(\$ Z \backslash t i m e s ~ Y \$)$
$0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 0 \& 1 \&(\$ Z \backslash o t i m e s ~ Z \$)$
\end\{block\} }
\end\{blockarray\} }
output_shape $=$ zeros $(4, N(M-2)+2) ; \%$ a (4 4 times $N(M-2)+2)$ matrix

```
% Linear programming in Matlab: linprog(f, A, b, A', b')
C = zeros((N(M-2)+2)^2,1); % the objective function
I = - eye((N(M-2)+2)^2) ; % A inequality
O = zeros((N(M-2)+2)^2,1) ; % b inequality
```

for $R=2.4: 0.001: 2.46$
\% Output

```
init = 3;
for k = 2: M-1
for phi = 0:360/N:360-360/N
output_shape(:, init) = [1; R*width(k)*\operatorname{cos(phi); R*width(k)*sin(phi); height(k)];}
init = init+1;
end
end
```

output_shape $(:, 1)=[1 ; 0 ; 0 ; 1] ; \%((:, 1)$ means rows of a specific col (in this case i
output_shape (:,2) = [1; 0; 0;-1];
out_shape $=$ kron(output-shape,output-shape) ; \% (out_shape is A equality in linear prog
\% Input
number_R = 0 ;
optimal_R = 0;
for ll = 1:M
for $\mathrm{mm}=1: \mathrm{M}$
input $=$ kron([1; width(ll); 0; height(ll)], [1; width(mm); 0; height(mm)]);
input_CZ = CZ*input; \%(input_CZ is b equality in linear programming)
Xopt = linprog(C, I, O, out_shape, input_CZ);
if Xopt >= 0
number_R = number_R + 1;
if $\mathrm{mm}==\mathrm{M}$ and $\mathrm{ll}==\mathrm{M}$ and number_ $R==\mathrm{M}^{\wedge} 2$
optimal_R = R;
end
else

> break
end
end
if Xopt >= 0
continue
else
break
end
end
if number_R == M^2 break
else
continue
end
end
\% For the linear programming, the formula that we need to use is as follow: linprog(f where
$\% C=\operatorname{zeros}\left((N(M-2)+2)^{\wedge} 2,1\right) ; \%$ which is the objective function,
$\% I=-\operatorname{eye}\left((N *(M-2)+2)^{\wedge} 2\right) ; \%$ which is the $A$ inequality,
$\% \quad 0=\operatorname{zeros}\left((N *(M-2)+2)^{\wedge} 2,1\right)$; $\%$ which is the $b$ inequality,
\% The out_shape is the A equality,
\% The input_CZ is b equality,
$\% l b=u b=0$; that are the lower and the upper bound respectively.
\% So by replacing the above formula by the equivalent notation, we have Xopt $=\operatorname{linprog}(C, I, O$, out-shape, input-CZ).

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## Bibliography

[1] Scott Aaronson and Daniel Gottesman. "Improved simulation of stabilizer circuits". In: Physical Review A 70.5 (2004), p. 052328.
[2] Scott Aaronson and Sam Gunn. "On the classical hardness of spoofing linear crossentropy benchmarking". In: arXiv preprint arXiv:1910.12085 (2019).
[3] Hussain Anwar et al. "Families of pure PEPS with efficiently simulatable local hidden variable models for most measurements". In: arXiv preprint arXiv:1412.3780 (2014).
[4] Hussain Anwar et al. "Generalised versions of separable decompositions applicable to bipartite entangled quantum states". In: New Journal of Physics 21.9 (2019), p. 093031.
[5] Hussain Anwar et al. "Smallest state spaces for which bipartite entangled quantum states are separable". In: New Journal of Physics 17.9 (2015), p. 093047.
[6] Frank Arute et al. "Quantum supremacy using a programmable superconducting processor". In: Nature 574.7779 (2019), pp. 505-510.
[7] Sahar Atallah et al. "Efficient classical simulation of cluster state quantum circuits with alternative inputs". In: arXiv preprint arXiv:2201.07655 (2022).
[8] Remigiusz Augusiak, Maciej Demianowicz, and Antonio Acın. "Local hidden-variable models for entangled quantum states". In: Journal of Physics A: Mathematical and Theoretical 47.42 (2014), p. 424002.
[9] Boaz Barak, Chi-Ning Chou, and Xun Gao. "Spoofing linear cross-entropy benchmarking in shallow quantum circuits". In: arXiv preprint arXiv:2005.02421 (2020).
[10] Jonathan Barrett. "Information processing in generalized probabilistic theories". In: Physical Review A 75.3 (2007), p. 032304.
[11] Jonathan Barrett et al. "The computational landscape of general physical theories". In: npj Quantum Information 5.1 (2019), pp. 1-10.
[12] Sean D Barrett et al. "Transitions in the computational power of thermal states for measurement-based quantum computation". In: Physical Review A 80.6 (2009), p. 062328.
[13] Stephen D Bartlett and Barry C Sanders. "Efficient classical simulation of optical quantum information circuits". In: Physical review letters 89.20 (2002), p. 207903.
[14] Stephen D Bartlett et al. "Efficient classical simulation of continuous variable quantum information processes". In: Physical Review Letters 88.9 (2002), p. 097904.
[15] John S Bell. "On the einstein podolsky rosen paradox". In: Physics Physique Fizika 1.3 (1964), p. 195.
[16] Jacob Biamonte and Ville Bergholm. "Tensor networks in a nutshell". In: arXiv preprint arXiv:1708.00006 (2017).
[17] Sergio Boixo et al. "Characterizing quantum supremacy in near-term devices". In: Nature Physics 14.6 (2018), pp. 595-600.
[18] Stephen Boyd and Lieven Vandenberghe. Convex optimization. Cambridge university press, 2004.
[19] Sergey Bravyi and David Gosset. "Improved classical simulation of quantum circuits dominated by Clifford gates". In: Physical review letters 116.25 (2016), p. 250501.
[20] Daniel J Brod. "Efficient classical simulation of matchgate circuits with generalized inputs and measurements". In: Physical Review A 93.6 (2016), p. 062332.
[21] Nicolas Brunner et al. "Bell nonlocality". In: Reviews of Modern Physics 86.2 (2014), p. 419.
[22] Andrew M Childs. "Lecture notes on quantum algorithms". In: Lecture notes at University of Maryland (2017).
[23] Andrew M Childs and Wim Van Dam. "Quantum algorithms for algebraic problems". In: Reviews of Modern Physics 82.1 (2010), p. 1.
[24] John F Clauser et al. "Proposed experiment to test local hidden-variable theories". In: Physical review letters 23.15 (1969), p. 880.
[25] Andrew C Doherty, Pablo A Parrilo, and Federico M Spedalieri. "Distinguishing separable and entangled states". In: Physical Review Letters 88.18 (2002), p. 187904.
[26] Albert Einstein, Boris Podolsky, and Nathan Rosen. "Can quantum-mechanical description of physical reality be considered complete?" In: Physical review 47.10 (1935), p. 777.
[27] Daniel Gottesman. "The Heisenberg representation of quantum computers". In: arXiv preprint quant-ph/9807006 (1998).
[28] D Gross and J Eisert. "Novel schemes for measurement-based quantum computation". In: Physical review letters 98.22 (2007), p. 220503.
[29] Daivd Gross et al. "Measurement-based quantum computation beyond the one-way model". In: Physical Review A 76.5 (2007), p. 052315.
[30] Lov K Grover. "A fast quantum mechanical algorithm for database search". In: Proceedings of the twenty-eighth annual ACM symposium on Theory of computing. 1996, pp. 212-219.
[31] Lucien Hardy. "Quantum theory from five reasonable axioms". In: arXiv preprint quant-ph/0101012 (2001).
[32] Aram W Harrow and Ashley Montanaro. "Quantum computational supremacy". In: Nature 549.7671 (2017), pp. 203-209.
[33] Aram W Harrow and Michael A Nielsen. "Robustness of quantum gates in the presence of noise". In: Physical Review A 68.1 (2003), p. 012308.
[34] Michał Horodecki, Paweł Horodecki, and Ryszard Horodecki. "Separability of mixed states: necessary and sufficient conditions". In: Physics Letters A 223.1 (1996), pp. 18. ISSN: 0375-9601. DOI: https://doi.org/10.1016/S0375-9601 (96)00706-2. URL: https://www.sciencedirect.com/science/article/pii/S0375960196007062.
[35] Pawel Horodecki. "Separability criterion and inseparable mixed states with positive partial transposition". In: Physics Letters A 232.5 (1997), pp. 333-339.
[36] Peter Janotta and Haye Hinrichsen. "Generalized probability theories: what determines the structure of quantum theory?" In: Journal of Physics A: Mathematical and Theoretical 47.32 (2014), p. 323001.
[37] Richard Jozsa. "On the simulation of quantum circuits". In: arXiv preprint quantph/0603163 (2006).
[38] Richard Jozsa and Noah Linden. "On the role of entanglement in quantum-computational speed-up". In: Proceedings of the Royal Society of London. Series A: Mathematical, Physical and Engineering Sciences 459.2036 (2003), pp. 2011-2032.
[39] Richard Jozsa and Akimasa Miyake. "Matchgates and classical simulation of quantum circuits". In: Proceedings of the Royal Society A: Mathematical, Physical and Engineering Sciences 464.2100 (2008), pp. 3089-3106.
[40] Richard Jozsa and Maarten Van den Nest. "Classical simulation complexity of extended Clifford circuits". In: arXiv preprint arXiv:1305.6190 (2013).
[41] Igor L Markov and Yaoyun Shi. "Simulating quantum computation by contracting tensor networks". In: SIAM Journal on Computing 38.3 (2008), pp. 963-981.
[42] Ashley Montanaro. "Quantum algorithms: an overview". In: npj Quantum Information 2.1 (2016), pp. 1-8.
[43] Caterina E Mora et al. "Universal resources for approximate and stochastic measurementbased quantum computation". In: Physical Review A 81.4 (2010), p. 042315.
[44] Michele Mosca. "Quantum algorithms". In: arXiv preprint arXiv:0808.0369 (2008).
[45] Michael A Nielsen. "Cluster-state quantum computation". In: Reports on Mathematical Physics 57.1 (2006), pp. 147-161.
[46] Michael A Nielsen and Isaac Chuang. Quantum computation and quantum information. 2002.
[47] Asher Peres. "Separability criterion for density matrices". In: Physical Review Letters 77.8 (1996), p. 1413.
[48] David Perez-Garcia et al. "Matrix product state representations". In: arXiv preprint quant-ph/0608197 (2006).
[49] Martin Plávala. "General probabilistic theories: An introduction". In: arXiv preprint arXiv:2103.07469 (2021).
[50] John Preskill. "Lecture notes for physics 229: Quantum information and computation". In: California Institute of Technology 16.1 (1998), pp. 1-8.
[51] N Ratanje and S Virmani. "Exploiting non-quantum entanglement to widen applicability of limited-entanglement classical simulations of quantum systems". In: arXiv preprint arXiv:1201.0613 (2012).
[52] N Ratanje and Shashank Virmani. "Generalized state spaces and nonlocality in faulttolerant quantum-computing schemes". In: Physical Review A 83.3 (2011), p. 032309.
[53] Robert Raussendorf and Hans J Briegel. "A one-way quantum computer". In: Physical Review Letters 86.22 (2001), p. 5188.
[54] Robert Raussendorf, Daniel E Browne, and Hans J Briegel. "Measurement-based quantum computation on cluster states". In: Physical review A 68.2 (2003), p. 022312.
[55] Norbert Schuch et al. "Computational complexity of projected entangled pair states". In: Physical review letters 98.14 (2007), p. 140506.
[56] Peter W Shor. "Introduction to quantum algorithms". In: Proceedings of Symposia in Applied Mathematics. Vol. 58. 2002, pp. 143-160.
[57] Peter W Shor. "Polynomial-time algorithms for prime factorization and discrete logarithms on a quantum computer". In: SIAM review 41.2 (1999), pp. 303-332.
[58] Barbara M Terhal and David P DiVincenzo. "Classical simulation of noninteractingfermion quantum circuits". In: Physical Review A 65.3 (2002), p. 032325.
[59] Leslie G Valiant. "Expressiveness of matchgates". In: Theoretical Computer Science 289.1 (2002), pp. 457-471.
[60] Leslie G Valiant. "Quantum circuits that can be simulated classically in polynomial time". In: SIAM Journal on Computing 31.4 (2002), pp. 1229-1254.
[61] Maarten Van den Nest. "Universal quantum computation with little entanglement". In: Physical review letters 110.6 (2013), p. 060504.
[62] Guifré Vidal. "Efficient classical simulation of slightly entangled quantum computations". In: Physical review letters 91.14 (2003), p. 147902.
[63] Reinhard F Werner. "Quantum states with Einstein-Podolsky-Rosen correlations admitting a hidden-variable model". In: Physical Review A 40.8 (1989), p. 4277.
[64] Mithuna Yoganathan, Richard Jozsa, and Sergii Strelchuk. "Quantum advantage of unitary Clifford circuits with magic state inputs". In: Proceedings of the Royal Society A 475.2225 (2019), p. 20180427.
[65] Nadav Yoran. "Efficiently contractable quantum circuits cannot produce much entanglement". In: arXiv preprint arXiv:0802.1156 (2008).
[66] Nadav Yoran and Anthony J Short. "Classical simulation of limited-width clusterstate quantum computation". In: Physical review letters 96.17 (2006), p. 170503.

