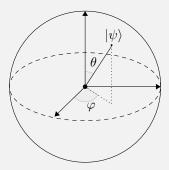
# Solutions to Quantum Computation and Quantum Information



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#### Introduction:

This document is a (work in progress) collection of comprehensive solutions to the exercises in Nielsen and Chuang's "Quantum Computation and Quantum Information". Each solution has the involved concepts (and hence rough pre-requisite knowledge) necessary for the problem in addition to the solution. Some problems may contain additional remarks about implications. Any problems denoted as (Research) are left as exercises to the reader. Starred exercises are considered to be more difficult (difficulty is assumed for the problems at the end of the chapter).

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# 1 Introduction and overview

### Exercise 1.1: Probabilistic Classical Algorithm

Suppose that the problem is not to distinguish between the constant and balanced functions with certainty, but rather, with some probability of error  $\epsilon < 1/2$ . What is the performance of the best classical algorithm for this problem?

#### Solution

Concepts Involved: Deutsch's Problem, Probability.

Recall that a Boolean function  $f:\{0,1\}^n \to \{0,1\}$  is said to be balanced if f(x)=1 for exactly half of all possible  $2^n$  values of x.

A single evaluation tells us no information about whether f is constant or balanced, so our success rate/error rate after a single evaluation is  $\epsilon = \frac{1}{2}$  (random guessing!). Therefore, consider the case where we do two evaluations. If we obtain two different results, then we immediately conclude that f is balanced. Suppose instead that we obtain two results that are the same. If f is balanced, then the probably that the first evaluation returns the given result is  $\frac{1}{2}$ , and the probability that the second evaluation returns the same result is  $\frac{2^n/2-1}{2^n-1}$  (as there are  $2^n/2$  of each result of 0 and 1,  $2^n$  total results,  $2^n/2-1$  of the given result left after the first evaluation, and  $2^n-1$  total uninvestigated cases after the first evaluation). Therefore, if f is balanced, this occurs with probability  $\frac{1}{2} \cdot \frac{2^n/2-1}{2^n-1}$ , which we can see is less than  $\frac{1}{2}$  as:

$$2^{n} < 2^{n+1} \implies 2^{n} - 2 < 2^{n+1} - 2 \implies \frac{2^{n}/2 - 1}{2^{n} - 1} < 1 \implies \frac{1}{2} \frac{2^{n}/2 - 1}{2^{n} - 1} < \frac{1}{2}$$

Hence, if we get the same result in two evaluations, we can conclude that f is constant with error  $\epsilon < \frac{1}{2}$ . We conclude that only 2 evaluations are required for this algorithm.

# Exercise 1.2

Explain how a device which, upon input of one of two non-orthogonal quantum states  $|\psi\rangle$  or  $|\varphi\rangle$  correctly identified the state, could be used to build a device which cloned the states  $|\psi\rangle$  and  $|\varphi\rangle$ , in violation of the no-cloning theorem. Conversely, explain how a device for cloning could be used to distinguish non-orthogonal quantum states.

### Solution

**Concepts Involved:** Quantum Distinguishability, Quantum Measurement.

Given access to a device which can distinguish non-orthogonal quantum states  $|\psi\rangle$ ,  $|\varphi\rangle$  (without measurement), we can then design a quantum circuit that would map  $|\psi\rangle\mapsto|\varphi\rangle$  (or vise versa), allowing us to clone the states as we like.

Conversely, given a cloning device, we could clone  $|\psi\rangle$  and  $|\varphi\rangle$  an arbitrary number of times. Then, performing repeated measurements of the two states in different measurement bases, we would (given enough measurements) be able to distinguish the two states based on the measurement statistics (there will of course be some error  $\epsilon$  based on probabilistic considerations, but given that we have access to as many measurements of the states as we like, we are able to make this error arbitrarily low).

# **Problem 1.1: (Feynman-Gates conversation)**

Construct a friendly imaginary discussion of about 2000 words between Bill Gates and Richard Feynman, set in the present, on the future of computation (Comment: You might like to try waiting until you've heard the rest of the book before attempting this question. See 'History and further reading' below for pointers to one possible answer for this question).

# Problem 1.2

What is the most significant discovery yet made in quantum computation and quantum information? Write an essay of about 2000 words to an educated lay audience about the discovery (Comment: As for the previous problem, you might like to try waiting until you've read the rest of the book before attempting this question.)

# 2 Introduction to quantum mechanics

### Exercise 2.1: Linear dependence: example

Show that (1,-1), (1,2) and (2,1) are linearly dependent.

### Solution

Concepts Involved: Linear Algebra, Linear Independence/Dependence.

We observe that:

$$\begin{bmatrix} 1 \\ -1 \end{bmatrix} + \begin{bmatrix} 1 \\ 2 \end{bmatrix} - \begin{bmatrix} 2 \\ 1 \end{bmatrix} = \begin{bmatrix} 1+1-2 \\ -1+2-1 \end{bmatrix} = \begin{bmatrix} 0 \\ 0 \end{bmatrix}$$

showing that the three vectors are linearly dependent by definition. Alternatively, we can apply theorem that states that for any vector space V with  $\dim V = n$ , any list of m > n vectors in V will be linearly dependent (here,  $V = \mathbb{R}^2, n = 2, m = 3$ ).

### Exercise 2.2: Matrix representations: example

Suppose V is a vector space with basis vectors  $|0\rangle$  and  $|1\rangle$ , and A is a linear operator from V to V such that  $A\,|0\rangle\,=\,|1\rangle$  and  $A\,|1\rangle\,=\,|0\rangle$ . Give a matrix representation for A, with respect to the input basis  $|0\rangle\,,\,|1\rangle$ , and the output basis  $|0\rangle\,,\,|1\rangle$ . Find input and output bases which give rise to a different matrix representation of A.

#### Solution

Concepts Involved: Linear Algebra, Matrix Representation of Operators.

Identifying 
$$|0\rangle\cong\begin{bmatrix}1\\0\end{bmatrix}$$
 and  $|1\rangle\cong\begin{bmatrix}0\\1\end{bmatrix}$ , we have that:

$$A = \begin{bmatrix} a_{00} & a_{01} \\ a_{10} & a_{11} \end{bmatrix}$$

Using the given relations, we have that:

$$A |0\rangle = 0 |0\rangle + 1 |1\rangle \implies \begin{bmatrix} a_{00} & a_{01} \\ a_{10} & a_{11} \end{bmatrix} \begin{bmatrix} 1 \\ 0 \end{bmatrix} = 0 \begin{bmatrix} 1 \\ 0 \end{bmatrix} + 1 \begin{bmatrix} 0 \\ 1 \end{bmatrix} \implies a_{00} = 0, a_{10} = 1$$

$$A |1\rangle = 1 |0\rangle + 0 |1\rangle \implies \begin{bmatrix} a_{00} & a_{01} \\ a_{10} & a_{11} \end{bmatrix} \begin{bmatrix} 0 \\ 1 \end{bmatrix} = 1 \begin{bmatrix} 1 \\ 0 \end{bmatrix} + 0 \begin{bmatrix} 0 \\ 1 \end{bmatrix} \implies a_{01} = 1, a_{11} = 0$$

Therefore with respect to the input basis  $\{ |0\rangle, |1\rangle \}$  and output basis  $\{ |0\rangle, |1\rangle \}$ , A has matrix represen-

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tation:

$$A\cong\begin{array}{ccc} |0\rangle \left[ & \langle 0| & \langle 1| \\ 0 & 1 \\ 1 & 0 \end{array} \right]$$

Suppose we instead choose the input and output basis to be  $\left\{ |+\rangle = \frac{|0\rangle + |1\rangle}{\sqrt{2}}, \; |-\rangle = \frac{|0\rangle - |1\rangle}{\sqrt{2}} \right\}$ . Identifying  $|+\rangle \cong \begin{bmatrix} 1 \\ 0 \end{bmatrix}$  and  $|-\rangle \cong \begin{bmatrix} 0 \\ 1 \end{bmatrix}$ , we have:

$$A = \begin{bmatrix} a_{++} & a_{+-} \\ a_{-+} & a_{--} \end{bmatrix}$$

Using the linearity of A, we have that:

$$A|+\rangle = \frac{1}{\sqrt{2}}A(|0\rangle + |1\rangle) = \frac{1}{\sqrt{2}}(A|0\rangle + A|1\rangle) = \frac{1}{\sqrt{2}}(|1\rangle + |0\rangle) = |+\rangle$$

and:

$$A\left|-\right\rangle = \frac{1}{\sqrt{2}}A(\left|0\right\rangle - \left|1\right\rangle) = \frac{1}{\sqrt{2}}(A\left|0\right\rangle - A\left|1\right\rangle) = \frac{1}{\sqrt{2}}(\left|1\right\rangle - \left|0\right\rangle) = -\left|-\right\rangle,$$

which can be used to determine the matrix elements:

$$\begin{split} A \left| + \right\rangle &= 1 \left| + \right\rangle + 0 \left| - \right\rangle \implies \begin{bmatrix} a_{++} & a_{+-} \\ a_{-+} & a_{--} \end{bmatrix} \begin{bmatrix} 1 \\ 0 \end{bmatrix} = 1 \begin{bmatrix} 1 \\ 0 \end{bmatrix} + 0 \begin{bmatrix} 0 \\ 1 \end{bmatrix} \implies a_{++} = 1, a_{+-} = 0 \\ A \left| - \right\rangle &= 0 \left| + \right\rangle - 1 \left| - \right\rangle \implies \begin{bmatrix} a_{++} & a_{+-} \\ a_{-+} & a_{--} \end{bmatrix} \begin{bmatrix} 0 \\ 1 \end{bmatrix} = 0 \begin{bmatrix} 1 \\ 0 \end{bmatrix} - 1 \begin{bmatrix} 0 \\ 1 \end{bmatrix} \implies a_{-+} = 0, a_{--} = -1 \end{split}$$

Therefore with respect to the input basis  $\{ |+\rangle \,,\, |-\rangle \}$  and output basis  $\{ |+\rangle \,,\, |-\rangle \}$ , A has matrix representation:

$$A \cong \begin{array}{c|c} |+\rangle & \langle +| & \langle -| \\ 1 & 0 \\ |-\rangle & 0 & -1 \end{array}$$

**Remark:** If we choose the input and output bases to be different, we can even represent the A operator as an identity matrix. Specifically, if the input basis to be chosen to be  $\{|0\rangle, |1\rangle\}$  and output basis as  $\{|1\rangle, |0\rangle\}$ , the matrix representation of A looks like:

$$A \cong \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix}.$$

### Exercise 2.3: Matrix representation for operator products

Suppose A is a linear operator from vector space V to vector space W, and B is a linear operator from vector space W to vector space X. Let  $|v_i\rangle$ ,  $|w_j\rangle$ ,  $|x_k\rangle$  be bases for the vector spaces V,W and X respectively. Show that the matrix representation for the linear transformation BA is the matrix product of the matrix representations for B and A, with respect to the appropriate bases.

#### Solution

Concepts Involved: Linear Algebra, Matrix Representation of Operators.

Taking the matrix of representations of A and B to the appropriate bases  $|v_i\rangle$ ,  $|w_j\rangle$ ,  $|x_k\rangle$  of V,W and X, we have that:

$$A|v_j\rangle = \sum_i A_{ij} |w_i\rangle, \quad B|w_i\rangle = \sum_k B_{ki} |x_k\rangle$$

Hence, looking at  $BA: V \mapsto X$ , we have that:

$$BA |v_{j}\rangle = B(A |v_{j}\rangle)$$

$$= B\left(\sum_{i} A_{ij} |w_{i}\rangle\right)$$

$$= \sum_{i} A_{ij} B |w_{i}\rangle$$

$$= \sum_{i} A_{ij} \left(\sum_{k} B_{ki} |x_{k}\rangle\right)$$

$$= \sum_{k} \sum_{i} B_{ki} A_{ij} |x_{k}\rangle$$

$$= \sum_{k} (BA)_{kj} |x_{k}\rangle$$

which shows that the matrix representation of BA is indeed the matrix product of the representations of B and A.

### Exercise 2.4: Matrix representation for identity

Show that the identity operator on a vector space V has a matrix representation which is one along the diagonal and zero everywhere else, if the matrix is taken with respect to the same input and output bases. This matrix is known as the *identity matrix*.

### Solution

Concepts Involved: Linear Algebra, Matrix Representation of Operators.

Let V be a vector space and  $|v_i\rangle$  be a basis of V. Let  $A:V\mapsto V$  be a linear operator, and let its matrix representation taken to be respect to  $|v_i\rangle$  as the input and output basis. We then have that for each

 $i \in \{1, \ldots, n\}$ :

$$A|v_i\rangle = 1|v_i\rangle + \sum_{j\neq i} 0|v_j\rangle = \sum_j \delta_{ij}|v_j\rangle$$

From which we obtain that  $\boldsymbol{A}$  has the matrix representation:

$$A \cong \begin{bmatrix} 1 & 0 & 0 & \cdots & 0 \\ 0 & 1 & 0 & & 0 \\ 0 & 0 & 1 & & 0 \\ \vdots & & & \ddots & \\ 0 & 0 & 0 & \dots & 1 \end{bmatrix}$$

Exercise 2.5

Verify that  $(\cdot, \cdot)$  just defined is an inner product on  $\mathbb{C}^n$ .

Solution

Concepts Involved: Linear Algebra, Inner Products.

Recall that on  $\mathbb{C}^n$ ,  $(\cdot,\cdot)$  was defined as:

$$((y_1,\ldots,y_n),(z_1,\ldots,z_n)) \equiv \sum_i y_i^* z_i = [y_1^*\ldots y_n^*] \begin{bmatrix} z_1 \\ \vdots \\ z_n \end{bmatrix}.$$

Furthermore, recall the three conditions for the function  $(\cdot,\cdot):V\times V\mapsto\mathbb{C}$  to be considered an inner product:

- (1)  $(\cdot, \cdot)$  is linear in the second argument.
- (2)  $(|v\rangle, |w\rangle) = (|w\rangle, |v\rangle)^*$ .
- (3)  $(|v\rangle, |v\rangle) \ge 0$  with equality if and only if  $|v\rangle = 0$ .

We check that  $(\cdot,\cdot):\mathbb{C}^n imes\mathbb{C}^n\mapsto\mathbb{C}$  satisfies the three conditions:

(1) We see that:

$$((y_1, \dots, y_n), \sum_k \lambda_k(z_1, \dots, z_n)_k) = \sum_i y_i^* \sum_k \lambda_k z_{i_k}$$

$$= \sum_k \lambda_k \sum_i y_i^* z_{i_k}$$

$$= \sum_k \lambda_k ((y_1, \dots, y_n), (z_1, \dots, z_n)_k)$$

(2) We have:

$$((y_1, \dots, y_n), (z_1, \dots, z_n)) = \sum_i y_i^* z_i$$

$$= \sum_i (y_i z_i^*)^*$$

$$= \left(\sum_i z_i^* y_i\right)^*$$

$$= (((z_1, \dots, z_n), (y_1, \dots, y_n)))^*$$

(3) We observe for 0 = (0, ... 0):

$$(\mathbf{0}, \mathbf{0}) = \sum_{i} 0 \cdot 0 = 0$$

For  $\mathbf{y}=(y_1,\ldots,y_n) \neq 0$  we have that at least one  $y_i$  (say,  $y_j$ ) is nonzero, and hence:

$$((y_1,\ldots,y_n),(y_1,\ldots,y_n)) = \sum_i y_i^2 \ge y_j^2 > 0$$

which proves the claim.

Exercise 2.6

Show that any inner product  $(\cdot,\cdot)$  is conjugate-linear in the first argument,

$$\left(\sum_{i} \lambda_{i} |w_{i}\rangle, |v\rangle\right) = \sum_{i} \lambda_{i}^{*}(|w_{i}\rangle, |v\rangle).$$

Solution

Concepts Involved: Linear Algebra, Inner Products

Applying properties (2) (conjugate symmetry), (1) (linearity in second argument), and (2) (again) in

succession, we have that:

$$\left(\sum_{i} \lambda_{i} |w_{i}\rangle, |v\rangle\right) = \left(|v\rangle, \sum_{i} \lambda_{i} |w_{i}\rangle\right)^{*}$$

$$= \left(\sum_{i} \lambda_{i} (|v\rangle, |w_{i}\rangle)\right)^{*}$$

$$= \sum_{i} \lambda_{i}^{*} (|v\rangle, |w_{i}\rangle)^{*}$$

$$= \sum_{i} \lambda_{i}^{*} (|w_{i}\rangle, |v\rangle)$$

# Exercise 2.7

Verify that  $|w\rangle=(1,1)$  and  $|v\rangle=(1,-1)$  are orthogonal. What are the normalized forms of these vectors?

### Solution

**Concepts Involved:** Linear Algebra, Inner Products, Orthogonality, Normalization Recall that two vectors  $|v\rangle, |w\rangle$  are orthogonal if  $\langle v|w\rangle = 0$ , and the norm of  $|v\rangle$  is given by  $\||v\rangle\| = \sqrt{\langle v|v\rangle}$ .

First we show the two vectors are orthogonal:

$$\langle w|v\rangle = 1 \cdot 1 + 1 \cdot (-1) = 0$$

The norms of  $|w\rangle$ ,  $|v\rangle$  are given by:

$$\begin{aligned} \left\| |w\rangle \right\| &= \sqrt{\langle w|w\rangle} = \sqrt{1^2 + 1^2} = \sqrt{2}, \\ \left\| |c\rangle \right\| &= \sqrt{\langle v|v\rangle} = \sqrt{1^2 + (-1)^2} = \sqrt{2} \end{aligned}$$

So the normalized forms of the vectors are:

$$\begin{split} \frac{|w\rangle}{\left\||w\rangle\right\|} &= \begin{bmatrix} \frac{1}{\sqrt{2}} \\ \frac{1}{\sqrt{2}} \end{bmatrix} \\ \frac{|v\rangle}{\left\||v\rangle\right\|} &= \begin{bmatrix} \frac{1}{\sqrt{2}} \\ \frac{-1}{\sqrt{2}} \end{bmatrix} \end{split}$$

# Exercise 2.8

Verify that the Gram-Schmidt procedure produces and orthonormal basis for V.

#### Solution

**Concepts Involved:** Linear Algebra, Linear Independence, Bases, Inner Products, Orthogonality, Normalization, Gram-Schmidt Procedure, Induction.

Recall that given  $|w_1\rangle,\ldots,|w_d\rangle$  as a basis set for a vector space V, the Gram-Schmidt procedure constructs a basis set  $|v_1\rangle,\ldots,|v_d\rangle$  by defining  $|v_1\rangle\equiv|w_1\rangle/\big||w_1\rangle\big||$  and then defining  $|v_{k+1}\rangle$  inductively for  $1\leq k\leq d-1$  as:

$$|v_{k+1}\rangle \equiv \frac{|w_{k+1}\rangle - \sum_{i=1}^{k} \langle v_i | w_{k+1}\rangle |v_i\rangle}{\left\||w_{k+1}\rangle - \sum_{i=1}^{k} \langle v_i | w_{k+1}\rangle |v_i\rangle\right\|}$$

It is evident that each of the  $|v_j\rangle$  have unit norm as they are defined in normalized form. It therefore suffices to show that each of the  $|v_1\rangle,\ldots,|v_d\rangle$  are orthogonal to each other, and that this set of vectors forms a basis of V. We proceed by induction. For k=1, we have that:

$$|v_2\rangle = \frac{|w_2\rangle - \langle v_1|w_2\rangle |v_1\rangle}{\||w_2\rangle + \langle v_1|w_2\rangle |v_1\rangle\|}$$

Therefore:

$$\langle v_1 | v_2 \rangle = \frac{\langle v_1 | w_2 \rangle - \langle v_1 | w_2 \rangle \langle v_1 | v_1 \rangle}{\| |w_2 \rangle + \langle v_1 | w_2 \rangle |v_1 \rangle \|} = \frac{\langle v_1 | w_2 \rangle - \langle v_1 | w_2 \rangle}{\| |w_2 \rangle + \langle v_1 | w_2 \rangle |v_1 \rangle \|} = 0$$

so the two vectors are orthogonal. Furthermore, they are linearly independent; if they were linearly dependent, we could write  $|v_1\rangle = \lambda |v_2\rangle$  for some  $\lambda \in \mathbb{C}$ , but then multiplying both sides by  $\langle v_1|$  we get:

$$\langle v_1 | v_1 \rangle = \lambda \langle v_1 | v_2 \rangle \implies 1 = 0$$

which is a contradiction. This concludes the base case. For the inductive step, let  $k \ge 1$  and suppose that  $|v_1\rangle, \ldots, |v_k\rangle$  are orthogonal and linearly independent. We then have that:

$$|v_{k+1}\rangle = \frac{|w_{k+1}\rangle - \sum_{i=1}^{k} \langle v_i | w_{k+1}\rangle |v_i\rangle}{\left\||w_{k+1}\rangle - \sum_{i=1}^{k} \langle v_i | w_{k+1}\rangle |v_i\rangle\right\|}$$

Then for any  $j \in \{1, \dots k\}$ , we have that:

$$\langle v_j | v_{k+1} \rangle = \frac{\langle v_j | w_{k+1} \rangle - \sum_{i=1}^k \langle v_i | w_{k+1} \rangle \langle v_j | v_i \rangle}{\left\| | w_{k+1} \rangle - \sum_{i=1}^k \langle v_i | w_{k+1} \rangle | v_i \rangle \right\|} = \frac{\langle v_j | w_{k+1} \rangle - \langle v_j | w_{k+1} \rangle}{\left\| | w_{k+1} \rangle - \sum_{i=1}^k \langle v_i | w_{k+1} \rangle | v_i \rangle \right\|} = 0$$

where in the second equality we use the fact that  $\langle v_j | v_i \rangle = \delta_{ij}$  for  $i,j \in \{1,\dots k\}$  by the inductive hypothesis. We therefore find that  $|v_{k+1}\rangle$  is orthogonal to all of  $|v_1\rangle,\dots,|v_k\rangle$ . Furthermore,  $|v_1\rangle,\dots,|v_k\rangle,|v_{k+1}\rangle$  is lienarly independent. Suppose for the sake of contradiction that this was false. Then, there would exist  $\lambda_1,\dots \lambda_k$  not all nonzero such that:

$$\lambda_1|v_1\rangle + \ldots + \lambda_k|v_k\rangle = |v_{k+1}\rangle$$

but then multiplying both sides by  $\langle v_{k+1}|$  we have that:

$$\lambda_1 \langle v_{k+1} | v_1 \rangle + \ldots + \lambda_k \langle v_{k+1} | v_k \rangle = \langle v_{k+1} | v_{k+1} \rangle \implies 0 = 1$$

by orthonormality. This gives a contradiction, and hence  $|v_1\rangle,\ldots,|v_k\rangle,|v_{k+1}\rangle$  are linearly independent, finishing the inductive step. Therefore,  $|v_1\rangle,\ldots,|v_d\rangle$  is an orthonormal list of vectors which is linearly independent. Since  $|w_1\rangle,\ldots,|w_d\rangle$  is a basis for V, then V has dimension d. Hence,  $|v_1\rangle,\ldots,|v_d\rangle$  being a linearly independent list of d vectors in V is a basis of V. We conclude that it is an orthonormal basis of V, as claimed.

### Exercise 2.9: Pauli operators and the outer product

The Pauli matrices (Figure 2.2 on page 65) can be considered as operators with respect to an orthonormal basis  $|0\rangle, |1\rangle$  for a two-dimensional Hilbert space. Express each of the Pauli operators in the outer product notation.

### Solution

**Concepts Involved:** Linear Algebra, Matrix Representation of Operators, Outer Products. Recall that if A has matrix representation:

$$A \cong \begin{bmatrix} a_{00} & a_{01} \\ a_{10} & a_{11} \end{bmatrix}$$

with respect to  $|0\rangle, |1\rangle$  as the input/output bases, then we can express A in outer product notation as:

$$A = a_{00} |0\rangle\langle 0| + a_{01} |0\rangle\langle 1| + a_{10} |1\rangle\langle 0| + a_{11} |1\rangle\langle 1|$$

Furthermore, recall the representation of the Pauli matrices with respect to the orthonormal basis  $|0\rangle, |1\rangle$ :

$$I = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix} \quad X = \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} \quad Y = \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} \quad Z = \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix}$$

We immediately see that:

$$I = |0\rangle\langle 0| + |1\rangle\langle 1|$$

$$X = |0\rangle\langle 1| + |1\rangle\langle 0|$$

$$Y = -i|0\rangle\langle 1| + i|1\rangle\langle 0|$$

$$Z = |0\rangle\langle 0| - |1\rangle\langle 1|$$

Exercise 2.10

Suppose  $|v_i\rangle$  is an orthonormal basis for an inner product space V. What is the matrix representation for the operator  $|v_i\rangle\langle v_i|$ , with respect to the  $|v_i\rangle$  basis?

#### Solution

Concepts Involved: Linear Algebra, Matrix Representation of Operators, Outer Products.

The matrix representation of  $|v_j\rangle\langle v_j|$  with respect to the  $|v_i\rangle$  basis is a matrix with 1 in the jth column and row (i.e. the (j,j)th entry in the matrix) and 0 everywhere else.

#### Exercise 2.11

Find the eigenvectors, eigenvalues, and diagonal representations of the Pauli matrices X, Y and Z.

### Solution

Concepts Involved: Linear Algebra, Eigenvalues, Eigenvectors, Diagonalization.

Given an operator A on a vector space V, recall that an eigenvector  $|v\rangle$  of A and its corresponding eigenvalue  $\lambda$  are defined by:

$$A|v\rangle = \lambda |v\rangle$$

Furthermore, recall the diagonal representation of A is given by

$$A = \sum_{i} \lambda_i |i\rangle\langle i|$$

Where  $|i\rangle$  form an orthonormal set of eigenvectors for A, and  $\lambda_i$  are the corresponding eigenvalues.

We start with X. Solving for the eigenvalues, we have:

$$\det(X - I\lambda) = 0 \implies \det \begin{bmatrix} -\lambda & 1\\ 1 & -\lambda \end{bmatrix} = 0 \implies \lambda^2 - 1 = 0$$

From which we obtain  $\lambda_1=1, \lambda_2=-1$ . Solving for the eigenvectors, we then have that:

$$(X - I\lambda_1)|v_1\rangle = \mathbf{0} \implies \begin{bmatrix} -1 & 1\\ 1 & -1 \end{bmatrix} \begin{bmatrix} v_{11}\\ v_{12} \end{bmatrix} = \begin{bmatrix} 0\\ 0 \end{bmatrix} \implies v_{11} = 1, v_{12} = 1$$
$$(X - I\lambda_2)|v_2\rangle = \mathbf{0} \implies \begin{bmatrix} 1 & 1\\ 1 & 1 \end{bmatrix} \begin{bmatrix} v_{21}\\ v_{22} \end{bmatrix} = \begin{bmatrix} 0\\ 0 \end{bmatrix} \implies v_{21} = 1, v_{22} = -1$$

Hence we find that  $|v_1\rangle=|0\rangle+|1\rangle$ ,  $|v_2\rangle=|0\rangle-|1\rangle$ . Normalizing these eigenvectors (Also see Exercise 2.7), we divide by  $||v_1\rangle||=||v_2\rangle||=\sqrt{2}$ , giving us:

$$|v_1\rangle = |+\rangle = \frac{|0\rangle + |1\rangle}{\sqrt{2}}, \quad |v_2\rangle = |-\rangle = \frac{|0\rangle - |1\rangle}{\sqrt{2}}.$$

The diagonal representation of X is then given by:

$$X = \lambda_1 |v_1\rangle\langle v_1| + \lambda_2 |v_2\rangle\langle v_2| = |+\rangle\langle +|-|-\rangle\langle -|$$

We do the same for Y. Solving for the eigenvalues:

$$\det(A - I\lambda) = 0 \implies \det \begin{bmatrix} -\lambda & -i \\ i & -\lambda \end{bmatrix} = 0 \implies \lambda^2 - 1 = 0$$

From which we obtain  $\lambda_1=1, \lambda_2=-1$ . Solving for the eigenvectors, we then have that:

$$(Y - I\lambda_1)|v_1\rangle = \mathbf{0} \implies \begin{bmatrix} -1 & -i \\ i & -1 \end{bmatrix} \begin{bmatrix} v_{11} \\ v_{12} \end{bmatrix} = \begin{bmatrix} 0 \\ 0 \end{bmatrix} \implies v_{11} = 1, v_{12} = i$$

$$(Y - I\lambda_2)|v_2\rangle = \mathbf{0} \implies \begin{bmatrix} 1 & -i \\ i & 1 \end{bmatrix} \begin{bmatrix} v_{21} \\ v_{22} \end{bmatrix} = \begin{bmatrix} 0 \\ 0 \end{bmatrix} \implies v_{21} = 1, v_{22} = -i$$

We therefore have that  $|v_1\rangle = |0\rangle + i|1\rangle, |v_2\rangle = |0\rangle - i|1\rangle$ . Normalizing by dividing by  $||v_1\rangle|| = ||v_2\rangle||$ , we obtain that:

$$|v_1\rangle = |y_+\rangle = \frac{|0\rangle + i|1\rangle}{\sqrt{2}}, \quad |v_2\rangle = |y_-\rangle = \frac{|0\rangle - i|1\rangle}{\sqrt{2}}.$$

The diagonal representation of Y is then given by:

$$Y = |y_{+}\rangle\langle y_{+}| - |y_{-}\rangle\langle y_{-}|$$

For Z, the process is again the same. We give the results and omit the details:

$$\lambda_1 = 1, |v_1\rangle = |0\rangle$$
  $\lambda_2 = -1, |v_2\rangle = |1\rangle$ 

$$Z = |0\rangle\langle 0| - |1\rangle\langle 1|$$

# Exercise 2.12

Prove that the matrix

$$\begin{bmatrix} 1 & 0 \\ 1 & 1 \end{bmatrix}$$

is not diagonalizable.

### Solution

Concepts Involved: Linear Algebra, Eigenvalues, Eigenvectors, Diagonalization.

Solving for the eigenvalues of the matrix, we have:

$$\det\begin{bmatrix} 1 - \lambda & 0 \\ 1 & 1 - \lambda \end{bmatrix} = 0 \implies (1 - \lambda)^2 = 0 \implies \lambda_1, \lambda_2 = 1$$

But since the eigenvalue 1 is degenerate, the matrix only has one eigenvector; it therefore cannot be diagonalized.  $\Box$ 

### Exercise 2.13

If  $|w\rangle$  and  $|v\rangle$  are any two vectors, show that  $(|w\rangle\langle v|)^{\dagger} = |v\rangle\langle w|$ .

### Solution

Concepts Involved: Linear Algebra, Adjoints.

We observe that:

$$((|w\rangle\langle v|)^{\dagger}|x\rangle, |y\rangle) = (|x\rangle, (|w\rangle\langle v|)|y\rangle) = (|x\rangle, \langle v|y\rangle |w\rangle) = \langle x|\langle v|y\rangle |w\rangle$$

$$= \langle x|w\rangle \langle v|y\rangle$$

$$= \langle x|w\rangle (|v\rangle, |y\rangle)$$

$$= (\langle x|w\rangle^* |v\rangle, |y\rangle)$$

$$= (\langle w|x\rangle |v\rangle, |y\rangle)$$

$$= ((|v\rangle\langle w|)|x\rangle, |y\rangle)$$

Where in the third-to last equality we use the conjugate linearity in the first argument (see Exercise 2.6) and in the second-to last equality we use that  $\langle a|b\rangle^*=\langle b|a\rangle$ . Comparing the first and last expressions, we conclude that  $(|w\rangle\langle v|)^\dagger=|v\rangle\langle w|$ .

### Exercise 2.14: Anti-linearity of the adjoint

Show that the adjoint operator is anti-linear,

$$\left(\sum_{i} a_{i} A_{i}\right)^{\dagger} = \sum_{i} a_{i}^{*} A_{i}^{\dagger}.$$

### Solution

Concepts Involved: Linear Algebra, Adjoints.

We observe that:

$$\left(\left(\sum_{i} a_{i} A_{i}\right)^{\dagger} |a\rangle, |b\rangle\right) = \left(|a\rangle, \sum_{i} a_{i} A_{i} |b\rangle\right) = \sum_{i} a_{i} \left(|a\rangle, A_{i} |b\rangle\right) \\
= \sum_{i} a_{i} \left(A_{i}^{\dagger} |a\rangle, |b\rangle\right) \\
= \left(\sum_{i} a_{i}^{*} A_{i}^{\dagger} |a\rangle, |b\rangle\right)$$

where we invoke the definition of the adjoint in the first and third equalities, the linearity in the second argument in the second equality, and the conjugate linearity in the first argument in the last equality. The claim is proven by comparing the first and last expressions.

### Exercise 2.15

Show that  $(A^{\dagger})^{\dagger} = A$ .

#### Solution

Concepts Involved: Linear Algebra, Adjoints

Applying the definition of the Adjoint twice (and using the conjugate symmetry of the inner product) we have that:

$$((A^{\dagger})^{\dagger}|a\rangle,|b\rangle) = (|a\rangle,A^{\dagger}|b\rangle) = (A^{\dagger}|b\rangle,|a\rangle)^* = (|b\rangle,A|a\rangle)^* = ((A|a\rangle,|b\rangle)^*)^* = (A|a\rangle,|b\rangle).$$

The claim follows by comparison of the first and last expressions.

#### Exercise 2.16

Show that any projector P satisfies the equation  $P^2 = P$ .

#### Solution

Concepts Involved: Linear Algebra, Projectors.

Let  $|1\rangle, \dots, |k\rangle$  be an orthonormal basis for the subspace W of V. Then, using the definition of the projector onto W, we have that:

$$P^2 = P \cdot P = \left(\sum_{i=1}^k |i\rangle\langle i|\right) \left(\sum_{i'=1}^k |i'\rangle\langle i'|\right) = \sum_{i=1}^k \sum_{i'=1}^k |i\rangle\langle i|i'\rangle\langle i'| = \sum_{i=1}^k \sum_{i'=1}^k |i\rangle\delta_{ii'}\langle i'| = \sum_{i=1}^k |i\rangle\langle i| = P$$

where in the fourth/fifth equality we use the orthonormality of the basis to collapse the double sum.

### Exercise 2.17

Show that a normal matrix is Hermitian if and only if it has real eigenvalues.

#### Solution

Concepts Involved: Linear Algebra, Hermitian Operators, Normal Operators, Spectral Decomposition.

Let A be a Normal and Hermitian matrix. Then, it has a diagonal representation  $A = \sum_i \lambda_i |i\rangle\langle i|$  where  $|i\rangle$  is an orthonormal basis for V and each  $|i\rangle$  is an eigenvector of A with eigenvalue  $\lambda_i$ . By the

Hermicity of A, we have that  $A = A^{\dagger}$ . Therefore, we have that:

$$A^{\dagger} = \left(\sum_{i} \lambda_{i} |i\rangle\langle i|\right)^{\dagger} = \sum_{i} \lambda_{i}^{*} |i\rangle\langle i| = A = \sum_{i} \lambda_{i} |i\rangle\langle i|$$

where we use the results of Exercises 2.13 and 2.14 in the second equality. Comparing the third and last expressions, we have that  $\lambda_i = \lambda_i *$  and hence the eigenvalues are real.

Let A be a Normal matrix with real eigenvalues. Then, A has diagonal representation  $A = \sum_i \lambda_i |i\rangle\langle i|$  where  $\lambda_i$  are all real. We therefore have that:

$$A^\dagger = \left(\sum_i \lambda_i \ket{i}\!\!ra{i}
ight)^\dagger = \lambda_i^* \ket{i}\!\!ra{i} = \sum_i \lambda_i \ket{i}\!\!ra{i} = A$$

where in the third equality we use that  $\lambda_i^* = \lambda_i$ . We conclude that A is Hermitian.

# Exercise 2.18

Show that all eigenvalues of a unitary matrix have modulus 1, that is, can be written in the form  $e^{i\theta}$  for some real  $\theta$ .

### Solution

Concepts Involved: Linear Algebra, Unitary Operators, Spectral Decomposition

Let U be a unitary matrix. It is then normal as  $U^\dagger U = U^\dagger U = I$ . It therefore has spectral decomposition  $U = \sum_i \lambda_i \, |i\rangle\langle i|$  where  $|i\rangle$  is an orthonormal basis of V, and  $|i\rangle$  are the eigenvectors of U with eigenvalues  $\lambda_i$ . We then have that:

$$UU^{\dagger} = I \implies \left(\sum_{i} \lambda_{i} |i\rangle\langle i|\right) \left(\sum_{i'} \lambda_{i'} |i'\rangle\langle i'|\right)^{\dagger} = I$$

$$\implies \left(\sum_{i} \lambda_{i} |i\rangle\langle i|\right) \left(\sum_{i'} \lambda_{i'}^{*} |i'\rangle\langle i'|\right) = I$$

$$\implies \sum_{i} \sum_{i'} \lambda_{i} \lambda_{i'} |i\rangle\langle i|i'\rangle\langle i'| = I$$

$$\implies \sum_{i} \sum_{i'} \lambda_{i} \lambda_{i'}^{*} |i\rangle\delta_{ii'}\langle i'| = I$$

$$\implies \sum_{i} \lambda_{i} \lambda_{i'} |i\rangle\langle i| = I$$

$$\implies \sum_{i} |\lambda_{i}|^{2} |i\rangle\langle i| = \sum_{i} 1 |i\rangle\langle i|$$

From which we obtain that  $|\lambda_i|^2 = 1$ , and hence  $|\lambda_i| = 1$ , proving the claim.

# Exercise 2.19: Pauli matrices: Hermitian and unitary

Show that the Pauli matrices are Hermitian and unitary.

### Solution

Concepts Involved: Linear Algebra, Hermitian Matrices, Unitary Matrices

We check I, X, Y, Z in turn.

$$I^\dagger = \left( \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix}^T \right)^* = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix}^* = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix} = I$$

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

$$X^{\dagger} = \left( \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix}^T \right)^* = \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix}^* = \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} = X$$

$$X^{\dagger}X = XX = \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix} = I$$

$$Y^{\dagger} = \left( \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix}^T \right)^* = \begin{bmatrix} 0 & i \\ -i & 0 \end{bmatrix}^* = \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} = Y$$

$$Y^{\dagger}Y = YY = \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix}$$

$$Z^{\dagger} = \left( \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix}^T \right)^* = \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix}^* = \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} = Z$$

$$Z^{\dagger}Z = \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix} = I$$

### Exercise 2.20: Basis changes

Suppose A' and A'' are matrix representations of an operator A on a vector space A on a vector space V with respect to two different orthonormal bases,  $|v_i\rangle$  and  $|w_i\rangle$ . Then the elements of A' and A'' are  $A'_{ij} = \langle v_i | A | v_j \rangle$  and  $A''_{ij} = \langle w_i | A | w_j \rangle$ . Characterize the relationship between A' and A''.

### Solution

Concepts Involved: Linear Algebra, Matrix Representations of Operators, Completeness Relation

Using the completeness relation twice, we get:

$$A'_{ij} = \langle v_i | A | v_j \rangle = \langle v_i | IAI | v_j \rangle = \langle v_i | \left( \sum_{i'} |w_{i'}\rangle \langle w_{i'}| \right) A \left( \sum_{j'} |w_{j'}\rangle \langle w_{j'}| \right) |v_j\rangle$$

$$= \sum_{i'} \sum_{j'} \langle v_i | w_{i'}\rangle \langle w_{i'}| A |w_{j'}\rangle \langle w_{j'}|v_j\rangle$$

$$= \sum_{i'} \sum_{j'} \langle v_i | w_{i'}\rangle A''_{ij} \langle w_{j'}|v_j\rangle$$

### Exercise 2.21

Repeat the proof of the spectral decomposition in Box 2.2 for the case when M is Hermitian, simplifying the proof wherever possible.

### Solution

Concepts Involved: Linear Algebra, Hermitian Operators, Spectral Decomposition.

Note that the converse of Theorem 2.1 does not hold if we replace "normal" with "Hermitian". Diagonalizability does not imply Hermicity, with a concrete example being  $S=|0\rangle\langle 0|+i\,|1\rangle\langle 1|$ . So, we just prove the forwards direction.

We proceed by induction on the dimension d of V. The d=1 case is trivial as M is already diagonal in any representation in this case. Let  $\lambda$  be an eigenvalue of M, P the projector onto the  $\lambda$  subspace, and Q the projector onto the orthogonal complement. Then M=(P+Q)M(P+Q)=PMP+QMP+PMQ+QMQ. Obviously  $PMP=\lambda P$ . Furthermore, QMP=0, as M takes the subspace P into itself. We claim that PMQ=0 also. To see this, we recognize that  $(PMQ)^\dagger=Q^\dagger M^\dagger P^\dagger=QMP=0$ . and hence PMQ=0. Thus M=PMP+QMQ. QMQ is normal, as  $(QMQ)^\dagger=Q^\dagger M^\dagger Q^\dagger=QMQ$  (and Hermiticity implies that the operator is normal). By induction, QMQ is diagonal with respect to some orthonormal basis for the subspace Q, and PMP is already diagonal with respect to some orthonormal basis for P. It follows that M=PMP+QMQ is diagonal with respect to some orthonormal basis for the total vector space.

#### Exercise 2.22

Prove that two eigenvectors of a Hermitian operator with different eigenvalues are necessarily orthogonal.

#### Solution

Concepts Involved: Linear Algebra, Eigenvalues, Eigenvectors, Hermitian Operators.

Let A be a Hermitian operator, and let  $|v_1\rangle, |v_2\rangle$  be two eigenvectors of A with corresponding eigenvalues  $\lambda_1, \lambda_2$  such that  $\lambda_1 \neq \lambda_2$ . We then have that:

$$\langle v_1 | A | v_2 \rangle = \langle v_1 | \lambda_2 | v_2 \rangle = \lambda_2 \langle v_1 | v_2 \rangle$$
$$\langle v_1 | A | v_2 \rangle = \langle v_1 | A^{\dagger} | v_2 \rangle = \langle v_1 | \lambda_1 | v_2 \rangle = \lambda_1 \langle v_1 | v_2 \rangle$$

where we use the Hermiticity of A in the second line. Substracting the first line from the second, we have that:

$$0 = (\lambda_2 - \lambda_1) \langle v_1 | v_2 \rangle.$$

Since  $\lambda_1 \neq \lambda_2$  by assumption, the only way this equality is satisfied is if  $\langle v_1|v_2\rangle=0$ . Hence,  $|v_1\rangle, |v_2\rangle$  are orthogonal.

# Exercise 2.23

Show that the eigenvalues of a projector P are all either 0 or 1.

### Solution

Concepts Involved: Linear Algebra, Eigenvalues, Eigenvectors, Projectors.

Let P be a projector, and  $|v\rangle$  be an eigenvector of P with corresponding eigenvalue  $\lambda$ . From Exercise 2.16, we have that  $P^2=P$ , and using this fact, we observe:

$$\begin{split} P|v\rangle &= \lambda|v\rangle \\ P|v\rangle &= P^2|v\rangle = PP|v\rangle = P\lambda|v\rangle = \lambda P|v\rangle = \lambda^2|v\rangle. \end{split}$$

Subtracting the first line from the second, we get:

$$\mathbf{0} = (\lambda^2 - \lambda)|v\rangle = \lambda(\lambda - 1)|v\rangle.$$

Since  $|v\rangle$  is not the zero vector, we therefore obtain that either  $\lambda=0$  or  $\lambda=1$ .

### Exercise 2.24: Hermiticity of positive operator

Show that a positive operator is necessarily Hermitian. (Hint: Show that an arbitrary operator A can be written A=B+iC where B and C are Hermitian.)

#### **Solution**

Concepts Involved: Linear Algebra, Hermitian Operators, Positive Operators

Let A be an operator. We first make the observation that we can write A as:

$$A = \frac{A}{2} + \frac{A}{2} + \frac{A^{\dagger}}{2} - \frac{A^{\dagger}}{2} = \frac{A + A^{\dagger}}{2} + i\frac{A - A^{\dagger}}{2i}.$$

So let  $B=\frac{A+A^{\dagger}}{2}$  and  $C=\frac{A-A^{\dagger}}{2i}.$  B and C are Hermitian, as:

$$B^\dagger = \left(\frac{A+A^\dagger}{2}\right)^\dagger = \frac{A^\dagger + (A^\dagger)^\dagger}{2} = \frac{A^\dagger + A}{2} = B$$

$$C^{\dagger} = \left(\frac{A - A^{\dagger}}{2i}\right)^{\dagger} = \frac{A^{\dagger} - (A^{\dagger})^{\dagger}}{-2i} = \frac{A - A^{\dagger}}{2i} = C$$

so we have hence proven that we can write A=B+iC for hermitian B,C for any operator A. Now, assume that A is positive. We then have that for any vector  $|v\rangle$ :

$$\langle v|A|v\rangle \geq 0.$$

Using the identity derived above, we have that:

$$\langle v|B|v\rangle + i\langle v|C|v\rangle \ge 0.$$

The positivity forces C=0. Therefore, A=B and hence A is Hermitian.

### Exercise 2.25

Show that for any operator  $A, A^{\dagger}A$  is positive.

#### Solution

Concepts Involved: Linear Algebra, Adjoints, Positive Operators

Let A be an operator. Let  $|v\rangle$  be an arbitrary vector, and then we then have that:

$$\left(|v\rangle,A^{\dagger}A|v\rangle\right)=\left((A^{\dagger})^{\dagger}|v\rangle,A|v\rangle\right)=\left(A|v\rangle,A|v\rangle\right).$$

By the property of inner products, the expression must be greater than zero.

# Exercise 2.26

Let  $|\psi\rangle=(|0\rangle+|1\rangle)/\sqrt{2}$ . Write out  $|\psi\rangle^{\otimes 2}$  and  $|\psi\rangle^{\otimes 3}$  explicitly, both in terms of tensor products like  $|0\rangle|1\rangle$  and using the Kronecker product.

# Solution

Concepts Involved: Linear Algebra, Tensor Products, Kronecker Products.

Using the definition of the tensor product, we have:

$$|\psi\rangle^{\otimes 2} = \frac{|0\rangle|0\rangle + |0\rangle|1\rangle + |1\rangle|0\rangle + |1\rangle|1\rangle}{2} \cong \begin{bmatrix} \frac{1}{\sqrt{2}} \\ \frac{1}{\sqrt{2}} \\ \frac{1}{\sqrt{2}} \\ \frac{1}{\sqrt{2}} \\ \frac{1}{\sqrt{2}} \end{bmatrix} = \begin{bmatrix} \frac{1}{2} \\ \frac{1}{2} \\ \frac{1}{2} \\ \frac{1}{2} \\ \frac{1}{2} \end{bmatrix}$$

$$= |\psi\rangle \otimes |\psi\rangle^{\otimes 2} \cong \begin{bmatrix} \frac{1}{2} \\ \frac{$$

# Exercise 2.27

Calculate the matrix representation of the tensor products of the Pauli operators (a) X and Z; (b) I and X; (c) X and I. Is the tensor product commutative?

### Solution

Concepts Involved: Linear Algebra, Tensor Products, Kronecker Products.

Using the Kronecker product, we have:

(a)

$$X \otimes Z = \begin{bmatrix} 0Z & 1Z \\ 1Z & 0Z \end{bmatrix} = \begin{bmatrix} 0 \begin{bmatrix} 1 & 0 \\ 0 & -1 \\ 1 & 0 \\ 0 & -1 \end{bmatrix} & 1 \begin{bmatrix} 1 & 0 \\ 0 & -1 \\ 1 & 0 \\ 0 & -1 \end{bmatrix} \end{bmatrix} = \begin{bmatrix} 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & -1 \\ 1 & 0 & 0 & 0 \\ 0 & -1 & 0 & 0 \end{bmatrix}$$

(b)

$$I \otimes X = \begin{bmatrix} 1X & 0X \\ 0X & 1X \end{bmatrix} = \begin{bmatrix} 1 \begin{bmatrix} 0 & 1 \\ 1 & 0 \\ 0 & 1 \\ 1 & 0 \end{bmatrix} & 0 \begin{bmatrix} 0 & 1 \\ 1 & 0 \\ 0 & 1 \\ 1 & 0 \end{bmatrix} \end{bmatrix} = \begin{bmatrix} 0 & 1 & 0 & 0 \\ 1 & 0 & 0 & 0 \\ 0 & 0 & 0 & 1 \\ 0 & 0 & 1 & 0 \end{bmatrix}$$

(c)

$$X \otimes I = \begin{bmatrix} 0I & 1I \\ 1I & 0I \end{bmatrix} = \begin{bmatrix} 0 \begin{bmatrix} 1 & 0 \\ 0 & 1 \\ 1 & 0 \\ 0 & 1 \end{bmatrix} & 1 \begin{bmatrix} 1 & 0 \\ 0 & 1 \\ 1 & 0 \\ 0 & 1 \end{bmatrix} \end{bmatrix} \begin{bmatrix} 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \\ 0 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \end{bmatrix}$$

Comparing (b) and (c), we conclude that the tensor product is not commutative.

### Exercise 2.28

Show that the transpose, complex conjugation and adjoint operations distribute over the tensor product,

$$(A \otimes B)^* = A^* \otimes B^*; (A \otimes B)^T = A^T \otimes B^T; (A \otimes B)^{\dagger} = A^{\dagger} \otimes B^{\dagger}.$$

### Solution

Concepts Involved: Linear Algebra, Adjoints, Tensor Products, Kronecker Products.

Using the Kronecker product representation of  $a \otimes B$ , we have:

$$(A \otimes B)^* = \begin{bmatrix} A_{11}B & A_{12}B & \dots & A_{1n}B \\ A_{21}B & A_{22}B & \dots & A_{2n}B \\ \vdots & \vdots & \vdots & \vdots \\ A_{m1}B & A_{m2}B & \dots & A_{mn}B \end{bmatrix}^* = \begin{bmatrix} A_{11}^*B^* & A_{12}^*B^* & \dots & A_{1n}^*B^* \\ A_{21}^*B^* & A_{22}^*B^* & \dots & A_{2n}^*B^* \\ \vdots & \vdots & \vdots & \vdots \\ A_{m1}^*B^* & A_{m2}^*B^* & \dots & A_{mn}^*B^* \end{bmatrix} = A^* \otimes B^*$$

$$(A \otimes B)^T = \begin{bmatrix} A_{11}B & A_{12}B & \dots & A_{1n}B \\ A_{21}B & A_{22}B & \dots & A_{2n}B \\ \vdots & \vdots & \vdots & \vdots \\ A_{m1}B & A_{m2}B & \dots & A_{mn}B \end{bmatrix}^T = \begin{bmatrix} A_{11}B^T & A_{21}B^T & \dots & A_{n1}B^T \\ A_{12}B^T & A_{22}B & \dots & A_{n2}B^T \\ \vdots & \vdots & \vdots & \vdots \\ A_{1m}B^T & A_{2m}B^T & \dots & A_{nm}B^T \end{bmatrix} = A^T \otimes B^T.$$

The relation for the distributivity of the hermitian conjugate over the tensor product then follows from the former two relations:

$$(A \otimes B)^{\dagger} = ((A \otimes B)^T)^* = (A^T \otimes B^T)^* = (A^T)^* \otimes (B^T)^* = A^{\dagger} \otimes B^{\dagger}$$

### Exercise 2.29

Show that the tensor product of two unitary operators is unitary.

#### Solution

Concepts Involved: Linear Algebra, Unitary Operators, Tensor Products

Suppose A,B are unitary. Then,  $A^{\dagger}A=I$  and  $B^{\dagger}B=I$ . Using the result of the Exercise 2.28, we then have that:

$$(A \otimes B)^{\dagger}(A \otimes B) = (A^{\dagger} \otimes B^{\dagger})(A \otimes B) = (A^{\dagger}A \otimes B^{\dagger}B) = I \otimes I$$

# Exercise 2.30

Show that the tensor product of two Hermitian operators is Hermitian.

#### Solution

Concepts Involved: Linear Algebra, Hermitian Operators, Tensor Products

Suppose A,B are Hermitian. Then,  $A^{\dagger}=A$  and  $B^{\dagger}=B$ . Then, using the result of Exercise 2.28, we have:

$$(A \otimes B)^{\dagger} = A^{\dagger} \otimes B^{\dagger} = A \otimes B$$

# Exercise 2.31

Show that the tensor product of two positive operators is positive.

### Solution

Concepts Involved: Linear Algebra, Positive Operators

Suppose A,B are positive operators. We then have that  $\langle v|A|v\rangle \geq 0$  and  $\langle w|B|w\rangle \geq 0$ . Therefore, for any  $|v\rangle \otimes |w\rangle$ :

$$(|v\rangle \otimes |w\rangle, A \otimes B(|v\rangle \otimes |w\rangle)) = \langle v|A|v\rangle \langle w|B|w\rangle \ge 0$$

#### Exercise 2.32

Show that the tensor product of two projectors is a projector.

#### Solution

Concepts Involved: Linear Algebra, Projectors

Let  $P_1, P_2$  be projectors. We then have that  $P_1^2 = P_1$  and  $P_2^2 = P_2$  by Exercise 2.16. Therefore:

$$(P_1 \otimes P_2)^2 = (P_1 \otimes P_2)(P_1 \otimes P_2) = P_1^2 \otimes P_2 = P_1 \otimes P_2$$

so  $P_1 \otimes P_2$  is a projector.

#### Exercise 2.33

The Hadamard operator on one qubit may be written as

$$H = \frac{1}{\sqrt{2}}[(|0\rangle + |1\rangle)\langle 0| + (|0\rangle - |1\rangle)\langle 1|]$$

Show explicitly that the Hadamard transform on n qubits,  $H^{\otimes n}$ , may be written as

$$H^{\otimes n} = \frac{1}{\sqrt{2^n}} \sum_{x,y} (-1)^{x \cdot y} |x\rangle \langle y|$$

Write out an explicit matrix representation for  $H^{\otimes 2}$ 

#### Solution

Concepts Involved: Linear algebra, Matrix Representation of Operators, Outer Products.

Looking at the form of the Hadamard operator on one qubit, we observe that:

$$H = \frac{1}{\sqrt{2}} \left[ |0\rangle\langle 0| + |0\rangle\langle 1| + |1\rangle\langle 0| - |1\rangle\langle 1| \right]$$

Hence:

$$H = \frac{1}{\sqrt{2}} \sum_{x,y} (-1)^{x \cdot y} |x\rangle \langle y|$$

Where x, y run over 0 and 1. Taking the n-fold tensor product of this expression, we get:

$$H^{\otimes} = \frac{1}{\sqrt{2}} \sum_{x,y} (-1)^{x \cdot y} |x\rangle \langle y| \otimes \frac{1}{\sqrt{2}} \sum_{x,y} (-1)^{x \cdot y} |x\rangle \langle y| \otimes \dots \otimes \frac{1}{\sqrt{2}} \sum_{x,y} (-1)^{x \cdot y} |x\rangle \langle y|$$
$$= \frac{1}{\sqrt{2^n}} \sum_{\mathbf{x}, \mathbf{y}} (-1)^{\mathbf{x} \cdot \mathbf{y}} |\mathbf{x}\rangle \langle \mathbf{y}|$$

Where x, y are length n-binary strings. This proves the claim.

Now explicitly writing  $H^{\otimes 2}$ , we have:

Note that here,  $\mathbf{x}, \mathbf{y}$  are binary length 2 strings. The sum goes through all pairwise combinations of  $\mathbf{x}, \mathbf{y} \in \{00, 01, 10, 11\}$ .

**Remark:** Sylvester's Construction gives an interesting recursive construction of Hadamard matrices. See https://en.wikipedia.org/wiki/Hadamard\_matrix. Discussion on interesting (related) open problem concerning the maximal determinant of matrices consisting of entries of 1 and -1 can be found here https://en.wikipedia.org/wiki/Hadamard%27s\_maximal\_determinant\_problem.

### Exercise 2.34

Find the square root and logarithm of the matrix

$$\begin{bmatrix} 4 & 3 \\ 3 & 4 \end{bmatrix}$$

# Solution

Concepts Involved: Linear Algebra, Spectral Decomposition, Operator Functions

We begin by diagonalizing the matrix (which we call A) as to be able to apply the definition of operator functions. By inspection, A is Hermitian as it is equal to its conjugate transpose, so the spectral decomposition exists. Solving for the eigenvalues, we consider the characteristic equation:

$$\det(A - \lambda I) = 0 \implies \det\begin{bmatrix} 4 - \lambda & 3 \\ 3 & 4 - \lambda \end{bmatrix} = 0 \implies (4 - \lambda)^2 - 9 = 0 \implies \lambda^2 - 8\lambda + 7 = 0$$

Using the quadratic equation, we get  $\lambda_1=1, \lambda_2=7$ . Using this to find the eigenvectors of the matrix, we have:

$$\begin{bmatrix} 4-1 & 3 \\ 3 & 4-1 \end{bmatrix} \begin{bmatrix} v_{11} \\ v_{12} \end{bmatrix} = \mathbf{0} \implies v_{11} = 1, v_{12} = -1$$

$$\begin{bmatrix} 4-7 & 3 \\ 3 & 4-7 \end{bmatrix} \begin{bmatrix} v_{21} \\ v_{22} \end{bmatrix} = \mathbf{0} \implies v_{21} = 1, v_{22} = 1$$

Hence our normalized eigenvectors are:

$$|v_1\rangle = \begin{bmatrix} \frac{1}{\sqrt{2}} \\ \frac{-1}{\sqrt{2}} \end{bmatrix}, \quad |v_2\rangle = \begin{bmatrix} \frac{1}{\sqrt{2}} \\ \frac{1}{\sqrt{2}} \end{bmatrix}$$

Therefore the spectral composition of the matrix is given by:

$$A = 1 |v_1\rangle\langle v_1| + 7 |v_2\rangle\langle v_2|$$

Calculating the square root of A, we then have:

$$\sqrt{A} = \sqrt{1} |v_1\rangle\langle v_1| + \sqrt{7} |v_2\rangle\langle v_2| = \frac{1}{2} \begin{bmatrix} 1 & -1 \\ -1 & 1 \end{bmatrix} + \frac{\sqrt{7}}{2} \begin{bmatrix} 1 & 1 \\ 1 & 1 \end{bmatrix} = \frac{1}{2} \begin{bmatrix} 1 + \sqrt{7} & -1 + \sqrt{7} \\ -1 + \sqrt{7} & 1 + \sqrt{7} \end{bmatrix}.$$

Calculating the logarithm of A, we have:

$$\log(A) = \log(1) |v_1\rangle\langle v_1| + \log(7) |v_2\rangle\langle v_2| = \frac{\log(7)}{2} \begin{bmatrix} 1 & 1\\ 1 & 1 \end{bmatrix}$$

# **Exercise 2.35: Exponential of Pauli matrices**

Let  ${f v}$  be any real, three-dimensional unit vector and  ${f heta}$  a real number. Prove that

$$\exp(i\theta \mathbf{v} \cdot \boldsymbol{\sigma}) = \cos(\theta)I + i\sin(\theta)\mathbf{v} \cdot \boldsymbol{\sigma}$$

Where  $\mathbf{v}\cdot\boldsymbol{\sigma}\equiv\sum_{i=1}^3v_i\sigma_i$ . This exercise is generalized in Problem 2.1 on page 117.

#### Solution

Concepts Involved: Linear Algebra, Spectral Decomposition, Operator Functions.

Recall that  $\sigma_1 \equiv X, \sigma_2 \equiv Y$ , and  $\sigma_3 \equiv Z$ .

First, we compute  $\mathbf{v} \cdot \boldsymbol{\sigma}$  in matrix form:

$$\mathbf{v} \cdot \boldsymbol{\sigma} = v_1 \sigma_1 + v_2 \sigma_2 + v_3 \sigma_3 = v_1 \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} + v_2 \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} + v_3 \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} = \begin{bmatrix} v_3 & v_1 - iv_2 \\ v_1 + iv_2 & -v_3 \end{bmatrix}$$

In order to compute the complex exponential of this matrix, we will want to find its spectral decomposition.

Using the characterstic equation to find the eigenvalues, we have:

$$\det(\mathbf{v} \cdot \boldsymbol{\sigma} - I\lambda) = 0 \implies \det \begin{bmatrix} v_3 - \lambda & v_1 - iv_2 \\ v_1 + iv_2 & -v_3 - \lambda \end{bmatrix} = 0$$

$$\implies (v_3 - \lambda)(-v_3 - \lambda) - (v_1 - iv_2)(v_1 + iv_2) = 0$$

$$\implies \lambda^2 - v_3^2 - v_1^2 - v_2^2 = \lambda^2 - (v_1^2 + v_2^2 + v_3^2) = 0$$

$$\implies \lambda^2 - 1 = 0$$

$$\implies \lambda_1 = 1, \lambda_2 = -1$$

where in the second-to-last implication we use the fact that  ${\bf v}$  is a unit vector. Letting  $|v_1\rangle, |v_2\rangle$  be the associated eigenvectors,  ${\bf v} \cdot {\boldsymbol \sigma}$  has spectral decomposition:

$$\mathbf{v} \cdot \boldsymbol{\sigma} = |v_1\rangle\langle v_1| - |v_2\rangle\langle v_2|$$

Applying the complex exponentiation operator, we then have that:

$$\exp(i\theta \mathbf{v} \cdot \boldsymbol{\sigma}) = \exp(i\theta) |v_1\rangle\langle v_1| + \exp(-i\theta) |v_2\rangle\langle v_2|.$$

Using Euler's formula, we then have that:

$$\exp(i\theta \mathbf{v} \cdot \boldsymbol{\sigma}) = (\cos \theta + i \sin \theta) |v_1\rangle\langle v_1| + (\cos \theta - i \sin \theta) |v_2\rangle\langle v_2|$$
$$= \cos(\theta) (|v_1\rangle\langle v_1| + |v_2\rangle\langle v_2|) + i \sin(\theta) (|v_1\rangle\langle v_1| - |v_2\rangle\langle v_2|)$$
$$= \cos(\theta) I + i \sin(\theta) \mathbf{v} \cdot \boldsymbol{\sigma}.$$

Where in the last line we use the completeness relation and the spectral decomposition.

# Exercise 2.36

Show that the Pauli matrices except for I have trace zero.

### Solution

Concepts Involved: Linear Algebra, Trace.

We have that:

$$tr(X) = tr \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} = 0 + 0 = 0$$

$$tr(Y) = tr \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} = 0 + 0 = 0$$

$$tr(Z) = tr \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} = 1 - 1 = 0$$

# Exercise 2.37: Cyclic property of the trace

If A and B are two linear operators show that

$$tr(AB) = tr(BA)$$

#### Solution

Concepts Involved: Linear Algebra, Trace.

Let A, B be linear operators. Then, C = AB has matrix representation with entries  $C_{ij} = \sum_k A_{ik} B_{kj}$  and D = BA has matrix representation with entries  $D_{ij} = \sum_k B_{ik} A_{kj}$ . We then have that:

$$\operatorname{tr}(AB) = \operatorname{tr}(C) = \sum_{i} C_{ii} = \sum_{i} \sum_{k} A_{ik} B_{ki} = \sum_{k} \sum_{i} B_{ki} A_{ik} = \sum_{k} D_{kk} = \operatorname{tr}(D) = \operatorname{tr}(BA)$$

# Exercise 2.38: Linearity of the trace

If A and B are two linear operators, show that

$$tr(A+B) = tr(A) + tr(B)$$

and if z is an arbitrary complex number show that

$$\operatorname{tr}(zA) = z \operatorname{tr}(A).$$

#### Solution

Concepts Involved: Linear Algebra, Trace.

From the definition of trace, we have that:

$$\operatorname{tr}(A+B) = \sum_{i} (A+B)_{ii} = \sum_{i} A_{ii} + B_{ii} = \sum_{i} A_{ii} + \sum_{i} B_{ii} = \operatorname{tr}(A) + \operatorname{tr}(B)$$

$$\operatorname{tr}(zA) = \sum_{i} (zA)_{ii} = z \sum_{i} A_{ii} = z \operatorname{tr}(A)$$

### Exercise 2.39: The Hilbert-Schmidt inner product on operators

The set  $L_V$  of linear operators on Hilbert space V is obviously a vector space - the sum of two linear operators is a linear operator, zA is a linear operator if A is a linear operator and z is a complex number, and there is a zero element 0. An important additional result is that the vector space  $L_V$  can be given a natural inner product structure, turning it into a Hilbert space.

(1) Show that the function  $(\cdot,\cdot)$  on  $L_V \times L_V$  defined by

$$(A,B) \equiv \operatorname{tr}(A^{\dagger}B)$$

is an inner product function. This inner product is known as the *Hilbert-Schmidt* or *trace* inner product.

- (2) If V has d dimensions show that  $L_V$  has dimension  $d^2$ .
- (3) Find an orthonormal basis of Hermitian matrices for the Hilbert space  $L_V$ .

#### Solution

Concepts Involved: Linear Algebra, Trace, Inner Products, Hermitian Operators, Bases

(1) We show that  $(\cdot, \cdot)$  satisfies the three properties of an inner product. Showing that it is linear in the second argument, we have that:

$$\left(A, \sum_{i} \lambda_{i} B_{i}\right) = \operatorname{tr}\left(A \sum_{i} \lambda_{i} B_{i}\right) = \sum_{i} \lambda_{i} \operatorname{tr}(A B_{i}) = \sum_{i} \lambda_{i} (A, B_{i})$$

where in the second to last equality we use the result of Exercise 2.38. To see that it is conjugate-symmetric, we have that:

$$(A,B)=\operatorname{tr}\Bigl(A^\dagger B\Bigr)=\operatorname{tr}\Bigl((B^\dagger A)^\dagger\Bigr)=\operatorname{tr}\Bigl(B^\dagger A\Bigr)^*=(B,A)^*$$

Finally, to show positive definitness, we have that:

$$(A, A) = \operatorname{tr}(A^{\dagger}A) = \sum_{i} \sum_{k} A_{ik}^{\dagger} A_{ki} = \sum_{i} \sum_{k} A_{ki}^{*} A_{ki} = \sum_{i} \sum_{k} |A_{ki}|^{2} \ge 0$$

so we conclude that  $(\cdot, \cdot)$  is an inner product function.

- (2) Suppose V has d dimensions. Then, the elements of  $L_V$  which consist of linear operators  $A:V\mapsto V$  have representations as  $d\times d$  matrices. There are  $d^2$  such linearly independent matrices (take the matrices with 1 in one of the  $d^2$  entries and 0 elsewhere), and we conclude that  $L_V$  has  $d^2$  linearly independent vectors and hence dimension  $d^2$ .
- (3) As discussed in the previous part of the question, one possible basis for this vector space would be  $|v_i\rangle\langle v_j|$  where  $|v_k\rangle$  form an orthonormal basis of V with  $i,j\in\{1,\ldots d\}$ . These of course are just matrices with 1 in one entry and 0 elsewhere. It is easy to see that this is a basis as for any

 $A \in L_V$  we can write  $A = \sum_{ij} \lambda_{ij} |v_i\rangle\langle v_j|$ . We can verify that these are orthonormal; suppose  $|v_{i_1}\rangle\langle v_{j_1}| \neq |v_{i_2}\rangle\langle v_{j_2}|$ . Then, we have that:

$$(|v_{i_1}\rangle\langle v_{j_1}|, |v_{i_2}\rangle\langle v_{j_2}|) = \operatorname{tr}((|v_{i_1}\rangle\langle v_{j_1}|)^{\dagger} |v_{i_2}\rangle\langle v_{j_2}|)$$
$$= \operatorname{tr}(|v_{j_1}\rangle\langle v_{i_1}|v_{i_2}\rangle\langle v_{j_2}|)$$

If  $|v_{i_1}\rangle \neq |v_{i_2}\rangle$ , then the above expression reduces to  $\operatorname{tr}(0)=0$ . If  $|v_{i_1}\rangle = |v_{i_2}\rangle$ , then it follows that  $|v_{j_1}\rangle \neq |v_{j_2}\rangle$  (else this would contradict  $|v_{i_1}\rangle\langle v_{j_1}| \neq |v_{i_2}\rangle\langle v_{j_2}|$ ) and in this case we have that:

$$(|v_{i_1}\rangle\langle v_{j_1}|, |v_{i_2}\rangle\langle v_{j_2}|) = \operatorname{tr}(|v_{j_1}\rangle\langle v_{i_1}|v_{i_2}\rangle\langle v_{j_2}|)$$
$$= \operatorname{tr}(|v_{j_1}\rangle\langle v_{j_2}|)$$
$$= 0$$

So we therefore have that the inner product of two non-identical elements in the basis is zero. Furthermore, we have that:

$$(|v_{i_1}\rangle\langle v_{j_1}|, |v_{i_1}\rangle\langle v_{j_1}|) = \operatorname{tr}(|v_{i_1}\rangle\langle v_{j_1}| |v_{i_1}\rangle\langle v_{j_1}|) = \operatorname{tr}(|v_{i_1}\rangle\langle v_{i_1}|) = 1$$

so we confirm that this basis is orthonormal. However, evidently this basis is *not* Hermitian as if  $i \neq j$ , then  $(|v_i\rangle\langle v_j|)^\dagger = |v_j\rangle\langle v_i| \neq |v_i\rangle\langle v_j|$ . To fix this, we can modify our basis slightly. We keep the diagonal entries as is (as these are indeed Hermitian!) but for the off-diagonals, we replace every pair of basis vectors  $|v_i\rangle\langle v_j|$ ,  $|v_j\rangle\langle v_i|$  with:

$$\frac{|v_i\rangle\langle v_j|+|v_j\rangle\langle v_i|}{\sqrt{2}}, \quad i\frac{|v_i\rangle\langle v_j|-|v_j\rangle\langle v_i|}{\sqrt{2}}.$$

A quick verification shows that these are indeed Hermitian:

$$\left(\frac{|v_i\rangle\langle v_j| + |v_j\rangle\langle v_i|}{\sqrt{2}}\right)^{\dagger} = \frac{(|v_i\rangle\langle v_j|)^{\dagger} + (|v_j\rangle\langle v_i|)^{\dagger}}{\sqrt{2}} = \frac{|v_i\rangle\langle v_j| + |v_j\rangle\langle v_i|}{\sqrt{2}} \\
\left(i\frac{|v_i\rangle\langle v_j| - |v_j\rangle\langle v_i|}{\sqrt{2}}\right)^{\dagger} = -i\frac{(|v_i\rangle\langle v_j|)^{\dagger} - (|v_j\rangle\langle v_i|)^{\dagger}}{\sqrt{2}} = i\frac{|v_i\rangle\langle v_j| - |v_j\rangle\langle v_i|}{\sqrt{2}}$$

It now suffices to show that these new vectors (plus the diagonals) form a basis and are orthonormal. To see that these form a basis, observe that:

$$\frac{1}{\sqrt{2}} \frac{|v_i\rangle\langle v_j| + |v_j\rangle\langle v_i|}{\sqrt{2}} - \frac{i}{\sqrt{2}} \left( i \frac{|v_i\rangle\langle v_j| - |v_j\rangle\langle v_i|}{\sqrt{2}} \right) = |v_i\rangle\langle v_j|$$

$$\frac{1}{\sqrt{2}} \frac{|v_i\rangle\langle v_j| + |v_j\rangle\langle v_i|}{\sqrt{2}} + \frac{i}{\sqrt{2}} \left( i \frac{|v_i\rangle\langle v_j| - |v_j\rangle\langle v_i|}{\sqrt{2}} \right) = |v_j\rangle\langle v_i|$$

and since we know that  $|v_i\rangle\langle v_j|$  for all  $i,j\in\{1,\ldots d\}$  form a basis, this newly defined set of vectors must be a basis as well. Furthermore, since the new basis vectors are constructed from orthogonal  $|v_i\rangle\langle v_j|$ , the newly defined vectors will be orthogonal to each other if  $i_1,j_1\neq i_2,j_2$ . The only things

left to check is that for any choice of i, j that:

$$\frac{|v_i\rangle\langle v_j|+|v_j\rangle\langle v_i|}{\sqrt{2}} \text{ and } i\frac{|v_i\rangle\langle v_j|-|v_j\rangle\langle v_i|}{\sqrt{2}}.$$

are orthogonal, and that these vectors are normalized. Checking the orthogonality, we have:

$$\left(\frac{|v_i\rangle\langle v_j| + |v_j\rangle\langle v_i|}{\sqrt{2}}, i\frac{|v_i\rangle\langle v_j| - |v_j\rangle\langle v_i|}{\sqrt{2}}\right) = \operatorname{tr}\left(\left(\frac{|v_i\rangle\langle v_j| + |v_j\rangle\langle v_i|}{\sqrt{2}}\right)\left(i\frac{|v_i\rangle\langle v_j| - |v_j\rangle\langle v_i|}{\sqrt{2}}\right)\right)$$

$$= \frac{i}{\sqrt{2}}\operatorname{tr}\left(|v_j\rangle\langle v_j| - |v_i\rangle\langle v_i|\right)$$

$$= 0.$$

And checking the normalization, we have that:

$$\left(\frac{|v_i\rangle\langle v_j| + |v_j\rangle\langle v_i|}{\sqrt{2}}, \frac{|v_i\rangle\langle v_j| + |v_j\rangle\langle v_i|}{\sqrt{2}}\right) = \operatorname{tr}\left(\left(\frac{|v_i\rangle\langle v_j| + |v_j\rangle\langle v_i|}{\sqrt{2}}\right)\left(\frac{|v_i\rangle\langle v_j| + |v_j\rangle\langle v_i|}{\sqrt{2}}\right)\right)$$

$$= \frac{1}{2}\operatorname{tr}\left(|v_i\rangle\langle v_i| + |v_j\rangle\langle v_j|\right)$$

$$= 1$$

$$\left(i\frac{|v_{i}\rangle\langle v_{j}| - |v_{j}\rangle\langle v_{i}|}{\sqrt{2}}, i\frac{|v_{i}\rangle\langle v_{j}| - |v_{j}\rangle\langle v_{i}|}{\sqrt{2}}\right) = \operatorname{tr}\left(\left(i\frac{|v_{i}\rangle\langle v_{j}| - |v_{j}\rangle\langle v_{i}|}{\sqrt{2}}\right)\left(i\frac{|v_{i}\rangle\langle v_{j}| - |v_{j}\rangle\langle v_{i}|}{\sqrt{2}}\right)\right)$$

$$= -\frac{1}{2}\operatorname{tr}\left(-|v_{i}\rangle\langle v_{i}| - |v_{j}\rangle\langle v_{j}|\right)$$

$$= 1$$

#### Exercise 2.40: Commutation relations for the Pauli matrices

Verify the commutation relations

$$[X, Y] = 2iZ;$$
  $[Y, Z] = 2iX;$   $[Z, X] = 2iY$ 

There is an elegant way of writing this using  $\epsilon_{jkl}$ , the antisymmetric tensor on three indices, for which  $\epsilon_{jkl}=0$  except for  $\epsilon_{123}=\epsilon_{231}=\epsilon_{312}=1$ , and  $\epsilon_{321}=\epsilon_{213}=\epsilon_{132}=-1$ :

$$[\sigma_j, \sigma_k] = 2i \sum_{l=1}^{3} \epsilon_{jkl} \sigma_l$$

#### Solution

Concepts Involved: Linear Algebra, Commutators.

We verify the proposed relations via computation in the computational basis:

$$[X,Y] = XY - YX = \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} - \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} = \begin{bmatrix} i & 0 \\ 0 & -i \end{bmatrix} - \begin{bmatrix} -i & 0 \\ 0 & i \end{bmatrix} = 2iZ$$
 
$$[Y,Z] = YZ - ZY = \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} - \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} = \begin{bmatrix} 0 & i \\ i & 0 \end{bmatrix} - \begin{bmatrix} 0 & -i \\ -i & 0 \end{bmatrix} = 2iX$$
 
$$[Z,X] = ZX - XZ = \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} - \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} = \begin{bmatrix} 0 & 1 \\ -1 & 0 \end{bmatrix} - \begin{bmatrix} 0 & -1 \\ 1 & 0 \end{bmatrix} = 2iY$$

#### Exercise 2.41: Anti-commutation relations for the Pauli matrices

Verify the anticommutation relations

$$\{\sigma_i, \sigma_j\} = 0$$

Where  $i \neq j$  are both chosen from the set 1, 2, 3. Also verify that (i = 0, 1, 2, 3)

$$\sigma_i^2 = I$$

#### Solution

Concepts Involved: Linear Algebra, Anticommutators.

We again verify the proposed relations via computation in the computational basis:

$$\{X,Y\} = XY + YX = \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} + \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} = \begin{bmatrix} i & 0 \\ 0 & -i \end{bmatrix} + \begin{bmatrix} -i & 0 \\ 0 & i \end{bmatrix} = \begin{bmatrix} 0 & 0 \\ 0 & 0 \end{bmatrix}$$

$$\{Y,Z\} = YZ + ZY = \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} + \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} = \begin{bmatrix} 0 & i \\ i & 0 \end{bmatrix} + \begin{bmatrix} 0 & -i \\ -i & 0 \end{bmatrix} = \begin{bmatrix} 0 & 0 \\ 0 & 0 \end{bmatrix}$$

$$\{Z,X\} = ZX + XZ = \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} + \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} = \begin{bmatrix} 0 & 1 \\ -1 & 0 \end{bmatrix} + \begin{bmatrix} 0 & -1 \\ 1 & 0 \end{bmatrix} = \begin{bmatrix} 0 & 0 \\ 0 & 0 \end{bmatrix} .$$

This proves the first claim as  $\{A,B\} = AB + BA = BA + AB = \{B,A\}$  and the other 3 relations are

equivalent to the ones already proven. Verifying the second claim, we have:

$$I^{2} = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix} = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix}$$

$$X^{2} = \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix}$$

$$Y^{2} = \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix}$$

$$Z^{2} = \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix}$$

**Remark:** Note that we can write this result consicely as  $\{\sigma_j,\sigma_k\}=2\delta_{ij}I$ 

# Exercise 2.42

Verify that

$$AB = \frac{[A, B] + \{A, B\}}{2}$$

# Solution

Concepts Involved: Linear Algebra, Commutators, Anticommutators.

By algebraic manipulation we obtain:

$$AB = \frac{AB + AB}{2} + \frac{BA - BA}{2} = \frac{(AB - BA) + (AB + BA)}{2} = \frac{[A, B] + \{A, B\}}{2}$$

#### Exercise 2.43

Show that for j, k = 1, 2, 3,

$$\sigma_j \sigma_k = \delta_{jk} I + i \sum_{l=1}^3 \epsilon_{jkl} \sigma_l.$$

#### Solution

Concepts Involved: Linear Algebra, Commutators, Anticommutators.

Applying the results of Exercises 2.40, 2.41, and 2.42, we have:

$$\sigma_{j}\sigma_{k} = \frac{[\sigma_{j}, \sigma_{k}] + \{\sigma_{j}, \sigma_{k}\}}{2}$$

$$= \frac{2i\sum_{l=1}^{3} \epsilon_{jkl}\sigma_{l} + 2\delta_{ij}I}{2}$$

$$= \delta_{ij}I + i\sum_{l=1}^{3} \epsilon_{jkl}\sigma_{l}$$

# Exercise 2.44

Suppose  $[A,B]=0,\{A,B\}=0$ , and A is invertible. Show that B must be 0.

### Solution

Concepts Involved: Linear Algebra, Commutators, Anticommutators.

By assumption, we have that:

$$[A, B] = AB - BA = 0$$
  
 $\{A, B\} = AB + BA = 0.$ 

Adding the first line to the second we have:

$$2AB = 0 \implies AB = 0.$$

 ${\cal A}^{-1}$  exists by the invertibility of  ${\cal A}$ , so multiplying by  ${\cal A}^{-1}$  on the left we have:

$$A^{-1}AB = A^{-1}0 \implies IB = 0 \implies B = 0.$$

# Exercise 2.45

Show that  $[A, B]^{\dagger} = [B^{\dagger}, A^{\dagger}].$ 

# Solution

Concepts Involved: Linear Algebra, Commutators, Adjoints.

Using the properties of the adjoint, we have:

$$[A, B]^{\dagger} = (AB - BA)^{\dagger} = (AB)^{\dagger} - (BA)^{\dagger} = B^{\dagger}A^{\dagger} - A^{\dagger}B^{\dagger} = [B^{\dagger}, A^{\dagger}]$$

# Exercise 2.46

Show that [A, B] = -[B, A].

### Solution

Concepts Involved: Linear Algebra, Commutators

By the definition of the commutator:

$$[A, B] = AB - BA = -(BA - AB) = -[B, A]$$

# Exercise 2.47

Suppose A and B are Hermitian. Show that i[A,B] is Hermitian.

### Solution

Concepts Involved: Linear Algebra, Commutators, Hermitian Operators

Suppose A, B are Hermitian. Using the results of Exercises 2.45 and 2.46, we have:

$$(i[A, B])^{\dagger} = -i([A, B])^{\dagger} = -i[B^{\dagger}, A^{\dagger}] = i[A^{\dagger}, B^{\dagger}] = i[A, B].$$

#### Exercise 2.48

What is the polar decomposition of a positive matrix P? Of a unitary matrix U? Of a Hermitian matrix, H?

#### Solution

**Concepts Involved:** Linear Algebra, Polar Decomposition, Positive Operators, Unitary Operators, Hermitian Operators

If P is a positive matrix, then no calculation is required; P=IP=PI is the polar decomposition (as I is unitary and P is positive). If U is a unitary matrix, then  $J=\sqrt{U^{\dagger}U}=\sqrt{I}=I$  and  $K=\sqrt{UU^{\dagger}}=\sqrt{I}=I$  so the polar decomposition is U=UI=IU (where U is unitary and I is positive). If H is hermitian, we then have that:

$$J = \sqrt{H^\dagger H} = \sqrt{H^2} = \sqrt{\sum_i \lambda_i^2 \left| i \middle\rangle\!\!\left\langle i \right|} = \sum_i \left| \lambda_i \right| \left| i \middle\rangle\!\!\left\langle i \right|$$

and  $K=\sqrt{HH^\dagger}=\sum_i|\lambda_i|\,|i\rangle\!\langle i|$  in the same way. Hence the polar decomposition is  $H=U\sum_i|\lambda_i|\,|i\rangle\!\langle i|=\sum_i|\lambda_i|\,|i\rangle\!\langle i|\,U.$ 

#### Exercise 2.49

Express the polar decomposition of a normal matrix in the outer product representation.

#### Solution

Concepts Involved: Linear Algebra, Polar Decomposition, Outer Products

Let A be a normal matrix. Then, A has spectral decomposition  $A = \sum_i \lambda_i \, |i\rangle\!\langle i|$ . Therefore, we have that:

$$A^{\dagger}A = AA^{\dagger} = \sum_{i} \sum_{i'} \lambda_{i} \lambda_{i'}^{*} |i\rangle\langle i| |i'\rangle\langle i'| = \sum_{i} \sum_{i'} \lambda_{i} \lambda_{i'}^{*} |i\rangle\langle i'| \delta_{ii'} = \sum_{i} |\lambda_{i}|^{2} |i\rangle\langle i|$$

We then have that:

$$J = \sqrt{A^{\dagger}A} = \sqrt{\sum_{i} |\lambda_{i}|^{2} |i\rangle\langle i|} = \sum_{i} |\lambda_{i}| |i\rangle\langle i|$$

and  $K = \sum_i |\lambda_i| |i\rangle\langle i|$  identically. Furthermore, U is unitary, so it also has a spectral decomposition of  $\sum_i \mu_j |j\rangle\langle j|$ . Hence we have the polar decomposition in the outer product representation as:

$$\begin{split} A &= UJ = KU \\ A &= U \sum_{i} |\lambda_{i}| |i\rangle\langle i| \sum_{j} = \sum_{i} |\lambda_{i}| |i\rangle\langle i| \sum_{j} U \\ A &= \sum_{j} \sum_{i} \mu_{j} |\lambda_{i}| |j\rangle\langle j| i\rangle\langle i| = \sum_{i} \sum_{j} |\lambda_{i}| \mu_{j} |i\rangle\langle i| j\rangle\langle j| \end{split}$$

## 

## Exercise 2.50

Find the left and right polar decompositions of the matrix

$$\begin{bmatrix} 1 & 0 \\ 1 & 1 \end{bmatrix}$$

#### Solution

Concepts Involved: Linear Algebra, Polar Decomposition.

Let  $A=\begin{bmatrix}1&0\\1&1\end{bmatrix}$  . We start with the left polar decomposition, and hence find  $J=\sqrt{A^\dagger A}$ . In order to do

this, we find the spectral decompositions of  $A^{\dagger}A$  and  $AA^{\dagger}$ .

$$\det \begin{pmatrix} A^{\dagger}A - I\lambda \end{pmatrix} = 0 \implies \det \begin{pmatrix} \begin{bmatrix} 1 & 1 \\ 0 & 1 \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 1 & 1 \end{bmatrix} - \begin{bmatrix} \lambda & 0 \\ 0 & \lambda \end{bmatrix} \end{pmatrix} = 0 \implies \det \begin{bmatrix} 2 - \lambda & 1 \\ 1 & 1 - \lambda \end{bmatrix} = 0$$

$$\implies \lambda^2 - 3\lambda + 1 = 0$$

$$\implies \lambda_1 = \frac{3 + \sqrt{5}}{2}, \lambda_2 = \frac{3 - \sqrt{5}}{2}$$

Solving for the eigenvectors, we have:

$$\begin{bmatrix} 2 - \frac{3+\sqrt{5}}{2} & 1\\ 1 & 1 - \frac{3+\sqrt{5}}{2} \end{bmatrix} |v_1\rangle = \mathbf{0} \implies |v_1\rangle = \begin{bmatrix} 1+\sqrt{5}\\ 2 \end{bmatrix}$$
$$\begin{bmatrix} 2 - \frac{3-\sqrt{5}}{2} & 1\\ 1 & 1 - \frac{3-\sqrt{5}}{2} \end{bmatrix} |v_2\rangle = \mathbf{0} \implies |v_2\rangle = \begin{bmatrix} 1-\sqrt{5}\\ 2 \end{bmatrix}$$

Normalizing, we get:

$$|v_1\rangle = \frac{1}{\sqrt{10 + 2\sqrt{5}}} \begin{bmatrix} 1 + \sqrt{5} \\ 2 \end{bmatrix}, \quad |v_2\rangle = \frac{1}{\sqrt{10 - 2\sqrt{5}}} \begin{bmatrix} 1 - \sqrt{5} \\ 2 \end{bmatrix}$$

The spectral decomposition of  $A^{\dagger}A$  is therefore:

$$A^{\dagger}A = \lambda_1 |v_1\rangle\langle v_1| + \lambda_2 |v_2\rangle\langle v_2|$$

Calculating J, we therefore have:

$$J = \sqrt{A^{\dagger}A} = \sqrt{\lambda_1} |v_1\rangle\langle v_1| + \sqrt{\lambda_2} |v_2\rangle\langle v_2| = \frac{1}{\sqrt{5}} \begin{bmatrix} 3 & 1\\ 1 & 2 \end{bmatrix}$$

The last equality is not completely trivial, but the algebra is tedious so we invite the reader to use a symbolic calculator, as we have. We make the observation that:

$$A = UJ \implies U = AJ^{-1}$$

So calculating  $J^{-1}$ , we have:

$$J^{-1} = \frac{1}{\sqrt{\lambda_1}} |v_1\rangle\langle v_1| + \frac{1}{\sqrt{\lambda_2}} |v_2\rangle\langle v_2| = \frac{1}{\sqrt{5}} \begin{bmatrix} 2 & -1 \\ -1 & 3 \end{bmatrix}$$

Where we again have used the help of a symbolic calculator. Calculating U, we then have that:

$$U = AJ^{-1} = \begin{bmatrix} 1 & 0 \\ 1 & 1 \end{bmatrix} \frac{1}{\sqrt{5}} \begin{bmatrix} 2 & -1 \\ -1 & 3 \end{bmatrix} = \frac{1}{\sqrt{5}} \begin{bmatrix} 2 & -1 \\ 1 & 2 \end{bmatrix}$$

Hence the left polar decomposition of A is given by:

$$A = UJ = \begin{pmatrix} \frac{1}{\sqrt{5}} \begin{bmatrix} 2 & -1\\ 1 & 2 \end{bmatrix} \end{pmatrix} \begin{pmatrix} \frac{1}{\sqrt{5}} \begin{bmatrix} 3 & 1\\ 1 & 2 \end{bmatrix} \end{pmatrix}$$

We next solve for the right polar decomposition. We could repeat the procedure of solving for the spectral decomposition of  $AA^{\dagger}$ , but we take a shortcut; since we know the K that satisfies:

$$A = KU$$

will be unique, and U is unitary, we can simply multiply both sides of the above equation on the right by  $U^{-1}=U^\dagger$  to obtain K. Hence:

$$K = AU^\dagger = \begin{bmatrix} 1 & 0 \\ 1 & 1 \end{bmatrix} \frac{1}{\sqrt{5}} \begin{bmatrix} 2 & 1 \\ -1 & 2 \end{bmatrix} = \frac{1}{\sqrt{5}} \begin{bmatrix} 2 & 1 \\ 1 & 3 \end{bmatrix}.$$

Therefore the right polar decomposition of A is given by:

$$A = KU = \left(\frac{1}{\sqrt{5}} \begin{bmatrix} 2 & 1\\ 1 & 3 \end{bmatrix}\right) \left(\frac{1}{\sqrt{5}} \begin{bmatrix} 2 & -1\\ 1 & 2 \end{bmatrix}\right)$$

Exercise 2.51

Verify that the Hadamard gate H is unitary.

Solution

Concepts Involved: Linear Algebra, Unitary Operators

We observe that:

$$H^{\dagger}H = \frac{1}{\sqrt{2}} \begin{bmatrix} 1 & 1\\ 1 & -1 \end{bmatrix} \frac{1}{\sqrt{2}} \begin{bmatrix} 1 & 1\\ 1 & -1 \end{bmatrix} = \begin{bmatrix} 1 & 0\\ 0 & 1 \end{bmatrix}$$

showing that H is indeed unitary.

**Remark:** The above calculation shows that H is also Hermitian and Idempotent.

Exercise 2.52

Verify that  $H^2 = I$ .

#### Solution

Concepts Involved: Linear Algebra

See the calculation and remark in the previous exercise.

## Exercise 2.53

What are the eigenvalues and eigenvectors of H?

## Solution

Concepts Involved: Linear Algebra, Eigenvalues, Eigenvectors

Using the characteristic equation to find the eigenvalues, we have:

$$\det(H - I\lambda) = 0 \implies \det\begin{bmatrix} \frac{1}{\sqrt{2}} - \lambda & \frac{1}{\sqrt{2}} \\ \frac{1}{\sqrt{2}} & -\frac{1}{\sqrt{2}} - \lambda \end{bmatrix} = 0 \implies \lambda^2 - 1 = 0$$
$$\implies \lambda_1 = 1, \lambda_2 = -1$$

Finding the eigenvectors, we then have:

$$\begin{bmatrix} \frac{1}{\sqrt{2}} - 1 & \frac{1}{\sqrt{2}} \\ \frac{1}{\sqrt{2}} & -\frac{1}{\sqrt{2}} - 1 \end{bmatrix} |v_1\rangle = \mathbf{0} \implies |v_1\rangle = \begin{bmatrix} 1 + \frac{1}{\sqrt{2}} \\ \frac{1}{\sqrt{2}} \end{bmatrix}$$
$$\begin{bmatrix} \frac{1}{\sqrt{2}} + 1 & \frac{1}{\sqrt{2}} \\ \frac{1}{\sqrt{2}} & -\frac{1}{\sqrt{2}} + 1 \end{bmatrix} |v_2\rangle = \mathbf{0} \implies |v_2\rangle = \begin{bmatrix} -1 + \frac{1}{\sqrt{2}} \\ \frac{1}{\sqrt{2}} \end{bmatrix}$$

Normalizing, we have:

$$|v_1\rangle = \frac{1}{\sqrt{2+\sqrt{2}}} \begin{bmatrix} 1 + \frac{1}{\sqrt{2}} \\ \frac{1}{\sqrt{2}} \end{bmatrix}, |v_2\rangle = \frac{1}{\sqrt{2-\sqrt{2}}} \begin{bmatrix} -1 + \frac{1}{\sqrt{2}} \\ \frac{1}{\sqrt{2}} \end{bmatrix}$$

## Exercise 2.54

Suppose A and B are commuting Hermitian operators. Prove that  $\exp(A)\exp(B)=\exp(A+B)$ . (Hint: Use the results of Section 2.1.9.)

#### Solution

Concepts Involved: Linear Algebra, Operator Functions, Simultaneous Diagonalization

Since A,B commute, they can be simulatneously diagonalized; that is, there exists some orthonormal basis  $|i\rangle$  of V such that  $A=\sum_i a_i\,|i\rangle\langle i|$  and  $B=\sum_i b_i\,|i\rangle\langle i|$ . Hence, using the definition of operator

functions, we have that:

$$\exp(A) \exp(B) = \exp\left(\sum_{i} a_{i} |i\rangle\langle i|\right) \exp\left(\sum_{i'} b_{i'} |i'\rangle\langle i'|\right)$$

$$= \sum_{i} \sum_{i'} \exp(a_{i}) \exp(b_{i'}) |i\rangle\langle i|i'\rangle\langle i'|$$

$$= \sum_{i} \sum_{i'} \exp(a_{i}) \exp(b_{i'}) |i\rangle\langle i'|\delta_{ii'}$$

$$= \sum_{i} \exp(a_{i}) \exp(b_{i}) |i\rangle\langle i|$$

$$= \sum_{i} \exp(a_{i} + b_{i}) |i\rangle\langle i|$$

$$= \exp\left(\sum_{i} (a_{i} + b_{i}) |i\rangle\langle i|\right)$$

$$= \exp(A + B)$$

Exercise 2.55

Prove that  $U(t_1, t_2)$  defined in Equation (2.91) is unitary.

#### Solution

Concepts Involved: Linear Algebra, Unitary Operators, Spectral Decomposition, Operator Functions.

Since the Hamiltonian H is Hermitian, it is normal and hence has spectral decomposition:

$$H = \sum_E E \, |E\rangle\!\langle E|$$

where all E are real by the Hermicity of H, and  $|E\rangle$  is an orthonormal basis of the Hilbert space. We then have that:

$$\begin{split} U(t_1, t_2) &\equiv \exp\left[\frac{-iH(t_2 - t_1)}{\hbar}\right] = \exp\left[\frac{-i\sum_E E \left|E\right\rangle\!\!\left\langle E\right|\left(t_2 - t_1\right)}{\hbar}\right] \\ &= \sum_E \exp\!\left(\frac{-iE(t_2 - t_1)}{\hbar}\right) \left|E\right\rangle\!\!\left\langle E\right| \end{split}$$

Hence calculating  $U^{\dagger}(t_1, t_2)$  we have:

$$U^{\dagger}(t_1, t_2) = \left(\sum_{E} \exp\left(\frac{-iE(t_2 - t_1)}{\hbar}\right) |E\rangle\langle E|\right)^{\dagger} = \sum_{E} \left(\exp\left(\frac{-iE(t_2 - t_1)}{\hbar}\right)\right)^{*} \left(|E\rangle\langle E|\right)^{\dagger}$$
$$= \sum_{E} \exp\left(\frac{iE(t_2 - t_1)}{\hbar}\right) |E\rangle\langle E|$$

Therefore computing  $U^{\dagger}(t_1, t_2)U(t_2, t_1)$  we have:

$$U^{\dagger}(t_{2}, t_{1})U(t_{2}, t_{1}) = \left(\sum_{E} \exp\left(\frac{-iE(t_{2} - t_{1})}{\hbar}\right) |E\rangle\langle E|\right) \left(\sum_{E'} \exp\left(\frac{iE'(t_{2} - t_{1})}{\hbar}\right) |E'\rangle\langle E'|\right)$$

$$= \sum_{E} \sum_{E'} \exp\left(\frac{-iE(t_{2} - t_{1})}{\hbar}\right) \exp\left(\frac{iE'(t_{2} - t_{1})}{\hbar}\right) \delta_{EE'} |E\rangle\langle E'|$$

$$= \sum_{E} \exp\left(\frac{-iE(t_{2} - t_{1})}{\hbar}\right) \exp\left(\frac{iE(t_{2} - t_{1})}{\hbar}\right) |E\rangle\langle E|$$

$$= \sum_{E} |E\rangle\langle E|$$

$$= I$$

where in the second equality we use the fact that the eigenstates are orthogonal. We conclude that U is unitary.  $\Box$ 

## Exercise 2.56

Use the spectral decomposition to show that  $K \equiv -i \log(U)$  is Hermitian for any unitary U, and thus  $U = \exp(iK)$  for some Hermitian K.

## Solution

**Concepts Involved:** Linear Algebra, Hermitian Operators, Unitary Operators, Spectral Decomposition, Operator Functions.

Suppose U is unitary. Then, U is normal and hence has spectral decomposition:

$$U = \sum_{j} \lambda_{j} |j\rangle\langle j|$$

where  $|j\rangle$  are eigenvectors of U with eigenvalues  $\lambda_j$ , and  $|j\rangle$  forms an orthonormal basis of the Hilbert space. By Exercise 2.18, all eigenvalues of unitary operators have eigenvalues of modulus 1, so we can let  $\lambda_j = \exp(i\theta_j)$  where  $\theta_j \in \mathbb{R}$  and hence write the above as:

$$U = \sum_{j} \exp(i\theta_j) |j\rangle\langle j|$$

We then have that:

$$K \equiv -i\log(U) = -i\log\left(\sum_{j} \exp(i\theta_{j})|j\rangle\langle j|\right) = \sum_{j} -i\log\left(\exp(i\theta_{j})\right)|j\rangle\langle j| = \sum_{j} -i(i\theta_{j})|j\rangle\langle j|$$
$$= \sum_{j} \theta_{j}|j\rangle\langle j|$$

We then observe that:

$$K^{\dagger} = \left(\sum_{j} \theta_{j} |j\rangle\!\langle j|\right)^{\dagger} = \sum_{j} \theta_{j} |j\rangle\!\langle j|$$

as the  $\theta_j$ s are real and  $(|j\rangle\langle j|)^\dagger=|j\rangle\langle j|$ . Hence K is Hermitian. Then, multiplying both sides in  $K=-i\log(U)$  by i and exponentiating both sides, we obtain the desired relation.

## Exercise 2.57: Cascaded measurements are single measurements

Suppose  $\{L_l\}$  and  $\{M_m\}$  are two sets of measurement operators. Show that a measurement defined by the measurement operators  $\{L_l\}$  followed by a measurement defined by the measurement operators  $\{M_m\}$  is physically equivalent to a single measurement defined by measurement operators  $\{N_{lm}\}$  with the representation  $N_{lm} \equiv M_m L_l$ .

#### Solution

Concepts Involved: Linear Algebra, Quantum Measurement.

Suppose we have (normalized) initial quantum state  $|\psi_0\rangle$ . Then, the state after measurement of  $L_l$  is given by definition to be:

$$|\psi_0\rangle \mapsto |\psi_1\rangle = \frac{L_l|\psi_0\rangle}{\sqrt{\langle\psi_0|L_l^{\dagger}L_l|\psi_0\rangle}}.$$

The state after measurement of  $M_m$  on  $|\psi_1\rangle$  is then given to be:

$$\begin{split} |\psi_1\rangle \mapsto |\psi_2\rangle &= \frac{M_m |\psi_1\rangle}{\sqrt{\langle\psi_1|M_m^\dagger M_m|\psi_1\rangle}} = \frac{M_m \left(\frac{L_l|\psi_0\rangle}{\sqrt{\langle\psi_0|L_l^\dagger L_l|\psi_0\rangle}}\right)}{\sqrt{\left(\frac{L_l^\dagger \langle\psi_0|}{\sqrt{\langle\psi_0|L_l^\dagger L_l|\psi_0\rangle}}\right)M_m^\dagger M_m \left(\frac{L_l|\psi_0\rangle}{\sqrt{\langle\psi_0|L_l^\dagger L_l|\psi_0\rangle}}\right)}} \\ &= \frac{M_m L_l|\psi_0\rangle}{\sqrt{\langle\psi_0|L_l^\dagger L_l|\psi_0\rangle}} \frac{\sqrt{\langle\psi_0|L_l^\dagger L_l|\psi_0\rangle}}{\sqrt{\langle\psi_0|L_l^\dagger M_m^\dagger M_m L_l^\dagger|\psi_0\rangle}} \\ &= \frac{M_m L_l|\psi_0\rangle}{\sqrt{\langle\psi_0|L_l^\dagger M_m^\dagger M_m L_l^\dagger|\psi_0\rangle}}. \end{split}$$

Conversely, the state of  $|\psi_0\rangle$  after measurement of  $N_{lm}=M_mL_l$  is given by:

$$|\psi_0\rangle \mapsto |\psi_3\rangle = \frac{M_m L_l |\psi_0\rangle}{\sqrt{\langle \psi_0 | L_l^{\dagger} M_m^{\dagger} M_m L_l |\psi_0\rangle}}.$$

We see that  $|\psi_2\rangle = |\psi_3\rangle$  (that is, the cascaded measurement produces the same result as the single measurement), proving the claim.

#### Exercise 2.58

Suppose we prepare a quantum system in an eigenstate  $|\psi\rangle$  of some observable M with corresponding eigenvalue m. What is the average observed value of M, and the standard deviation?

### Solution

Concepts Involved: Linear Algebra, Quantum Measurement, Expectation, Standard Deviation.

By the definition of expectation, we have that:

$$\langle M \rangle_{|\psi\rangle} = \langle \psi | M | \psi \rangle = \langle \psi | m | \psi \rangle = m \langle \psi | \psi \rangle = m$$

Where in the second equality we use that  $|\psi\rangle$  is an eigenstate of M with eigenvalue m, and in the last equality we use that  $|\psi\rangle$  is a normalized quantum state. Next, calculating  $\langle M^2\rangle_{|\psi\rangle}$ , we have:

$$\left\langle M^2\right\rangle_{|\psi\rangle} = \langle \psi|M^2|\psi\rangle = \langle \psi|MM|\psi\rangle = \langle \psi|M^\dagger M|\psi\rangle = \langle \psi|m^*m|\psi\rangle = \langle \psi|m^2|\psi\rangle = m^2\langle \psi|\psi\rangle = m^2.$$

Note that we have used the fact that M is Hermitian (it is an observable) to use that  $M^{\dagger}=M$  and  $m^*=m$  as all eigenvalues of Hermitian operators are real. Now calculating the standard deviation, we have:

$$\Delta(M) = \sqrt{\langle M^2 \rangle - \langle M \rangle^2} = \sqrt{m^2 - (m)^2} = 0$$

## Exercise 2.59

Suppose we have qubit in the state  $|0\rangle$ , and we measure the observable X. What is the average value of X? What is the standard deviation of X?

## Solution

**Concepts Involved:** Linear Algebra, Quantum Measurement, Projective Measurement, Expectation, Standard Deviation.

By the definition of expectation, we have:

$$\langle X \rangle_{|0\rangle} = \langle 0|X|0\rangle = \langle 0|1\rangle = 0$$

Next calculating  $\langle X^2 \rangle_{|0\rangle}$ , we have:

$$\left\langle X^2 \right\rangle_{|0\rangle} = \left\langle 0|XX|0\right\rangle = \left\langle 1|1\right\rangle = 1$$

Hence the standard deviation of X is given by:

$$\Delta(X) = \sqrt{\langle X^2 \rangle - \langle X \rangle^2} = \sqrt{1 - 0} = 1$$

# Exercise 2.60

Show that  ${\bf v}\cdot {\boldsymbol \sigma}$  has eigenvalues  $\pm 1$ , and that the projectors into the corresponding eigenspaces are given by  $P_{\pm}=(I\pm {\bf v}\cdot {\boldsymbol \sigma})/2.$ 

#### Solution

Concepts Involved: Eigenvalues, Projectors.

Let  $|v\rangle$  be a unit vector. We already showed in Exercise 2.35 that  $\mathbf{v}\cdot\boldsymbol{\sigma}$  has eigenvalues  $\lambda_+=1,\lambda_-=-1$ . We next prove a general statement; namely, that for a observable on a 2-dimensional Hilbert space with eigenvalues  $\lambda_\pm=\pm 1$  has projectors

$$P_{\pm} = \frac{I \pm O}{2}$$

To see this is the case, let  $P_+=|o_+\rangle\langle o_+|$ ,  $P_-=|o_-\rangle\langle o_-|$ ,  $I=|o_+\rangle\langle o_+|+|o_-\rangle\langle o_-|$ , and  $O=|o_+\rangle\langle o_+|-|o_-\rangle\langle o_-|$ . We then have that:

$$\begin{split} \frac{I+O}{2} &= \frac{|o_+\rangle\!\langle o_+| - |o_-\rangle\!\langle o_-|}{2} = |o_+\rangle\!\langle o_+| = P_+ \\ \frac{I-O}{2} &= \frac{|o_-\rangle\!\langle o_-| + |o_-\rangle\!\langle o_-| - |o_+\rangle\!\langle o_+| + |o_-\rangle\!\langle o_-|}{2} = |o_-\rangle\!\langle o_-| = P_- \rangle \end{split}$$

Hence the general statement is proven. Applying this to  $O = \mathbf{v} \cdot \boldsymbol{\sigma}$  (which is indeed Hermitian and hence an observable as each of X, Y, Z are Hermitian), we get that:

$$P_{\pm} = \frac{I \pm \mathbf{v} \cdot \boldsymbol{\sigma}}{2}$$

as claimed.

## Exercise 2.61

Calculate the probability of obtaining the result +1 for a measurement of  $\mathbf{v} \cdot \boldsymbol{\sigma}$ , given that the state prior to measurement is  $|0\rangle$ . What is the state of the system after measurement if +1 is obtained?

#### Solution

Concepts Involved: Quantum Measurement, Projective Measurement.

The probability of obtaining the result +1 is given by:

$$p(+) = \langle 0|P_+|0\rangle = \langle 0|\frac{I + \mathbf{v} \cdot \boldsymbol{\sigma}}{2}|0\rangle$$

We recall from from Exercise 2.35 that:

$$\mathbf{v} \cdot \boldsymbol{\sigma} = \begin{bmatrix} v_3 & v_1 - iv_2 \\ v_1 + iv_2 & -v_3 \end{bmatrix} = v_3 |0\rangle\langle 0| + (v_1 - iv_2) |0\rangle\langle 1| + (v_1 + iv_2) |1\rangle\langle 0| - v_3 |1\rangle\langle 1|.$$

Hence computing p(+), we get:

$$\begin{split} p(+) &= \langle 0| \left(\frac{1}{2}|0\rangle + \frac{1}{2} \left(v_3|0\rangle + (v_1 + iv_2)|1\rangle\right)\right) \\ &= \langle 0| \left(\frac{1+v_3}{2}|0\rangle + \frac{v_1+iv_2}{2}|1\rangle\right) \\ &= \frac{1+v_3}{2} \langle 0|0\rangle + \frac{v_1+iv_2}{2} \langle 0|1\rangle = \frac{1+v_3}{2} \end{split}$$

so the probability of measuring the +1 outcome is  $\frac{1+v_3}{2}$ . The state after the measurement of the +1 outcome is given by:

$$|0\rangle \mapsto \frac{P_{+}|0\rangle}{\sqrt{p(+)}} = \frac{\frac{1+v_{3}}{2}|0\rangle + \frac{v_{1}+iv_{2}}{2}|1\rangle}{\sqrt{\frac{1+v_{3}}{2}}} = \frac{1}{\sqrt{2(1+v_{3})}} \left( (1+v_{3})|0\rangle + (v_{1}+iv_{2})|1\rangle \right)$$

## Exercise 2.62

Show that any measurement where the measurement operators and the POVM elements coincide is a projective measurement.

#### Solution

Concepts Involved: Quantum Measurement, Projective Measurement, POVM Measurement.

Suppose we have that the measurement operators  $M_m$  are equal to the POVM elements  $E_m$ . In this case, we have that:

$$M_m = E_m \equiv M_m^{\dagger} M_m$$

 $M_m^{\dagger}M_m$  is positive by Exercise 2.25, so it follows that  $M_m$  is positive and hence Hermitian by Exercise 2.24. Hence,  $M_m^{\dagger}=M_m$ , and therefore:

$$M_m = M_m^{\dagger} M_m = M_m^2$$

From which we conclude that  $M_m$  are projective measurement operators.

## Exercise 2.63

Suppose a measurement is described by measurement operators  $M_m$ . Show that there exist unitary operators  $U_m$  such that  $M_m = U_m \sqrt{E_m}$ , where  $E_m$  is the POVM associated to the measurement.

#### Solution

Concepts Involved: Quantum Measurement, POVM Measurement, Polar Decomposition.

Since  $M_m$  is a linear operator, by the left polar decomposition there exists unitary U such that:

$$M_m = U\sqrt{M_m^{\dagger} M_m} = U\sqrt{E_m},$$

where in the last equality we use that  $M_m^{\dagger}M_m=E_m.$ 

#### Exercise 2.64

(\*) Suppose Bob is given a quantum state chosen from a set  $|\psi_1\rangle,\ldots,|\psi_m\rangle$  of linearly independent states. Construct a POVM  $\{E_1,E_2,\ldots,E_{m+1}\}$  such that if outcome  $E_i$  occurs,  $1\leq i\leq m$ , then Bob knows with certainty that he was given the state  $|\psi_i\rangle$ . (The POVM must be such that  $\langle\psi_i|E_i|\psi_i\rangle>0$  for each i.)

#### Solution

Concepts Involved: POVM Measurement, Orthogonality

Let  $\mathcal H$  be the Hilbert space where the given states lie, and let V be the m-dimensional subspace spanned by  $|\psi_1\rangle,\ldots,|\psi_m\rangle$ . For each  $i\in\{1,\ldots,m\}$ , let  $W_i$  be the subspace of V spanned by  $\{|\psi_j\rangle:j\neq i\}$ . Let  $W_i^\perp$  be the orthogonal complement of  $W_i$  which consists of all states in  $\mathcal H$  orthogonal to all states in  $W_i$ . We then have that any vector in V can be written as the sum of a vector in  $W_i$  and  $W_i^\perp\cap V$  (see for example Theorem 6.47 in Axler's Linear Algebra Done Right). Therefore, for any  $|\psi_i\rangle$  we can write:

$$|\psi_i\rangle = |w_i\rangle + |p_i\rangle$$

Where  $|w_i\rangle \in W_i$  and  $|p_i\rangle \in W_i^{\perp} \cap V$ . Define  $E_i = \frac{|p_i\rangle\langle p_i|}{m}$ . By construction, we have that for any  $|\psi\rangle \in \mathcal{H}$ :

$$\langle \psi | E_i | \psi \rangle = \frac{\left| \langle \psi | p_i \rangle \right|^2}{m} \ge 0$$

so the  $E_i$ s are positive are required. Furthermore, defining  $E_{i+1} = I - \sum_{i=1}^m E_i$  we again see that for any  $|\psi\rangle \in \mathcal{H}$ :

$$\langle \psi | E_{i+1} | \psi \rangle = \langle \psi | I | \psi \rangle - \sum_{i=1}^{m} \langle \psi | E_i | \psi \rangle = 1 - \sum_{i=1}^{m} \langle \psi | E_i | \psi \rangle \ge 1 - \sum_{i=1}^{m} \frac{1}{m} = 0$$

so  $E_{i+1}$  is also positive as required. Finally, to see that the  $E_1, \ldots E_m$  have the desired properties, observe by construction that since  $|p_i\rangle \in W_i^\perp \cap V$ , it follows that  $\langle \psi_j|p_i\rangle = 0$  for any  $j\neq i$  (as the  $|p_i\rangle$  will be orthogonal to all the vectors in  $\{|\psi_j\rangle:j\neq i\}$  by construction). Calculating  $\langle \psi_i|E_i|\psi_i\rangle$ , we observe that:

$$\langle \psi_i | E_i | \psi_i \rangle = \left( \langle w_i | + \langle p_i | \right) \frac{|p_i \rangle \langle p_i|}{m} \left( |w_i \rangle + |p_i \rangle \right) = \frac{\left| \langle p_i | p_i \rangle \right|^2}{m} = \frac{1}{m} \ge 0$$

so if Bob measures  $E_i$ , he can be certain that he was given the state  $|\psi_i\rangle$ .

## Exercise 2.65

Express the states  $(|0\rangle + |1\rangle)/\sqrt{2}$  and  $(|0\rangle - |1\rangle)/\sqrt{2}$  is a basis in which they are *not* the same up to a relative phase shift.

#### Solution

Concepts Involved: Linear Algebra, Phase

Let us define our basis to be  $|+\rangle:=(|0\rangle+|1\rangle)/\sqrt{2}$  and  $|-\rangle:=(|0\rangle-|1\rangle)/\sqrt{2}$ . Our two states are then just the basis vectors of this basis  $(|+\rangle,|-\rangle)$  and are not the same up to relative phase shift.

## Exercise 2.66

Show that the average value of the observable  $X_1Z_2$  for a two qubit system measured in the state  $(|00\rangle+|11\rangle)/\sqrt{2}$  is zero.

#### Solution

Concepts Involved: Quantum Measurement, Expectation, Composite Systems

Computing the expectation value of  $X_1Z_2$ , we get:

$$\langle X_1 Z_2 \rangle = \left( \frac{\langle 00| + \langle 11| \rangle}{\sqrt{2}} \right) X_1 Z_2 \left( \frac{|00\rangle + |11\rangle}{\sqrt{2}} \right)$$

$$= \left( \frac{\langle 00| + \langle 11| \rangle}{\sqrt{2}} \right) \left( \frac{X_1 Z_2 |00\rangle + X_1 Z_2 |11\rangle}{\sqrt{2}} \right)$$

$$= \left( \frac{\langle 00| + \langle 11| \rangle}{\sqrt{2}} \right) \left( \frac{|10\rangle - |01\rangle}{\sqrt{2}} \right)$$

$$= \frac{1}{2} \left( \langle 00|10\rangle - \langle 00|01\rangle + \langle 11|10\rangle - \langle 11|01\rangle \right)$$

$$= \frac{1}{2} \left( 0 + 0 + 0 + 0 \right)$$

$$= 0.$$

#### Exercise 2.67

Suppose V is a Hilbert space with a subspace W. Suppose  $U:W\mapsto V$  is a linear operator which preserves inner products, that is, for any  $|w_1\rangle$  and  $|w_2\rangle$  in W,

$$\langle w_1|U^{\dagger}U|w_2\rangle = \langle w_1|w_2\rangle$$

Prove that there exists a unitary operator  $U':V\mapsto V$  which extends U. That is,  $U'|w\rangle=U|w\rangle$  for all  $|w\rangle$  in W, but U' is defined on the entire space V. Usually we omit the prime symbol ' and just write U to denote the extension.

### Solution

Concepts Involved: Linear Algebra, Inner Products, Unitary Operators

By assumption we have that U is unitary on W as  $\langle w_1|U^{\dagger}U|w_2\rangle=\langle w_1|w_2\rangle$  and hence  $U^{\dagger}U=I_W$ . Hence, it has spectral decomposition:

$$U = \sum_{j} \lambda_{j} |j\rangle\langle j|$$

where  $\{|j\rangle\}$  is an orthonormal basis of the subspace W. Then, let  $\{|j\rangle\} \cup \{|i\rangle\}$  be an orthonormal basis of the full space V. We then define:

$$U' = \sum_{j} \lambda_{j} |j\rangle\langle j| + \sum_{i} |i\rangle\langle i| = U + \sum_{i} |i\rangle\langle i|$$

We can then see that for any  $|w\rangle \in W$  that:

$$U'|w\rangle = \left(U + \sum_{i} |i\rangle\langle i|\right)|w\rangle = U|w\rangle + \sum_{j} |i\rangle\langle i|w\rangle = U|w\rangle$$

where in the last line we use that  $\langle i|w\rangle=0$  as  $|i\rangle$  are not in the subspace W. Finally, verifying the unitarity of U' we have that:

$$U'^{\dagger}U' = \left(\sum_{j} \lambda_{j}^{*} |j\rangle\langle j| + \sum_{i} |i\rangle\langle i|\right) \left(\sum_{j'} \lambda_{j} |j\rangle\langle j| + \sum_{i'} |i\rangle\langle i|\right)$$

$$= \sum_{j} \sum_{j'} |j\rangle\langle j|j'\rangle\langle j'| + \sum_{j} \sum_{i'} |j\rangle\langle j|i'\rangle\langle i'| + \sum_{i} \sum_{j'} |i\rangle\langle i|j'\rangle\langle j'| + \sum_{i} \sum_{i'} |i\rangle\langle i|i'\rangle\langle i'|$$

$$= \sum_{j} \sum_{j'} \langle j|j'\rangle\delta_{jj'} + \sum_{i} \sum_{i'} \langle i|i'\rangle\delta_{ii'}$$

$$= \sum_{j} |j\rangle\langle j| + \sum_{i} |i\rangle\langle i|$$

$$= I$$

#### Exercise 2.68

Prove that  $|\psi\rangle \neq |a\rangle|b\rangle$  for all single qubit states  $|a\rangle$  and  $|b\rangle$ .

## Solution

Concepts Involved: Linear Algebra, Composite Systems, Entanglement.

Recall that:

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

Suppose for the take of contradiction that  $|\psi\rangle=|a\rangle|b\rangle$  for some single qubit states  $|a\rangle$  and  $|b\rangle$ . Then, we have that  $|a\rangle=\alpha|0\rangle+\beta|1\rangle$  and  $|b\rangle=\gamma|0\rangle+\delta|1\rangle$  for some  $\alpha,\beta,\gamma,\delta\in\mathbb{C}$  such that  $|\alpha|^2+|\beta|^2=1$  and  $|\gamma|^2+|\delta|^2=1$ . We then have that:

$$|a\rangle|b\rangle = (\alpha|0\rangle + \beta|1\rangle)(\gamma|0\rangle + \delta|1\rangle) = \alpha\gamma|00\rangle + \alpha\delta|01\rangle + \beta\gamma|10\rangle + \beta\delta|11\rangle.$$

Where we have used the linearity of the tensor product (though we supress the  $\otimes$  symbols in the above expression). We then have that:

$$|\psi\rangle = \frac{|00\rangle + |11\rangle}{\sqrt{2}} = \alpha\gamma|00\rangle + \alpha\delta|01\rangle + \beta\gamma|10\rangle + \beta\delta|11\rangle$$

which forces  $\alpha\delta=0$  and  $\beta\gamma=0$ . However, we then have that at least one of  $\alpha\gamma$  or  $\beta\delta$  is also zero, and we thus reach a contradiction.

## Exercise 2.69

Verify that the Bell basis forms an orthonormal basis for the two qubit state space.

## Solution

Concepts Involved: Linear Algebra, Orthogonality, Bases, Composite Systems.

Recall that the bell basis is given by:

$$|B_{00}\rangle = \frac{|00\rangle + |11\rangle}{\sqrt{2}}, \quad |B_{01}\rangle \frac{|00\rangle - |11\rangle}{\sqrt{2}}, \quad |B_{10}\rangle = \frac{|01\rangle + |10\rangle}{\sqrt{2}}, \quad |B_{11}\rangle = \frac{|01\rangle - |10\rangle}{\sqrt{2}}$$

We first verify orthonormality. We observe that:

$$\langle B_{00} | B_{00} \rangle = \left( \frac{\langle 00 | + \langle 11 |}{\sqrt{2}} \right) \left( \frac{|00 \rangle + |11 \rangle}{\sqrt{2}} \right) = \frac{1}{2} \left( \langle 00 | 00 \rangle + \langle 00 | 11 \rangle + \langle 11 | 00 \rangle + \langle 11 | 11 \rangle \right) = 1$$

$$\langle B_{00} | B_{01} \rangle = \left( \frac{\langle 00 | + \langle 11 |}{\sqrt{2}} \right) \left( \frac{|00 \rangle - |11 \rangle}{\sqrt{2}} \right) = \frac{1}{2} \left( \langle 00 | 00 \rangle - \langle 00 | 11 \rangle + \langle 11 | 00 \rangle - \langle 11 | 11 \rangle \right) = 0$$

$$\langle B_{00} | B_{10} \rangle = \left( \frac{\langle 00 | + \langle 11 |}{\sqrt{2}} \right) \left( \frac{|01 \rangle + |10 \rangle}{\sqrt{2}} \right) = \frac{1}{2} \left( \langle 00 | 01 \rangle + \langle 00 | 10 \rangle + \langle 11 | 01 \rangle + \langle 11 | 10 \rangle \right) = 0$$

$$\langle B_{00} | B_{11} \rangle = \left( \frac{\langle 00 | - \langle 11 |}{\sqrt{2}} \right) \left( \frac{|01 \rangle - |10 \rangle}{\sqrt{2}} \right) = \frac{1}{2} \left( \langle 00 | 01 \rangle - \langle 00 | 10 \rangle + \langle 11 | 01 \rangle - \langle 11 | 10 \rangle \right) = 0$$

$$\langle B_{01} | B_{01} \rangle = \left( \frac{\langle 00 | - \langle 11 |}{\sqrt{2}} \right) \left( \frac{|01 \rangle + |10 \rangle}{\sqrt{2}} \right) = \frac{1}{2} \left( \langle 00 | 01 \rangle + \langle 00 | 10 \rangle - \langle 11 | 01 \rangle + \langle 11 | 10 \rangle \right) = 0$$

$$\langle B_{01} | B_{10} \rangle = \left( \frac{\langle 00 | - \langle 11 |}{\sqrt{2}} \right) \left( \frac{|01 \rangle + |10 \rangle}{\sqrt{2}} \right) = \frac{1}{2} \left( \langle 00 | 01 \rangle + \langle 00 | 10 \rangle - \langle 11 | 01 \rangle + \langle 11 | 10 \rangle \right) = 0$$

$$\langle B_{10} | B_{10} \rangle = \left( \frac{\langle 00 | - \langle 11 |}{\sqrt{2}} \right) \left( \frac{|01 \rangle + |10 \rangle}{\sqrt{2}} \right) = \frac{1}{2} \left( \langle 01 | 01 \rangle + \langle 01 | 10 \rangle + \langle 10 | 01 \rangle + \langle 10 | 10 \rangle \right) = 1$$

$$\langle B_{10} | B_{11} \rangle = \left( \frac{\langle 01 | + \langle 10 |}{\sqrt{2}} \right) \left( \frac{|01 \rangle + |10 \rangle}{\sqrt{2}} \right) = \frac{1}{2} \left( \langle 01 | 01 \rangle - \langle 01 | 10 \rangle + \langle 10 | 01 \rangle - \langle 10 | 10 \rangle \right) = 0$$

$$\langle B_{11} | B_{11} \rangle = \left( \frac{\langle 01 | - \langle 10 |}{\sqrt{2}} \right) \left( \frac{|01 \rangle - |10 \rangle}{\sqrt{2}} \right) = \frac{1}{2} \left( \langle 01 | 01 \rangle - \langle 01 | 10 \rangle + \langle 10 | 01 \rangle - \langle 10 | 10 \rangle \right) = 1$$

$$\langle B_{11} | B_{11} \rangle = \left( \frac{\langle 01 | - \langle 10 |}{\sqrt{2}} \right) \left( \frac{|01 \rangle - |10 \rangle}{\sqrt{2}} \right) = \frac{1}{2} \left( \langle 01 | 01 \rangle - \langle 01 | 10 \rangle - \langle 10 | 01 \rangle + \langle 10 | 10 \rangle \right) = 1$$

so orthonormality is verified. We know show that it is a basis. Recall that we can write any vector  $|\psi\rangle$  in the 2 qubit state space as:

$$|\psi\rangle = \alpha|00\rangle + \beta|01\rangle + \gamma|10\rangle + \delta|11\rangle$$

where  $\alpha, \beta, \gamma, \delta \in \mathbb{C}$  and  $|\alpha|^2 + |\beta|^2 + |\gamma|^2 + |\delta|^2 = 1$ . We then observe that this is equivalent to:

$$\left(\frac{\alpha+\delta}{\sqrt{2}}\right)|B_{00}\rangle + \left(\frac{\alpha-\delta}{\sqrt{2}}\right)|B_{01}\rangle + \left(\frac{\beta+\gamma}{\sqrt{2}}\right)|B_{10}\rangle + \left(\frac{\beta-\gamma}{\sqrt{2}}\right)|B_{11}\rangle \quad (*)$$

as:

$$\begin{split} \left(\frac{\alpha+\delta}{\sqrt{2}}\right) \frac{|00\rangle + |11\rangle}{\sqrt{2}} + \left(\frac{\alpha-\delta}{\sqrt{2}}\right) \frac{|00\rangle - |11\rangle}{\sqrt{2}} + \left(\frac{\beta+\gamma}{\sqrt{2}}\right) \frac{|01\rangle + |10\rangle}{\sqrt{2}} + \left(\frac{\beta-\gamma}{\sqrt{2}}\right) \frac{|01\rangle - |10\rangle}{\sqrt{2}} \\ &= \left(\frac{\alpha}{2} + \frac{\alpha}{2} + \frac{\delta}{2} - \frac{\delta}{2}\right) |00\rangle + \left(\frac{\alpha}{2} - \frac{\alpha}{2} + \frac{\delta}{2} + \frac{\delta}{2}\right) |11\rangle \\ &+ \left(\frac{\beta}{2} + \frac{\beta}{2} + \frac{\gamma}{2} - \frac{\gamma}{2}\right) |01\rangle + \left(\frac{\beta}{2} - \frac{\beta}{2} + \frac{\gamma}{2} + \frac{\gamma}{2}\right) |10\rangle \\ &= \alpha |00\rangle + \beta |01\rangle + \gamma |10\rangle + \delta |11\rangle = |\psi\rangle \end{split}$$

Hence (\*) shows that the Bell states form a basis.

#### Exercise 2.70

Suppose E is any positive operator acting on Alice's qubit Show that  $\langle \psi | E \otimes I | \psi \rangle$  takes the same value when  $|\psi \rangle$  is any of the four Bell states. Suppose some malevolent third party ('Eve') intercepts Alice's qubit on the way to Bob in the superdense coding protocol. Can Eve infer anything about which of the four possible bit strings 00,01,10,11 Alice is trying to send? If so, how, or if not, why not?

### Solution

Concepts Involved: Linear Algebra, Superdense Coding, Quantum Measurement

Let E be a positive operator. We have that E for a single qubit can be written as a linear combination of the Pauli matrices:

$$E = a_1 I + a_2 X + a_3 Y + a_4 Z$$

To see that this is the case, consider that the vector space of linear operators acting on a single qubit has dimension 4 (one easy way to see this is that the matrix representations of these operators have 4 entries). Hence, any set of 4 linearly independent linear operators form a basis for the space. As I, X, Y, Z are linearly independent, it follows that they form a basis of the space of linear operators on one qubit. Hence any E can be written as above (Remark: the above decomposition into Paulis is possible regardess of whether E is positive or not).

We then have that:

$$\langle B_{00}|E \otimes I|B_{00}\rangle = \left(\frac{\langle 00| + \langle 11|}{\sqrt{2}}\right) E \otimes I \left(\frac{|00\rangle + |11\rangle}{\sqrt{2}}\right)$$

$$= \left(\frac{\langle 00| + \langle 11|}{\sqrt{2}}\right) (a_1 I + a_2 X + a_3 Y + a_4 Z) \otimes I \left(\frac{|00\rangle + |11\rangle}{\sqrt{2}}\right)$$

$$= \left(\frac{\langle 00| + \langle 11|}{\sqrt{2}}\right) \left(a_1 \frac{|00\rangle + |11\rangle}{\sqrt{2}} + a_2 \frac{|10\rangle + |01\rangle}{\sqrt{2}} + a_3 \frac{i|10\rangle - i|01\rangle}{\sqrt{2}} + a_4 \frac{|00\rangle - |11\rangle}{\sqrt{2}}\right)$$

$$= \frac{1}{2} (a_1 + a_1 + a_4 - a_4)$$

$$= a_1$$

where in the second last equality we use the orthonormality of  $\{|00\rangle, |01\rangle, |10\rangle, |11\rangle\}$ . Repeating the same process for the other Bell states, we have:

$$\langle B_{01}|E \otimes I|B_{01}\rangle = \left(\frac{\langle 00| - \langle 11|}{\sqrt{2}}\right) (a_{1}I + a_{2}X + a_{3}Y + a_{4}Z) \otimes I\left(\frac{|00\rangle - |11\rangle}{\sqrt{2}}\right)$$

$$= \left(\frac{\langle 00| - \langle 11|}{\sqrt{2}}\right) \left(a_{1}\frac{|00\rangle - |11\rangle}{\sqrt{2}} + a_{2}\frac{|10\rangle - |01\rangle}{\sqrt{2}} + a_{3}\frac{i|10\rangle + i|01\rangle}{\sqrt{2}} + a_{4}\frac{|00\rangle + |11\rangle}{\sqrt{2}}\right)$$

$$= \frac{1}{2} (a_{1} + a_{1} + a_{4} - a_{4})$$

$$= a_{1}$$

$$\langle B_{10}|E \otimes I|B_{10}\rangle = \left(\frac{\langle 01| + \langle 10| }{\sqrt{2}}\right) \left(a_{1}I + a_{2}X + a_{3}Y + a_{4}Z\right) \otimes I\left(\frac{|01\rangle + |10\rangle}{\sqrt{2}}\right)$$

$$= \left(\frac{\langle 01| + \langle 10| }{\sqrt{2}}\right) \left(a_{1}\frac{|01\rangle + |10\rangle}{\sqrt{2}} + a_{2}\frac{|11\rangle + |00\rangle}{\sqrt{2}} + a_{3}\frac{i|11\rangle - i|00\rangle}{\sqrt{2}} + a_{4}\frac{|01\rangle - |10\rangle}{\sqrt{2}}\right)$$

$$= \frac{1}{2} \left(a_{1} + a_{1} + a_{4} - a_{4}\right)$$

$$= a_{1}$$

$$\langle B_{01}|E \otimes I|B_{01}\rangle = \left(\frac{\langle 01| - \langle 10|}{\sqrt{2}}\right) (a_{1}I + a_{2}X + a_{3}Y + a_{4}Z) \otimes I\left(\frac{|01\rangle - |10\rangle}{\sqrt{2}}\right)$$

$$= \left(\frac{\langle 01| - \langle 10|}{\sqrt{2}}\right) \left(a_{1}\frac{|01\rangle - |10\rangle}{\sqrt{2}} + a_{2}\frac{|11\rangle - |00\rangle}{\sqrt{2}} + a_{3}\frac{i|11\rangle + i|00\rangle}{\sqrt{2}} + a_{4}\frac{|01\rangle + |10\rangle}{\sqrt{2}}\right)$$

$$= \frac{1}{2} (a_{1} + a_{1} + a_{4} - a_{4})$$

$$= a_{1}$$

Now, suppose that Eve intercepts Alice's qubit. Eve cannot infer anything about which of the four possible bit strings that Alice is trying to send, as any single-qubit measurement that Eve can perform on the intercepted qubit will return the value:

$$\langle \psi | M^{\dagger} M \otimes I | \psi \rangle$$

Where M is the (single-qubit) measurement operator. But,  $M^{\dagger}M$  is positive, so by the above argument, the measurement outcome will be the same regardless of which Bell state  $|\psi\rangle$  is. Hence, Eve cannot obtain the information about the bit string.

## Exercise 2.71: Criterion to decide if a state is mixed or pure

Let  $\rho$  be a density operator. Show that  $\operatorname{tr}(\rho^2) \leq 1$ , with equality if and only if  $\rho$  is a pure state.

#### Solution

**Concepts Involved:** Linear Algebra, Trace, Density Operators, Pure States, Mixed States. Recall that a density operator  $\rho$  is pure if:

$$\rho = |\psi\rangle\langle\psi|$$

for some normalized quantum state vector  $|\psi\rangle$ .

Since  $\rho$  is a positive operator, by the spectral decomposition we have that:

$$\rho = \sum_{i} p_i |i\rangle\langle i|$$

where  $p_i \ge 0$  (due to positivity) and  $|i\rangle$  are orthonormal. Furthermore, by the property of density operators,

we have that  $tr(\rho) = 1$ , hence:

$$\operatorname{tr}(\rho) = \operatorname{tr}\left(\sum_{i} p_{i} |i\rangle\langle i|\right) = \sum_{i} p_{i} \operatorname{tr}(|i\rangle\langle i|) = \sum_{i} p_{i} = 1$$

where in the second equality we use the linearity of the trace. We obtain that  $0 \le p_i \le 1$  for each i. Calculating  $\rho^2$ , we have that:

$$\rho^2 = \left(\sum_i p_i |i\rangle\langle i|\right) \left(\sum_{i'} p_{i'} |i'\rangle\langle i'|\right) = \sum_i \sum_{i'} p_i p_{i'} |i\rangle\langle i| i'\rangle\langle i| = \sum_i \sum_{i'} p_i p_{i'} |i\rangle\langle i'| \,\delta_{ii'} = \sum_i p_i^2 |i\rangle\langle i|$$

Hence:

$$\operatorname{tr}\left(\rho^{2}\right) = \sum_{i} p_{i}^{2} \operatorname{tr}\left(|i\rangle\langle i|\right) = \sum_{i} p_{i}^{2} \leq \sum_{i} p_{i} = 1$$

where in the inequality we use the fact that  $p_i^2 \leq p_i$  as  $0 \leq p_i \leq 1$ . The inequality becomes an equality when  $p_i^2 = p_i$ , that is, when  $p_i = 0$  or  $p_i = 1$ . In order for  $\operatorname{tr}(\rho) = 1$  to hold, we have that  $p_i = 1$  for one i and zero for all others. Hence,  $\rho$  in this case is a pure state. Conversely, suppose  $\rho$  is a pure state. Then:

$$\operatorname{tr}(\rho^2) = \operatorname{tr}(|\psi\rangle\langle\psi|\psi\rangle\langle\psi|) = \operatorname{tr}(|\psi\rangle\langle\psi|) = 1.$$

## Exercise 2.72: Bloch sphere for mixed states

The Bloch sphere picture for pure states of a single qubit was introduced in Section 1.2. This description has an important generalization to mixed states as follows.

(1) Show that an arbitrary density matrix for a mixed state qubit may be written as

$$\rho = \frac{I + \mathbf{r} \cdot \boldsymbol{\sigma}}{2},$$

Where  $\mathbf{r}$  is a real three-dimensional vector such that  $\|\mathbf{r}\| \leq 1$ . This vector is known as the *Bloch vector* for the state  $\rho$ .

- (2) What is the Bloch vector representation for the state  $\rho = I/2$ ?
- (3) Show that a state  $\rho$  is pure if and only if  $\|\mathbf{r}\| = 1$ .
- (4) Show that for pure states the description of the Bloch vector we have given coincides with that in Section 1.2.

## Solution

Concepts Involved: Linear Algebra, Trace, Density Operators, Pure States, Mixed States

(1) Since  $\{I, X, Y, Z\}$  form an basis the vector space of single-qubit linear operators, we can write (for any  $\rho$ , regardless of whether it is a density operator or not):

$$\rho = a_1 I + a_2 X + a_3 Y + a_4 Z$$

for constants  $a_1, a_2, a_3, a_4 \in \mathbb{C}$ . Since  $\rho$  is a Hermitian operator, we find that each of these constants are actually real, as:

$$a_1I + a_2X + a_3Y + a_4Z = \rho = \rho^{\dagger} = a_1^*I^{\dagger} + a_2^*X^{\dagger} + a_3^*Y^{\dagger} + a_4^*Z^{\dagger} = a_1^*I + a_2^*X + a_3^*Y + a_4^*Z$$

Now, we require that  $tr(\rho) = 1$  for any density operator, hence:

$$\operatorname{tr}(\rho) = \operatorname{tr}(a_1 I + a_2 X + a_3 Y + a_4 Z) = a_1 \operatorname{tr}(I) + a_2 \operatorname{tr}(X) + a_3 \operatorname{tr}(Y) + a_4 \operatorname{tr}(Z) = 2a_1 = 1$$

from which we obtain that  $a_1 = \frac{1}{2}$ . Note that in the second equality we use the linearity of the trace, and in the third equality we use that tr(I) = 2 and  $tr(\sigma_i) = 0$  for  $i \in \{1, 2, 3\}$  (Exercise 2.36). Calculating  $\rho^2$ , we have that:

$$\begin{split} \rho^2 &= \frac{1}{4}I + \frac{a_2}{2}X + \frac{a_3}{2}Y + \frac{a_4}{2}Z + \frac{a_2}{2}X + a_2^2X^2 + a_2a_3XY + a_2a_4XZ \\ &+ \frac{a_3}{2}Y + a_3a_2YX + a_3^2Y^2 + a_3a_4YZ + \frac{a_4}{2}Z + a_4a_2ZX + a_4a_3ZY + a_4^2Z^2 \end{split}$$

Now, using that  $\{\sigma_i, \sigma_j\} = 0$  for  $i, j \in \{1, 2, 3\}$ ,  $i \neq j$  and that  $\sigma_i^2 = I$  for any  $i \in \{1, 2, 3\}$  (Exercise 2.41), the above simplifies to:

$$\rho^2 = \left(\frac{1}{4} + a_2^2 + a_3^2 + a_4^2\right)I + a_2X + a_3Y + a_4Z$$

Taking the trace of  $\rho^2$  we have that:

$$\operatorname{tr}\left(\rho^{2}\right) = \left(\frac{1}{4} + a_{2}^{2} + a_{3}^{2} + a_{4}^{2}\right)\operatorname{tr}(I) + a_{2}\operatorname{tr}(X) + a_{3}\operatorname{tr}(Y) + a_{4}\operatorname{tr}(Z) = 2\left(\frac{1}{4} + a_{2}^{2} + a_{3}^{2} + a_{4}^{2}\right)$$

From the previous exercise (Exercise 2.71) we know that  $tr(\rho^2) \le 1$ , so:

$$2\left(\frac{1}{4}+a_2^2+a_3^2+a_4^2\right) \leq 1 \implies a_2^2+a_3^2+a_4^2 \leq \frac{1}{4} \implies \sqrt{a_2^2+a_3^2+a_4^2} \leq \frac{1}{2}$$

Hence we can write:

$$\rho = \frac{I + r_x X + r_y Y + r_z Z}{2} = \frac{I + \mathbf{r} \cdot \boldsymbol{\sigma}}{2}$$

with  $\|\mathbf{r}\| \leq 1$ .

(2) The Bloch representation for the state  $\rho = \frac{I}{2}$  is the above form with  $\mathbf{r} = \mathbf{0}$ . This vector corresponds to the center of the Bloch sphere, which is a maximally mixed state  $(\operatorname{tr}(\rho^2))$  is minimized, with  $\operatorname{tr}(\rho^2) = \frac{1}{2}$ .

(3) From the calculation in part (1), we know that for any  $\rho$ :

$$\operatorname{tr}\!\left(\rho^2\right) = 2\left(\frac{1 + r_x^2 + r_y^2 + r_z^2}{4}\right) = \frac{1 + r_x^2 + r_y^2 + r_z^2}{2}$$

if  $\|\mathbf{r}\|=1$ , then  $r_x^2+r_y^2+r_z^2=1$ . Hence,  $\mathrm{tr}\big(\rho^2\big)=1$  and  $\rho$  is pure by Exercise 2.71. Conversely, suppose  $\rho$  is pure. Then,  $\mathrm{tr}\big(\rho^2\big)=1$ , so:

$$\frac{1 + r_x^2 + r_y^2 + r_z^2}{2} = 1 \implies r_x^2 + r_y^2 + r_z^2 = 1 \implies \|\mathbf{r}\| = 1.$$

(4) In section 1.2, we looked at states that lie on the surface of the Bloch sphere, which we parameterized as:

$$|\psi\rangle = \cos\!\left(\frac{\theta}{2}\right)\!|0\rangle + e^{i\varphi}\sin\!\left(\frac{\theta}{2}\right)\!|1\rangle.$$

Calculating the density operator corresponding to  $|\psi\rangle$ , we have:

$$\begin{split} \rho &= |\psi\rangle\!\langle\psi| = \cos^2\left(\frac{\theta}{2}\right) \, |0\rangle\!\langle 0| + \cos\left(\frac{\theta}{2}\right) \sin\left(\frac{\theta}{2}\right) e^{-i\varphi} \, |0\rangle\!\langle 1| \\ &+ \cos\left(\frac{\theta}{2}\right) \sin\left(\frac{\theta}{2}\right) e^{i\varphi} \, |1\rangle\!\langle 0| + \sin^2\left(\frac{\theta}{2}\right) \, |1\rangle\!\langle 1| \\ &= \cos^2\left(\frac{\theta}{2}\right) \, |0\rangle\!\langle 0| + \frac{\sin(\theta) e^{-i\varphi}}{2} \, |0\rangle\!\langle 1| + \frac{\sin(\theta) e^{i\varphi}}{2} \, |1\rangle\!\langle 0| + \sin^2\left(\frac{\theta}{2}\right) \, |1\rangle\!\langle 1| \end{split}$$

Conversely, we have that (in the computational basis) our proposed form of  $\rho = \frac{I + \mathbf{r} \cdot \boldsymbol{\sigma}}{2}$  can be represented as:

$$\rho = \frac{1+r_z}{2} \left| 0 \right\rangle\!\left\langle 0 \right| + \frac{r_x - ir_y}{2} \left| 0 \right\rangle\!\left\langle 1 \right| + \frac{r_x + ir_y}{2} \left| 1 \right\rangle\!\left\langle 0 \right| + \frac{1-r_z}{2} \left| 1 \right\rangle\!\left\langle 1 \right|$$

Solving for  $r_x, r_y, r_z$  by equating the two expressions for  $\rho$  (using Euler's formula and  $\sin(2\theta) = 2\cos(\theta)\sin(\theta)$ ), we have:

$$r_x = \cos(\varphi)\sin(\theta), \quad r_y = \sin(\varphi)\sin(\theta), \quad r_z = 2\cos^2\left(\frac{\theta}{2}\right) - 1 = \cos(\theta)$$

Calculating  $\|\mathbf{r}\|$  we have that:

$$\|\mathbf{r}\| = \sqrt{r_x^2 + r_y^2 + r_z^2} = \sqrt{\cos^2(\varphi)\sin^2(\theta) + \sin^2(\varphi)\cos^2(\theta) + \cos^2(\theta)}$$
$$= \sqrt{\sin^2(\theta) + \cos^2(\theta)}$$
$$= 1$$

so we see that indeed, the two definitions coincide for pure states (as  $\|\mathbf{r}\|=1$ ).

#### Exercise 2.73

(\*) Let  $\rho$  be a density operator. A minimal ensemble for  $\rho$  is an ensemble  $\left\{p_i,\,|\psi_i\rangle\right\}$  containing a number of elements equal to the rank of  $\rho$ . Let  $|\psi\rangle$  be any state in the support of  $\rho$ . (The support of a Hermitian operator A is the vector space spanned by the eigenvectors of A with non-zero eigenvalues.) Show that there is a minimal ensemble for  $\rho$  that contains  $|\psi\rangle$ , and moreover that in any such ensemble  $|\psi\rangle$  must appear with probability

$$p_i = \frac{1}{\langle \psi_i | \rho^{-1} | \psi_i \rangle},$$

where  $\rho^{-1}$  is defined to be the inverse of  $\rho$ , when  $\rho$  is considered as an operator acting only on the support of  $\rho$ . (This definition removes the problem that  $\rho$  may not have an inverse.)

#### Solution

**Concepts Involved:** Below, we will use the unitary freedom in the ensemble for density matrices which is also known as Uhlmann's theorem. Specifically recall that  $\rho = \sum_i p_i |\psi_i\rangle \langle \psi_i| = \sum_j q_j |\varphi_j\rangle \langle \varphi_j|$  for ensembles  $\{p_i, |\psi_i\rangle\}$  and  $\{q_j, |\varphi_j\rangle\}$  if and only if

$$\sqrt{p_i} \left| \psi_i \right\rangle = \sum_j u_{ij} \sqrt{q_j} \left| \varphi_j \right\rangle$$

for some unitary matrix  $u_{ij}$ .

Using the spectral decomposition of the density matrix we have

$$ho = \sum_{k=1}^r \lambda_k |k
angle \langle k| \quad ext{ with } \lambda_k > 0$$

where all the eigenvectors with eigenvalue 0 have been removed. Thus, the set of vectors  $S = \{|k\rangle\}_{k=1}^r$  forms a spanning set set for the support of  $\rho$ . An element in the support of  $\rho$  can thus be decomposed as

$$|\psi_i\rangle = \sum_k c_{ik} |k\rangle = \sum_k \langle k |\psi_i\rangle |k\rangle$$

Assuming that  $|\psi_i\rangle$  occurs with probability  $p_i$ , we can use the Uhlmann's theorem quoted above to arrive at the relation

$$\sqrt{p_i}|\psi_i\rangle \stackrel{?}{=} \sum_k u_{ik} \sqrt{\lambda_k} |k\rangle = \sqrt{p_i} \sum_k \langle k|\psi_i\rangle |k\rangle,$$

which allows us relate the elements of one of the columns (i th) of the unitary matrix to

$$u_{ik}\sqrt{\lambda_k} \stackrel{?}{=} \sqrt{p_i}\langle k|\psi_i\rangle.$$

Such a relation can always be satisfied for a unitary matrix with dimension r. As u is unitary, we have

$$\sum_{k} |u_{ik}|^{2} = 1 \implies 1 = \sum_{k} \frac{p_{i} \langle \psi_{i} | k \rangle \langle k | \psi_{i} \rangle}{\lambda_{k}}$$

$$\implies p_{i} = \sum_{k} \frac{\lambda_{k}}{\langle \psi_{i} | k \rangle \langle k | \psi_{i} \rangle}$$

$$= \frac{1}{\langle \psi_{i} | \sum_{k} \frac{1}{\lambda_{k}} | k \rangle \langle k | \psi_{i} \rangle}$$

$$= \frac{1}{\langle \psi_{i} | \rho^{-1} | \psi_{i} \rangle}.$$

$$\begin{aligned} |\psi_j\rangle &= \rho \rho^{-1} |\psi_i\rangle \\ &:= \sum_{i=1}^r p_i |\psi_i\rangle \langle \psi_i | \rho^{-1} |\psi_j\rangle \\ &= \sum_{i=1}^r p_i \langle \psi_i | \rho^{-1} |\psi_j\rangle |\psi_i\rangle. \end{aligned}$$

But now note that  $\{|\psi_j\rangle\}_{j=1}^r$  are linearly independent and  $|\psi_j\rangle := \sum_{i=1}^r \delta_{ij} |\psi_i\rangle$ .

$$\implies p_i \langle \psi_i | \rho^{-1} | \psi_i \rangle = 1.$$

Thus, the probability associated with the state  $|\psi_i\rangle$  in the ensemble is given by

$$p_i = \frac{1}{\langle \psi_i | \rho^{-1} | \psi_i \rangle} \ .$$

## Exercise 2.74

Suppose a composite of systems A and B is in the state  $|a\rangle$   $|b\rangle$ , where  $|a\rangle$  is a pure state of system A, and  $|b\rangle$  is a pure state of system B. Show that the reduced density operator of system A alone in a pure state.

## Solution

**Concepts Involved:** Linear Algebra, Density Operators, Reduced Density Operators, Partial Trace, Pure States.

Suppose we have  $|a\rangle|b\rangle\in A\otimes B$ . Then, the density operator of the combined system is given as  $\rho^{AB}=(|a\rangle|b\rangle)(\langle a|\langle b|)=|a\rangle\!\langle a|\otimes|b\rangle\!\langle b|.$  Calculating the reduced density operator of system A by tracing out system B, we have

$$\rho^{A} = \operatorname{tr}_{B}(\rho_{AB}) = \operatorname{tr}_{B}(|a\rangle\langle a| \otimes |b\rangle\langle b|) = |a\rangle\langle a| \operatorname{tr}(|b\rangle\langle b|) = |a\rangle\langle a| \langle b|b\rangle = |a\rangle\langle a|.$$

Hence we find that  $\rho^A = |a\rangle\langle a|$  is indeed a pure state.

## Exercise 2.75

For each of the four Bell states, find the reduced density operator for each qubit.

### Solution

Concepts Involved: Linear Algebra, Density Operators, Reduced Density Operators, Partial Trace.

For the bell state  $|B_{00}\rangle$ , we have the density operator:

$$\rho = \left(\frac{|00\rangle + |11\rangle}{\sqrt{2}}\right) \left(\frac{\langle 00| + \langle 11|}{\sqrt{2}}\right) = \frac{|00\rangle\langle 00| + |11\rangle\langle 00| + |00\rangle\langle 11| + |11\rangle\langle 11|}{2}$$

Obtaining the reduced density operator for qubit A, we have:

$$\rho^{A} = \operatorname{tr}_{B}(\rho) = \frac{\operatorname{tr}_{B}(|00\rangle\langle00|) + \operatorname{tr}_{B}(|00\rangle\langle11|) + \operatorname{tr}_{B}(|11\rangle\langle00|) + \operatorname{tr}_{B}(|11\rangle\langle11|)}{2}$$

$$= \frac{|0\rangle\langle0|\langle0|0\rangle + |0\rangle\langle1|\langle1|0\rangle + |1\rangle\langle0|\langle0|1\rangle + |1\rangle\langle1|\langle1|1\rangle}{2}$$

$$= \frac{|0\rangle\langle0| + |1\rangle\langle1|}{2}$$

$$= \frac{I}{2}$$

Obtaining the reduced density operator for qubit B, we have:

$$\begin{split} \rho^B &= \operatorname{tr}_A(\rho) = \frac{\operatorname{tr}_A(\,|00\rangle\!\langle 00|) + \operatorname{tr}_A(\,|00\rangle\!\langle 11|) + \operatorname{tr}_A(\,|11\rangle\!\langle 00|) + \operatorname{tr}_A(\,|11\rangle\!\langle 11|)}{2} \\ &= \frac{\langle 0|0\rangle\,|0\rangle\!\langle 0| + \langle 1|0\rangle\,|0\rangle\!\langle 1| + \langle 0|1\rangle\,|1\rangle\!\langle 0| + \langle 1|1\rangle\,|1\rangle\!\langle 1|}{2} \\ &= \frac{|0\rangle\!\langle 0| + \,|1\rangle\!\langle 1|}{2} \\ &= \frac{I}{2} \end{split}$$

We repeat a similar process for the other four bell states. For  $|B_{01}\rangle$ , we have:

$$\begin{split} \rho &= \left(\frac{|00\rangle - |11\rangle}{\sqrt{2}}\right) \left(\frac{\langle 00| - \langle 11|}{\sqrt{2}}\right) = \frac{|00\rangle\langle 00| - |00\rangle\langle 11| - |11\rangle\langle 00| + |11\rangle\langle 11|}{2} \\ \rho^A &= \frac{\operatorname{tr}_B(|00\rangle\langle 00|) - \operatorname{tr}_B(|00\rangle\langle 11|) - \operatorname{tr}_B(|11\rangle\langle 00|) + \operatorname{tr}_B(|11\rangle\langle 11|)}{2} \\ &= \frac{|0\rangle\langle 0| \langle 0|0\rangle - |0\rangle\langle 1| \langle 1|0\rangle - |1\rangle\langle 0| \langle 0|1\rangle + |1\rangle\langle 1| \langle 1|1\rangle}{2} \\ &= \frac{I}{2} \\ \rho^B &= \frac{\operatorname{tr}_A(|00\rangle\langle 00|) - \operatorname{tr}_A(|00\rangle\langle 11|) - \operatorname{tr}_A(|11\rangle\langle 00|) + \operatorname{tr}_A(|11\rangle\langle 11|)}{2} \\ &= \frac{\langle 0|0\rangle |0\rangle\langle 0| + \langle 1|0\rangle |0\rangle\langle 1| + \langle 0|1\rangle |1\rangle\langle 0| + \langle 1|1\rangle |1\rangle\langle 1|}{2} \\ &= \frac{I}{2} \end{split}$$

For  $|B_{10}\rangle$ , we have:

$$\rho = \left(\frac{|01\rangle + |10\rangle}{\sqrt{2}}\right) \left(\frac{\langle 01| + \langle 10|}{\sqrt{2}}\right) = \frac{|01\rangle\langle 01| + |01\rangle\langle 10| + |10\rangle\langle 01| + |10\rangle\langle 10|}{2}$$

$$\rho^{A} = \frac{\operatorname{tr}_{B}(|01\rangle\langle 01|) + \operatorname{tr}_{B}(|01\rangle\langle 10|) + \operatorname{tr}_{B}(|10\rangle\langle 01|) + \operatorname{tr}_{B}(|10\rangle\langle 10|)}{2}$$

$$= \frac{|0\rangle\langle 0|\langle 1| + \rangle |0\rangle\langle 1|\langle 0|1\rangle + |1\rangle\langle 0|\langle 1|0\rangle + |1\rangle\langle 1|\langle 0|0\rangle}{2}$$

$$= \frac{I}{2}$$

$$\rho^{B} = \frac{\operatorname{tr}_{A}(|01\rangle\langle 01|) + \operatorname{tr}_{A}(|01\rangle\langle 10|) + \operatorname{tr}_{A}(|10\rangle\langle 01|) + \operatorname{tr}_{A}(|10\rangle\langle 10|)}{2}$$

$$= \frac{\langle 0|0\rangle |1\rangle\langle 1| + \langle 1|0\rangle |1\rangle\langle 0| + \langle 0|1\rangle |0\rangle\langle 1| + \langle 1|1\rangle |1\rangle\langle 1|}{2}$$

$$= \frac{I}{2}$$

Finally, for  $|B_{11}\rangle$  we have:

$$\rho = \left(\frac{|01\rangle - |10\rangle}{\sqrt{2}}\right) \left(\frac{\langle 01| - \langle 10|}{\sqrt{2}}\right) = \frac{|01\rangle\langle 01| - |01\rangle\langle 10| - |10\rangle\langle 01| + |10\rangle\langle 10|}{2}$$

$$\rho^{A} = \frac{\operatorname{tr}_{B}(|01\rangle\langle 01|) - \operatorname{tr}_{B}(|01\rangle\langle 10|) - \operatorname{tr}_{B}(|10\rangle\langle 01|) + \operatorname{tr}_{B}(|10\rangle\langle 10|)}{2}$$

$$= \frac{|0\rangle\langle 0|\langle 1| - \rangle |0\rangle\langle 1|\langle 0|1\rangle - |1\rangle\langle 0|\langle 1|0\rangle + |1\rangle\langle 1|\langle 0|0\rangle}{2}$$

$$= \frac{I}{2}$$

$$\rho^{B} = \frac{\operatorname{tr}_{A}(|01\rangle\langle 01|) - \operatorname{tr}_{A}(|01\rangle\langle 10|) - \operatorname{tr}_{A}(|10\rangle\langle 01|) + \operatorname{tr}_{A}(|10\rangle\langle 10|)}{2}$$

$$= \frac{\langle 0|0\rangle |1\rangle\langle 1| - \langle 1|0\rangle |1\rangle\langle 0| - \langle 0|1\rangle |0\rangle\langle 1| + \langle 1|1\rangle |1\rangle\langle 1|}{2}$$

$$= \frac{I}{2}$$

## Exercise 2.76

Extend the proof of the Schmidt decomposition to the case where A and B may have state space of different dimensionality.

#### Solution

Concepts Involved: Schmidt Decomposition, Singular Value Decomposition.

Note that for this problem we will use a more general form of the Singular Value Decomposition than proven in Nielsen and Chuang (that may have been encountered in a linear algebra course). Given an arbitrary  $m \times n$  rectangular matrix A, there exists an  $m \times m$  unitary matrix U and  $n \times n$  unitary matrix V such that  $V = U \cap V$  where V = u is a V = u rectangular diagonal matrix with non-negative reals on the diagonal (see https://en.wikipedia.org/wiki/Singular\_value\_decomposition).

Let  $|m\rangle, |n\rangle$  be orthonormal bases for A and B. We can then write:

$$A = \sum_{mn} a_{mn} |m\rangle |n\rangle$$

for some  $m \times n$  matrix of complex numbers a. Using the generalized SVD, we can write:

$$A = \sum_{min} u_{mi} d_{ii} v_{in} |m\rangle |n\rangle$$

where  $d_{ii}$  is a rectangular diagonal matrix. We can then define  $|i_A\rangle = \sum_m u_{mi} |m\rangle$ ,  $|i_B\rangle = \sum_n u_{in} |n\rangle$ , and  $\lambda_i = d_{ii}$  to yield the Schmidt decomposition. Note that we take  $i = \min(m, n)$  and our sum only has as many terms as the dimensionality of the smaller space.

## Exercise 2.77

(\*) Suppose ABC is a three component quantum system. Show by example that there are quantum states  $|\psi\rangle$  of such systems which can not be written in the form

$$|\psi\rangle = \sum_{i} \lambda_i |i_A\rangle |i_B\rangle |i_C\rangle$$

where  $\lambda_i$  are real numbers, and  $|i_A\rangle, |i_B\rangle, |i_C\rangle$  are orthonormal bases of the respective systems.

#### Solution

Concepts Involved: Linear Algebra, Schmidt Decomposition.

Consider the state:

$$|\psi\rangle = |0\rangle \otimes |B_{00}\rangle = \frac{|000\rangle + |011\rangle}{\sqrt{2}}$$

we claim that this state cannot be written in the form:

$$|\psi\rangle = \sum_{i} \lambda_i |i_A\rangle |i_B\rangle |i_C\rangle$$

for orthonormal bases  $|i_A\rangle, |i_B\rangle, |i_C\rangle$ . Suppose for the sake of contradiction that we could write it in this form. We then make the observation that:

$$\rho^{A} = \operatorname{tr}_{BC}(|\psi\rangle\langle\psi|) = \sum_{i} \lambda_{i}^{2} |i_{A}\rangle\langle i_{A}|$$

$$\rho^{B} = \operatorname{tr}_{AC}(|\psi\rangle\langle\psi|) = \sum_{i} \lambda_{i}^{2} |i_{B}\rangle\langle i_{B}|$$

$$\rho^{C} = \operatorname{tr}_{AB}(|\psi\rangle\langle\psi|) = \sum_{i} \lambda_{i}^{2} |i_{C}\rangle\langle i_{C}|.$$

From this, we conclude that if it is possible to write  $|\psi\rangle$  in such a form, then the eigenvalues of the reduced density matrices must all agree and be equal to  $\lambda_i^2$ . Computing the density matrix of the proposted  $|\psi\rangle=|0\rangle\otimes|B_{00}\rangle$ , we have:

$$\rho = \frac{|000\rangle\langle000| + |000\rangle\langle011| + |011\rangle\langle000| + |011\rangle\langle011|}{2}$$

Computing the reduced density matrices  $\rho^A$  and  $\rho^B$ , we find that:

$$\rho^A = \operatorname{tr}_{BC}(\rho) = |0\rangle\langle 0|$$

$$\rho^B = \operatorname{tr}_{AC}(\rho) = \frac{|0\rangle\langle 0| + |1\rangle\langle 1|}{2}.$$

However, the former reduced density matrix has eigenvalues  $\lambda_1^2=1,\lambda_2^2=0$ , and the latter has  $\lambda_1^2=\frac{1}{2}$ ,  $\lambda_2^2=\frac{1}{2}$ . This contradicts the fact that the  $\lambda_i^2$ s must match.

**Remark:** Necessary and Sufficient conditions for the tripartite (and higher order) Schmidt decompositions can be found here https://arxiv.org/pdf/quant-ph/9504006.pdf.

#### Exercise 2.78

Prove that a state  $|\psi\rangle$  of a composite system AB is a product state if and only if it has a Schmidt number 1. Prove that  $|\psi\rangle$  is a product state if and only if  $\rho^A$  (and thus  $\rho^B$ ) are pure states.

#### Solution

**Concepts Involved:** Linear Algebra, Schmidt Decomposition, Schmidt Number, Reduced Density Operators.

Suppose  $|\psi\rangle$  is a product state. Then,  $|\psi\rangle=|0_A\rangle|0_B\rangle$  for some  $|0_A\rangle, |0_B\rangle$ , and we therefore have that  $|\psi\rangle$  has Schmidt number 1 (it is already written in Schmidt decomposition form, and has one nonzero  $\lambda$ ). Conversely, suppose  $|\psi\rangle$  has Schmidt number 1. Then,  $|\psi\rangle=1|0_A\rangle|0_B\rangle+0|1_A\rangle|1_B\rangle$  when writing  $|\psi\rangle$  in its Schmidt decomposition. Therefore,  $|\psi\rangle=|i_A\rangle|i_B\rangle$  and  $|\psi\rangle$  is a product state.

Next, take any  $|\psi\rangle$  and write out its Schmidt decomposition. We then get:

$$|\psi\rangle = \sum_{i} \lambda_i |i_A\rangle |i_B\rangle.$$

Hence:

$$\rho = \sum_{i} \lambda_{i}^{2} |i_{A}\rangle\langle i_{A}| \otimes |i_{B}\rangle\langle i_{B}|.$$

Taking the partial trace of  $\rho$  to obtain  $\rho^A$ , we have:

$$\rho^{A} = \operatorname{tr}_{B}(\rho) = \sum_{i} \lambda_{i}^{2} \operatorname{tr}_{B}(|i_{A}\rangle\langle i_{A}| \otimes |i_{B}\rangle\langle i_{B}|) = \sum_{i} \lambda_{i}^{2} |i_{A}\rangle\langle i_{A}| \operatorname{tr}(|i_{B}\rangle\langle i_{B}|) = \sum_{i} \lambda_{i}^{2} |i_{A}\rangle\langle i_{A}|.$$

Identically:

$$\rho^B = \operatorname{tr}_A(\rho) = \sum_i \lambda_i^2 |i_B\rangle\langle i_B|.$$

Now, suppose that  $|\psi\rangle$  is a product state. Then,  $|\psi\rangle$  has Schmidt number 1. Hence, only one of  $\lambda_1,\lambda_2$  is nonzero. Hence,  $\rho^A=|i_A\rangle\langle i_A|$  and  $\rho^B=|i_B\rangle\langle i_B|$ , so  $\rho^A,\rho^B$  are pure. Conversely, suppose  $\rho^A,\rho^B$  are pure. Then, we have that  $\rho^A=|i_A\rangle\langle i_A|$  and  $\rho^B=|i_B\rangle\langle i_B|$ , so it follows that one of  $\lambda_1,\lambda_2$  in the above equations for  $\rho^A,\rho^B$  must be zero. Therefore,  $|\psi\rangle$  has Schmidt number 1, and is hence a product state.

#### Exercise 2.79

Consider a composite system consisting of two qubits. Find the Schmidt decomposition of the states

$$\frac{|00\rangle+|11\rangle}{\sqrt{2}}; \quad \frac{|00\rangle+|01\rangle+|10\rangle+|11\rangle}{2}; \text{ and } \frac{|00\rangle+|01\rangle+|10\rangle}{\sqrt{3}}$$

#### Solution

Concepts Involved: Linear Algebra, Schmidt Decomposition, Reduced Density Matrices, Partial Trace.

For the first two expressions, by inspection we find that:

$$\frac{|00\rangle + |11\rangle}{\sqrt{2}} = \frac{1}{\sqrt{2}}|0\rangle|0\rangle + \frac{1}{\sqrt{2}}|1\rangle|1\rangle$$
$$\frac{|00\rangle + |01\rangle + |10\rangle + |11\rangle}{2} = 1|+\rangle|+\rangle + 0|-\rangle|-\rangle$$

For the third expression, we require a little more work. We first note that the existence of the Schmidt decomposition guarantees that the state  $|\psi\rangle=\frac{|00\rangle+|01\rangle+|10\rangle}{\sqrt{3}}$  can be written in the form  $|\psi\rangle=\sum_{i=1}^2\lambda_i|i_A\rangle|i_B\rangle$  for some choice of orthonormal bases  $|i_A\rangle,|i_B\rangle$ . By the definition of reduced density matrices/partial trace, we can make the observation that:

$$\rho^{A} = \operatorname{tr}_{B}(\rho) = \operatorname{tr}_{B}(|\psi\rangle\langle\psi|) = \sum_{i=1}^{2} \lambda_{i}^{2} |i_{A}\rangle\langle i_{A}| \operatorname{tr}_{B}(|i_{B}\rangle\langle i_{B}|) = \sum_{i=1}^{2} \lambda_{i}^{2} |i_{A}\rangle\langle i_{A}|$$

and similarly that  $\rho^B = \sum_{i=1}^2 \lambda_i^2 |i_B\rangle\langle i_B|$ . Hence, to find the Schmidt decomposition of  $|\psi\rangle$ , we can compute the reduced density matrices and then solve for their eigenvalues  $\lambda_i^2$  and eigenvectors  $|i\rangle$ . First solving for  $\rho$ , we have:

$$\rho = \frac{|00\rangle\langle00| + |00\rangle\langle01| + |00\rangle\langle10| + |01\rangle\langle00| + |01\rangle\langle01| + |01\rangle\langle10| + |10\rangle\langle00| + |10\rangle\langle01| + |10\rangle\langle10|}{3}$$

Solving for the reduced density matrix  $\rho^A$  we have:

$$\rho^{A} = \operatorname{tr}_{B}(\rho) = \frac{2|0\rangle\langle 0| + |0\rangle\langle 1| + |1\rangle\langle 0| + |1\rangle\langle 1|}{3} \cong \frac{1}{3} \begin{bmatrix} 2 & 1\\ 1 & 1 \end{bmatrix}$$

Solving for the eigenvalues and (normalized) eigenvectors, we have:

$$\lambda_1^2 = \frac{3+\sqrt{5}}{6}, |1_A\rangle = \frac{1}{\sqrt{10+2\sqrt{5}}} \left( (1+\sqrt{5})|0\rangle + 2|1\rangle \right)$$

$$\lambda_2^2 = \frac{3 - \sqrt{5}}{6}, |2_A\rangle = \frac{1}{\sqrt{10 - 2\sqrt{5}}} \left( (1 - \sqrt{5})|0\rangle + 2|1\rangle \right).$$

Next solving for the reduced density matrix  $\rho^B$ , we have:

$$\rho^B = \operatorname{tr}_A(\rho) = \frac{2|0\rangle\langle 0| + |0\rangle\langle 1| + |1\rangle\langle 0| + |1\rangle\langle 1|}{3} \cong \frac{1}{3} \begin{bmatrix} 2 & 1\\ 1 & 1 \end{bmatrix}.$$

This (of course) has the same eigenvectors and eigenvalues:

$$\lambda_1^2=\frac{3+\sqrt{5}}{6}, |1_B\rangle=\frac{1}{\sqrt{10+2\sqrt{5}}}\left((1+\sqrt{5})|0\rangle+2|1\rangle\right)$$

$$\lambda_2^2 = \frac{3 - \sqrt{5}}{6}, |2_B\rangle = \frac{1}{\sqrt{10 - 2\sqrt{5}}} \left( (1 - \sqrt{5})|0\rangle + 2|1\rangle \right).$$

Hence the Schmidt decomposition of  $|\psi\rangle$  is given by:

$$|\psi\rangle = \lambda_1 |1_A\rangle |1_B\rangle + \lambda_2 |2_A\rangle |2_B\rangle$$

where the expressions for the eigenvalues/eigenvectors are given above.

#### Exercise 2.80

Suppose  $|\psi\rangle$  and  $|\varphi\rangle$  are two pure states of a composite quantum system with components A and B, with identical Schmidt coefficients. Show that there are unitary transformations U on a system A and V on system B such that  $|\psi\rangle=(U\otimes V)|\varphi\rangle$ .

### Solution

Concepts Involved: Linear Algebra, Schmidt Decomposition, Unitary Operators.

We first prove a Lemma. Suppose we have two (orthonormal) bases  $\{|i\rangle\}$ ,  $\{|i'\rangle\}$  of a (n-dimensional) vector space A. We claim that the change of basis transformation U where  $|i'\rangle = U|i\rangle$  is unitary.

To see this is the case, let  $U=\sum_i|i'\rangle\langle i|$ . By orthonormality, we see that  $U|i\rangle=|i'\rangle$  as desired. Computing  $U^\dagger$ , we have that  $U^\dagger=\sum_i(|i'\rangle\langle i|)^\dagger=\sum_i|i\rangle\langle i'|$ . By orthonormality, we then see that  $U^\dagger U=\sum_i|i\rangle\langle i|=I$  and hence U is unitary.

We now move onto the actual problem. By assumption, we can write  $|\psi\rangle=\sum_i\lambda_i|i_A\rangle|i_B\rangle$  and  $|\varphi\rangle=\sum_j\lambda_j|j_A\rangle|j_B\rangle$  where  $\lambda_i=\lambda_j$  if i=j. By the lemma, there exists unitary change-of-basis matrices U,V such that  $|i_A\rangle=U|j_A\rangle$  and  $|i_B\rangle=V|j_B\rangle$ . Hence, we have that:

$$|\psi\rangle = \sum_{i} \lambda_{i} |i_{A}\rangle |i_{B}\rangle = \sum_{j} \lambda_{j} (U|j_{A}\rangle) (V|j_{B}\rangle) = (U \otimes V) \sum_{j} \lambda_{j} |j_{A}\rangle |j_{B}\rangle = (U \otimes V) |\varphi\rangle$$

which is what we wanted to prove.

### Exercise 2.81: Freedom in purifications

Let  $|AR_1\rangle$  and  $|AR_2\rangle$  be two purifications of a state  $\rho^A$  to a composite system AR. Prove that there exists a unitary transformation  $U_R$  acting on system R such that  $|AR_1\rangle = (I_A \otimes U_R)|AR_2\rangle$ .

#### Solution

Concepts Involved: Linear Algebra, Schmidt Decomposition, Purification, Unitary Operators

Let  $|AR_1\rangle, |AR_2\rangle$  be two purifications of  $\rho^A$  to a composite system AR. We can write the orthonormal decomposition of  $\rho^A$  as  $\rho^A = \sum_i p_i \, |i^A\rangle\langle i^A|$ , from which it follows that we can write:

$$|AR_1\rangle = \sum_{i} \sqrt{p_i} |i^A\rangle |i^R\rangle$$
$$|AR_2\rangle = \sum_{i} \sqrt{p_i} |i^A\rangle |i'^R\rangle$$

for some bases  $\{|i\rangle\}$ ,  $\{|i'\rangle\}$  of R. By the Lemma proven in the previous exercise, the transformation  $U_R$  such that  $|i\rangle = U_R|i'\rangle$  is unitary, so hence:

$$|AR_{1}\rangle = \sum_{i} \sqrt{p_{i}} |i^{A}\rangle |i^{R}\rangle = \sum_{i} \sqrt{p_{i}} |i^{A}\rangle (U_{R}|i'^{R}\rangle) = \sum_{i} \sqrt{p_{i}} (I_{A}|i^{A}\rangle) (U_{R}|i'^{R}\rangle)$$
$$= (I_{A} \otimes U_{R}) \sum_{i} \sqrt{p_{i}} |i^{A}\rangle |i'^{R}\rangle$$
$$= (I_{A} \otimes U_{R}) |AR_{2}\rangle$$

which proves the claim.

#### Exercise 2.82

Suppose  $\{p_i, |\psi_i\rangle\}$  is an ensemble of states generating a density matrix  $\rho = \sum_i p_i |\psi_i\rangle\langle\psi_i|$  for a quantum system A. Introduce a system R with orthonormal basis  $|i\rangle$ .

- (1) Show that  $\sum_{i} \sqrt{p_i} |\psi_i\rangle |i\rangle$  is a purification of  $\rho$ .
- (2) Suppose we measure R in the basis  $|i\rangle$ , obtained outcome i. With what probability do we obtain the result i, and what is the corresponding state of system A?
- (3) Let  $|AR\rangle$  be any purification of  $\rho$  to the system AR. Show that there exists an orthonormal basis  $|i\rangle$  in which R can be measured such that the corresponding post-measurement state for system A is  $|\psi_i\rangle$  with probability  $p_i$ .

## Solution

Concepts Involved: Linear Algebra, Purification, Schmidt Decomposition.

(1) To verify that  $\sum_i \sqrt{p_i} |\psi_i\rangle |i\rangle$  is a purification, we see that:

$$\operatorname{tr}_{R}\left(\left(\sum_{i}\sqrt{p_{i}}|\psi_{i}\rangle|i\rangle\right)\left(\sum_{j}\sqrt{p_{j}}\langle\psi_{j}|\langle j|\right)\right) = \sum_{i}\sum_{j}\sqrt{p_{i}p_{j}}|\psi_{i}\rangle\langle\psi_{j}|\operatorname{tr}_{R}(|i\rangle\langle j|)$$

$$= \sum_{i}\sum_{j}\sqrt{p_{i}p_{j}}|\psi_{i}\rangle\langle\psi_{j}|\delta_{ij}$$

$$= \sum_{i}\sqrt{p_{i}^{2}}|\psi_{i}\rangle\langle\psi_{i}|$$

$$= \sum_{i}p_{i}|\psi_{i}\rangle\langle\psi_{i}|$$

$$= \sum_{i}p_{i}|\psi_{i}\rangle\langle\psi_{i}|$$

(2) We measure the observable  $M_i = I_A \otimes \sum_i P_i = I_A \otimes \sum_i |i\rangle\langle i|$ . The probability of obtaining outcome i is given by  $p(i) = \langle AR | (I_A \otimes P_i) | AR \rangle$  (where  $|AR\rangle = \sum_i \sqrt{p_i} |\psi_i\rangle |i\rangle$ ), which we can calculate to be:

$$p(i) = \langle AR | (I_A \otimes |i\rangle\langle i|) | AR \rangle$$

$$= \left( \sum_j \sqrt{p_j} \langle \psi_j | \langle j| \right) (I_A \otimes |i\rangle\langle i|) \left( \sum_k \sqrt{p_k} |\psi_k\rangle |k\rangle \right)$$

$$= \sum_j \sum_j \sqrt{p_j} \sqrt{p_k} \langle \psi_j | \psi_k \rangle \delta_{ji} \delta_{ik}$$

$$= p_i$$

The post measurement state is given by:

$$\frac{(I_A \otimes P_i)|AR\rangle}{\sqrt{p(i)}} = \frac{(I_A \otimes |i\rangle\langle i|) \sum_j \sqrt{p_j} |\psi_j\rangle|j\rangle}{\sqrt{p_i}}$$

$$= \frac{\sum_j \sqrt{p_j} |\psi_j\rangle|j\rangle\delta_{ij}}{\sqrt{p_i}}$$

$$= \frac{\sqrt{p_i} |\psi_i\rangle|i\rangle}{\sqrt{p_i}}$$

$$= |\psi_i\rangle|i\rangle$$

so the corresponding state of system A is  $|\psi_i\rangle$ .

(3) Let  $|AR\rangle$  be any purification of  $\rho$  to the combined system AR. We then have that  $|AR\rangle$  has Schmidt Decomposition:

$$|AR\rangle = \sum_{i} \lambda_{i} |i_{A}\rangle |i_{R}\rangle$$

for orthonormal bases  $|i_A\rangle, |i_R\rangle$  of A and R respectively. Define a linear transformation U such that

 $\lambda_i |i_A\rangle = \sum_j U_{ij} p_j |\psi_j\rangle$ . We then have that:

$$|AR\rangle = \sum_{i} \left( \sum_{j} U_{ij} p_{j} |\psi_{j}\rangle \right) |i_{R}\rangle = \sum_{j} p_{j} |\psi_{j}\rangle \sum_{i} U_{ij} |i_{R}\rangle.$$

We note that we can move the  $U_{ij}$  to system R as R has the same state space as A by construction. Letting  $|j\rangle = \sum_i U_{ij} |i_R\rangle$  be our orthonormal basis of R, the claim follows (by part (2) of the question.

#### Problem 2.1: Functions of the Pauli matrices

Let  $f(\cdot)$  be any function from complex numbers to complex numbers. Let  $\mathbf n$  be a normalized vector in three dimensions, and let  $\theta$  be real. Show

$$f(\theta \mathbf{n} \cdot \boldsymbol{\sigma}) = \frac{f(\theta) + f(-\theta)}{2} I + \frac{f(\theta) - f(-\theta)}{2} \mathbf{n} \cdot \boldsymbol{\sigma}$$

#### Solution

Concepts Involved: Linear Algebra, Spectral Decomposition, Operator Functions.

From Exercise 2.35, we recall that  $\mathbf{n} \cdot \boldsymbol{\sigma}$  has spectral decomposition  $\mathbf{n} \cdot \boldsymbol{\sigma} = |n_+\rangle\langle n_+| - |n_-\rangle\langle n_-|$ . We then have that (by the definition of operator functions):

$$f(\theta \mathbf{n} \cdot \boldsymbol{\sigma}) = f\left(\theta(|n_{+}\rangle\langle n_{+}| - |n_{-}\rangle\langle n_{-}|)\right) = f(\theta)|n_{+}\rangle\langle n_{+}| + f(-\theta)|n_{-}\rangle\langle n_{-}|.$$

We then use the fact proven in the solution to Exercise 2.60 that we can write the projectors  $P_{\pm}=|n_{\pm}\rangle\langle n_{\pm}|$  in terms of the operator  $\mathbf{n}\cdot\boldsymbol{\sigma}$  as:

$$|n_{\pm}\rangle\langle n_{\pm}| = \frac{I \pm \mathbf{n} \cdot \boldsymbol{\sigma}}{2}.$$

Hence making this substitution we have:

$$f(\theta \mathbf{n} \cdot \boldsymbol{\sigma}) = f(\theta) \left( \frac{I + \mathbf{n} \cdot \boldsymbol{\sigma}}{2} \right) + f(-\theta) \left( \frac{I - \mathbf{n} \cdot \boldsymbol{\sigma}}{2} \right).$$

Grouping terms, we obtain the desired relation:

$$f(\theta \mathbf{n} \cdot \boldsymbol{\sigma}) = \frac{f(\theta) + f(-\theta)}{2} I + \frac{f(\theta) - f(-\theta)}{2} \mathbf{n} \cdot \boldsymbol{\sigma}.$$

## Remark:

Arguably, the most used application of the above identity in quantum information is when  $f(\theta \mathbf{n} \cdot \boldsymbol{\sigma}) =$ 

 $\exp\{i(\theta/2)\mathbf{n}\cdot\boldsymbol{\sigma}\}$ . In this case (as in Exercise 2.35), we have

$$\begin{split} \exp \bigl\{ i(\theta/2) \mathbf{n} \cdot \boldsymbol{\sigma} \bigr\} &= \frac{\exp \bigl\{ \theta/2 \bigr\} + \exp \bigl\{ -\theta/2 \bigr\}}{2} I + \frac{\exp \bigl\{ \theta/2 \bigr\} - \exp \bigl\{ -\theta/2 \bigr\}}{2} \mathbf{n} \cdot \boldsymbol{\sigma} \\ &= \cos \left( \frac{\theta}{2} \right) I + i \sin \left( \frac{\theta}{2} \right) \mathbf{n} \cdot \boldsymbol{\sigma} \quad . \end{split}$$

## **Problem 2.2: Properties of Schmidt numbers**

Suppose  $|\psi\rangle$  is a pure state of a composite system with components A and B.

- (1) Prove that the Schmidt number of  $|\psi\rangle$  is equal to the rank of the reduced density matrix  $\rho_A \equiv {\rm tr}_B(|\psi\rangle\langle\psi|)$ . (Note that the rank of a Hermitian operator is equal to the dimension of its support.)
- (2) Suppose  $|\psi\rangle = \sum_j |\alpha_j\rangle |\beta_j\rangle$  is a representation for  $|\psi\rangle$ , where  $|\alpha_j\rangle$  and  $|\beta_j\rangle$  are (un-normalized) states for systems A and B, respectively. Prove that the number of terms in a such a decomposition is greater than or equal to the Schmidt number of  $|\psi\rangle$ ,  $\mathrm{Sch}(\psi)$ .
- (3) Suppose  $|\psi\rangle = \alpha |\varphi\rangle + \beta |\gamma\rangle$ . Prove that

$$Sch(\psi) \ge |Sch(\varphi) - Sch(\gamma)|$$

### Solution

**Concepts Involved:** Linear Algebra, Schmidt Decomposition, Schmidt Number, Reduced Density Operators.

(1) We write the Schmidt decomposed  $|\psi\rangle$ , and therefore the density matrix  $\rho_{\psi}$  as:

$$|\psi\rangle = \sum_{i} \lambda_{i} |i_{A}\rangle |i_{B}\rangle \implies |\psi\rangle\langle\psi| = \sum_{ii'} \lambda_{i}^{2} |i_{A}\rangle\langle i_{A}| \otimes |i_{B}\rangle\langle i_{B}|$$

Taking the partial trace of subsystem B in the  $|i_B\rangle$  basis, we obtain the reduced density matrix  $\rho_A$  to be:

$$\rho_A = \operatorname{tr}_B(|\psi\rangle\langle\psi|) = \sum_i \lambda_i^2 |i_A\rangle\langle i_A|$$

 $\operatorname{Sch}(\psi)$  of the  $\lambda_i$ s are nonzero, and therefore  $\rho_A$  has  $\operatorname{Sch}(\psi)$  nonzero eigenvalues - therefore the rank of its support is  $\operatorname{Sch}(\psi)$ .

(2) Suppose for the sake of contradiction that some decomposition  $|\psi\rangle = \sum_{j=1}^{N} |\alpha_j\rangle |\beta_j\rangle$  had less terms than the Schmidt decomposition of  $|\psi\rangle$ , i.e.  $N < \mathrm{Sch}(\psi)$ .

The density matrix of  $|\psi\rangle$  is:

$$\rho_{\psi} = |\psi\rangle\langle\psi| = \sum_{j=1,k=1}^{N} |\alpha_{j}\rangle\langle\alpha_{k}| \otimes |\beta_{j}\rangle\langle\beta_{k}|$$

Tracing out subsystem B, we obtain the reduced density matrix of subsystem A:

$$\rho_A = \operatorname{Tr}_B(\rho_\psi) = \sum_{j=1,k=1}^N \left| \alpha_j \middle\langle \alpha_k \middle| \middle\langle \beta_j \middle| \beta_k \middle\rangle \right|$$

where we have used that  $\mathrm{Tr}\big(|\beta_1\rangle\langle\beta_2|\big)=\langle\beta_1|\beta_2\rangle$ . From the above, it is clear that  $\rho_A$  has rank at most N, as the support of  $\rho_A$  is spanned by  $\{|\alpha_1\rangle,\ldots,|\alpha_N\rangle\}$ . But then the rank of  $\rho_A$  is less than  $\mathrm{Sch}(\psi)$ , which contradicts our finding in part (a).

(3) If  $Sch(\varphi) = Sch(\gamma)$  then there is nothing to prove as  $Sch(\psi)$  is non-negative by definition. Suppose then that  $Sch(\varphi) \neq Sch(\gamma)$ . WLOG suppose  $Sch(\varphi) > Sch(\gamma)$ . We can then write:

$$|\varphi\rangle = \frac{\beta}{\alpha} |\gamma\rangle - \frac{1}{\alpha} |\psi\rangle$$

If we Schmidt decompose  $|\varphi\rangle$  and  $|\psi\rangle$ , we have written  $|\varphi\rangle$  as the sum of  $\mathrm{Sch}(\gamma) + \mathrm{Sch}(\psi)$  (unnormalized) bipartite states. Applying the result from part (2) of this problem, we then have that:

$$Sch(\varphi) \leq Sch(\gamma) + Sch(\psi)$$

which we rearrange to obtain:

$$Sch(\psi) \ge Sch(\varphi) - Sch(\gamma) = |Sch(\varphi) - Sch(\gamma)|$$

which proves the claim.

# Problem 2.3: Tsirelson's inequality

Suppose  $Q = \mathbf{q} \cdot \boldsymbol{\sigma}$ ,  $R = \mathbf{r} \cdot \boldsymbol{\sigma}$ ,  $S = \mathbf{s} \cdot \boldsymbol{\sigma}$ ,  $T = \mathbf{t} \cdot \boldsymbol{\sigma}$ , where  $\mathbf{q}, \mathbf{r}, \mathbf{s}$ , and  $\mathbf{t}$  are real unit vectors in three dimensions. Show that

$$(Q \otimes S + R \otimes S + R \otimes T - Q \otimes T)^2 = 4I + [Q, R] \otimes [S, T]$$

Use this result to prove that

$$\langle Q \otimes S \rangle + \langle R \otimes S \rangle + \langle R \otimes T \rangle - \langle Q \otimes T \rangle < 2\sqrt{2}$$

so the violation of the Bell inequality found in Equation (2.230) is the maximum possible in quantum mechanics.

#### Solution

Concepts Involved: Tensor Products, Commutators

We first show that  $N^2 = I$  for any  $N = \mathbf{n} \cdot \boldsymbol{\sigma}$  where  $\mathbf{n}$  is a unit vector in three dimensions. We have that:

$$N^{2} = \left(\sum_{i=1}^{3} n_{i} \sigma_{i}\right)^{2}$$

$$= n_{1}^{2} \sigma_{1}^{2} + n_{2}^{2} \sigma_{2}^{2} + n_{3}^{2} \sigma_{3}^{2} + n_{1} n_{2} (\sigma_{1} \sigma_{2} + \sigma_{2} \sigma_{1}) + n_{1} n_{3} (\sigma_{1} \sigma_{3} + \sigma_{3} \sigma_{1}) + n_{2} n_{3} (\sigma_{2} \sigma_{3} + \sigma_{3} \sigma_{2})$$

By Exercise 2.41,  $\sigma_i^2 = I$  and  $\{\sigma_i, \sigma_i\} = 0$  for  $i \neq j$ , so the above reduces to:

$$N^{2} = n_{1}^{2}I + n_{2}^{2}I + n_{3}^{2}I = (n_{1}^{2} + n_{2}^{2} + n_{3}^{2})I = I$$

where we use the fact that n is of unit length. Using this fact, we have that:

$$(Q \otimes S + R \otimes S + R \otimes T - Q \otimes T)^2 = Q^2 \otimes S^2 + QR \otimes S^2 + QR \otimes ST - Q^2 \otimes ST + RQ \otimes S^2 + R^2 \otimes S^2 + R^2 \otimes ST - RQ \otimes ST + RQ \otimes TS + R^2 \otimes TS + R^2 \otimes T^2 - RQ \otimes T^2 - Q^2 \otimes TS - QR \otimes TS - QR \otimes T^2 + Q^2 \otimes T^2 = I \otimes I + QR \otimes I + QR \otimes ST - I \otimes ST + RQ \otimes I + I \otimes I + I \otimes ST - RQ \otimes ST + RQ \otimes TS + I \otimes TS + I \otimes I - RQ \otimes I - I \otimes TS - QR \otimes TS - QR \otimes I + I \otimes I = 4I \otimes I + RQ \otimes TS - RQ \otimes ST + QR \otimes ST - QR \otimes TS -$$

which proves the first equation. We have that  $\langle 4I \rangle = 4$   $\langle I \rangle = 4$ . Since each of Q,R,S,T have eigenvalues  $\pm 1$  (Exercise 2.35), we also ave that  $\langle [Q,R] \otimes [S,T] \rangle \leq 4$  as the tensor product of commutators consists of 4 terms, each of which has expectation less than or equal to 1. We therefore have by the linearity of expectation (Exercise A1.4) that:

$$\langle (Q \otimes S + R \otimes S + R \otimes T - Q \otimes T)^2 \rangle = \langle 4I + [QR] \otimes [S, T] \rangle \leq 8.$$

Furthermore, we have that:

$$\left\langle \left(Q\otimes S+R\otimes S+R\otimes T-Q\otimes T\right)\right\rangle ^{2}\leq \left\langle \left(Q\otimes S+R\otimes S+R\otimes T-Q\otimes T\right)^{2}\right\rangle$$

so combining the two inequalities we obtain:

$$\langle (Q \otimes S + R \otimes S + R \otimes T - Q \otimes T) \rangle^2 \le 8.$$

Taking square roots on both sides, we have:

$$\langle (Q \otimes S + R \otimes S + R \otimes T - Q \otimes T) \rangle \le 2\sqrt{2}$$

and again by the linearity of expectation:

$$\langle Q \otimes S \rangle + \langle R \otimes S \rangle + \langle R \otimes T \rangle - \langle Q \otimes T \rangle \le 2\sqrt{2}$$

which is the desired inequality.

# 3 Introduction to computer science

#### Exercise 3.1: Non-computable processes in Nature

How might we recognize that a process in Nature computes a function not computable by a Turing machine?

#### Solution

Concepts Involved: Turing Machines.

One criteria is natural phenomena that appear to be truly random; Turing machines as defined in the text are deterministic (though there are probabilistic variations that would solve this issue) and hence would not be able to compute a random function. From a more direct point, if a process in Nature was to be found to compute a known non-computable problem (e.g. solve the Halting problem or the Tiling problem) then we may conclude (trivially) that the process would not be computable. However since the domain of inputs that we could provide top such a natural process would have to be finite, there would be no concrete method in which one could actually test if such a process was truly computing a non-Turing computable function (as a Turing machine that works on a finite subset of inputs for an uncomputable problem could be devised).

## **Exercise 3.2: Turing numbers**

(\*) Show that single-tape Turing machines can each be given a number from the list  $1,2,3,\ldots$  in such a way that the number uniquely specifies the corresponding machine. We call this number the *Turing number* of the corresponding Turing machine. (Hint: Every positive integer has a unique prime factorization  $p_1^{a_1}p_2^{a_2}\ldots p_k^{a_k}$ , where  $p_i$  are distinct prime numbers, and  $a_1,\ldots a_k$  are non-negative integers.)

### Solution

**Concepts Involved:** Turing Machines, Cardinality.

Lemma 1. If  $\{A_n\}_{n=1}^{\infty}$  is a sequence of sets that are countably infinite (that is, they can be put in bijection with the natural numbers  $\mathbb{N}$ ) then their union  $A = \bigcup_{n=1}^{\infty} A_n$  is also countably infinite.

*Proof.* Write  $A_n = \{x_{n1}, x_{n2}, x_{n3}, \ldots\}$  (which we can do as each of the  $A_n$ s are countably infinite). Then, we form an array:

$$A_1 = x_{11} x_{12} x_{13} \cdots$$
 $A_2 = x_{21} x_{22} x_{23} \cdots$ 
 $A_3 = x_{31} x_{32} x_{33} \cdots$ 

Then, we can re-number the elements along the diagonal lines (i.e.  $x_{11}, x_{21}, x_{12}, x_{31}, x_{22}, x_{13}, \ldots$ ). This new enumeration corresponds to a countably infinite set. From there, we let  $T \subset \mathbb{N}$  be the remaining labels in the enumeration after removing the repeated elements from the sequence. Then, |T| = |A|, and hence A is at most countably infinite. A cannot be finite as  $A_1 \subset A$  and  $A_1$  is not finite. Hence A is countably infinite.

Lemma 2. The set of all finite sequences of elements from a countably infinite set B is also countably infinite.

*Proof.* Denote  $B_n$  the set of length n sequences consisting of elements from B. We show that  $B_n$  is countably infinite by induction.

Base Case:  $|B| = |B_1|$ , so  $B_1$  is countably infinite.

Inductive Step: Suppose that  $B_k$  is countably infinite for  $k \geq 1$ . Then, define  $B_{k,i}$  which is the set of k+1 length sequences, consisting of all sequences in  $B_k$  and then terminating with the ith element of B.  $\left|B_{k,i}\right| = \left|B_k\right|$  so  $B_{k,i}$  is countably infinite, and then  $B_{k+1} = \bigcup_{i=1}^{\infty} B_{k,i}$  is countably infinite by Lemma 1. This concludes the argument by induction.

Finally, the set of all finite sequences of elements from B is simply  $\bigcup_{k=1}^{\infty} B_k$ , which is a union of countably infinite sets and therefore countably infinite again by Lemma '.

Solution. A single-tape Turing machine can be uniquely specified by a finite integer sequence as follows. First, we state the number of states |Q|=N, then enumerate the states  $q_s=-1,q_h=0,q_1=1,\ldots,q_{N-2}=N-2$ . We then state the size of the alphabet  $|\Gamma|=L$  and enumerate  $|\Gamma|=L$  and enumer

As an example, the Turing machine in the text which computes the constant function (and let us say for the sake of example the tape starts with the number 7 on it, so 110 in binary followed by infinite blanks) has the associated finite integer sequence:

$$\underbrace{\frac{5}{|Q|}, \underbrace{-1}_{q_s}, \underbrace{0}_{q_h}, \underbrace{1}_{q_1}, \underbrace{2}_{q_2}, \underbrace{3}_{q_3}, -3, \underbrace{4}_{|\Sigma|}, \underbrace{-2}_{\triangleright}, \underbrace{-1}_{b}, 0, 1, -3, \underbrace{1, 1, 0}_{7}, -3, \underbrace{-1, -2, 1, -2, 1}_{\langle q_s, \triangleright, q_1, \triangleright, +1 \rangle},}_{\langle q_1, 0, q_1, b, +1 \rangle}, \underbrace{\frac{1, 0, 1, -1}{\langle q_1, 0, q_2, b, -1 \rangle}, \underbrace{2, -1, 2, -1, -1}_{\langle q_2, b, q_2, b, -1 \rangle}, \underbrace{2, -2, 3, -2, 1}_{\langle q_2, b, q_2, b, -1 \rangle}, \underbrace{3, -1, 0, 1, 0}_{\langle q_3, b, q_h, 1, 0 \rangle}$$

Because the number of states is finite, the size of the alphabet is finite, the number of non-blank squares on the initial tape is finite, and there is a finite number of lines in the program, the above algorithm produces a finite sequence of integers for a given Turing machine.

Now, observe that the integers are countably infinite, as:

$$f : \mathbb{N} \longrightarrow \mathbb{Z}$$

$$n \longmapsto \begin{cases} \frac{n}{2} & n \text{ even} \\ -\frac{n-1}{2} & n \text{ odd} \end{cases}$$

$$(1)$$

is a bijection. So too then is the set of all finite integer sequences by the Lemma 2. Since the finite integer sequences describing Turing machines as produced by the above prescription are a subset of all finite integer sequences, we must therefore be able to associate any single-tape Turing machine to a unique natural number, as claimed.

#### Exercise 3.3: Turing machine to reverse a bit string

Describe a Turing machine which takes a binary number x as input, and outputs the bits of x in reverse order. (*Hint*: In this exercise and the next it may help to use a multi-tape Turing machine and/or symbols other than  $\triangleright$ , 0, 1 and the blank.)

### Exercise 3.4: Turing machine to add modulo 2

Describe a Turing machine to add two binary numbers x and y modulo 2. The numbers are input on the Turing machine tape in binary, in the form x, followed by a single blank, followed by a y. If one number is not as long as the other then you may assume that it has been padded with leading 0s to make the two numbers the same length.

#### **Exercise 3.5: Halting problem with no inputs**

Show that given a Turing machine M there is no algorithm to determine whether M halts when the input to the machine is a blank tape.

#### **Exercise 3.6: Probabilistic halting problem**

Suppose we number the probabilistic Turing machines using a scheme similar to that found in Exercise 3.2 and define the probabilistic halting function  $h_p(x)$  to be 1 if machine x halts on input of x with probability at least 1/2 and 0 if machine x halts on input of x with probability less than 1/2. Show that there is no probabilistic Turing machine which can output  $h_p(x)$  with probability of correctness strictly greater than 1/2 for all x.

#### **Exercise 3.7: Halting oracle**

Suppose a black box is made available to us which takes a non-negative integer x as input, and then outputs the value of h(x), where  $h(\cdot)$  is the halting function defined in Box 3.2 on page 130. This type of black box is sometimes known as an oracle for the halting problem. Suppose we have a regular Turing machine which is augmented by the power to call the oracle. One way of accomplishing this is to use a two-tape Turing machine, and add an extra program instruction to the Turing machine which results in the oracle being called, and the value of h(x) being printed on the second tape, where x is the current contents of the second tape. It is clear that this model for computation is more powerful than the conventional Turing machine model, since it can be used to compute the halting function. Is the halting problem for this model of computation undecidable? That is, can a Turing machine aided by an oracle for the halting problem decide whether a program for the Turing machine with oracle will halt on a particular input?

#### Exercise 3.8: Universality of NAND

Show that the NAND gate can be used to simulate the AND, XOR, and NOT gates, provides wires, ancilla bits and FANOUT are available.

#### Solution

#### Concepts Involved: Logic Gates.

We start by showing how we can get a 1 qubit using two 0 ancilla bits and a NAND gate.

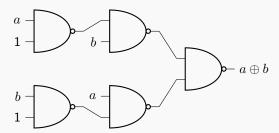
We will now show how to simulate the NOT, AND, and XOR gates. We note that we will use "1" to denote as shorthand a 1 bit constructed using two ancilla bits (as above). a/b represent the input bits. We start with the NOT gate.

$$a - 1 - 1 - 1$$

Next, we simulate the AND gate.

$$a - b - a \wedge b$$

For the XOR simulated gate, we note that we first use FANOUT twice to copy both input bits.



Having simulated the three gates using the NAND gate only, we conclude that the NAND is universal.

## Exercise 3.9

Prove that f(n) is O(g(n)) if and only if g(n) is  $\Omega(f(n))$ . Deduce that f(n) is  $\Theta(g(n))$  if and only if g(n) is  $\Theta(f(n))$ .

#### Solution

Concepts Involved: Asymptotic Notation.

Suppose f(n) is O(g(n)). Then, there exists c>0 such that for all  $n>n_0$ ,  $f(n)\leq cg(n)$ . Therefore, we have  $\frac{1}{c}>0$  such that for all  $n>n_0$ ,  $\frac{1}{c}f(n)\leq g(n)$ . Hence, g(n) is  $\Omega(f(n))$ . Conversely, if g(n) is  $\Omega(f(n))$ , there exists c>0 such that for all  $n>n_0$ ,  $cf(n)\leq g(n)$ . Hence, we have  $\frac{1}{c}>0$  such that for all  $n>n_0$ ,  $f(n)\leq \frac{1}{c}g(n)$  and hence f(n) is O(g(n)).

Therefore, if f(n) is  $\Theta(g(n))$  then f(n) is O(g(n)) and  $\Omega(g(n))$ , and by the above argument, g(n) is O(f(n)) and  $\Omega(f(n))$  and hence g(n) is  $\Theta(f(n))$ . The converse holds in the same way.

#### Exercise 3.10

Suppose g(n) is a polynomial of degree k. Show that g(n) is  $O(n^l)$  for any  $l \geq k$ .

#### Solution

Concepts Involved: Asymptotic Notation.

By assumption,  $g(n)=a_0+a_1n^1+a_2n^2+\ldots+a_kn^k$  with  $a_k\neq 0$ . For  $n\geq 1$  we have that  $n^l\geq n^k$  if  $l\geq k$ , and hence if  $l\geq k$  we have that  $a_in^l\geq a_in^i$  for all  $i\in\{0,\ldots,k\}$ . Therefore, we have that:

$$(a_0 + a_1 + \ldots + a_k)n^l \ge a_0 + a_1n^1 + \ldots + a_kn^k = g(n)$$

for  $n \ge 1$  and hence g(n) is  $O(n^l)$ .

#### Exercise 3.11

Show that  $\log n$  is  $O(n^k)$  for any k > 0.

#### Solution

Concepts Involved: Asymptotic Notation.

Let k > 0 and c > 0. By the definition of the exponential we have that:

$$\exp(cn^k) = \sum_{j=0}^{\infty} \frac{(an^k)^j}{j!} = \sum_{j=0}^{\infty} \frac{a^{kj} n^{kj}}{j!}$$

now, there exists some  $j_0 \in \mathbb{Z}$  for which  $kj_0 > 1$ . Since for  $n \geq 0$  the terms in the above sum are non-negative, we find:

$$\exp\!\left(cn^k\right) \ge \frac{c^{kj_0}n^{kj_0}}{j_0!}$$

Now, choose c sufficiently large such that  $c^{kj_0} \ge j_0!$ . We then find that:

$$\exp(cn^k) \ge \frac{c^{kj_0}n^{kj_0}}{j_0!} \ge n^{kj_0}$$

Then for  $n \ge 1$  it follows that  $n^{kj_0} \ge n$  as  $kj_0 \ge 1$  and so:

$$\exp(cn^k) \ge n$$

Since the logarithm is monotonic, we may take the log of both sides and preserve the inequality:

$$cn^k > \log n$$

So we have shown that for any k > 0, there exists c > 0 such that for all n > 1,  $cn^k \ge \log n$ . Hence,  $\log n$  is  $O(n^k)$  for any k > 0.

#### **Exercise 3.12:** $n^{\log n}$ is super-polynomial

Show that  $n^k$  is  $O(n^{\log n})$  for any k, but that  $n^{\log n}$  is never  $O(n^k)$ .

#### Solution

Concepts Involved: Asymptotic Notation.

First, note that for any k,  $e^k \le n$  for sufficiently large  $n > n_0$  and so  $k \le \log n$  by monotonicity of the logarithm. Therefore, for  $n > n_0$  it follows by monotonicity (of exponentiation) that  $n^k \le n^{\log n}$  and so  $n^k$  is  $O(n^{\log n})$ .

Now, consider an arbitrary a>0. It still follows for sufficiently large  $n>n_0$  that  $e^{ak}\leq n$  and so  $ak\leq \log n$  and  $n^an^k\leq n^{\log n}$ . But for any c>0  $n^a>c$  for sufficiently large n and so:

$$cn^k < n^{\log n}$$

So since for any c>0 there exists some  $n_0'$  for which  $n>n_0'$  implies  $cn^k\leq n^{\log n}$ , it follows that  $n^{\log n}$  is never  $O(n^k)$ .

## **Exercise 3.13:** $n^{\log n}$ is sub-exponential

Show that  $c^n$  is  $\Omega(n^{\log n})$  for any c>1, but that  $n^{\log n}$  is never  $\Omega(c^n)$ .

#### Solution

Concepts Involved: Asymptotic Notation.

First note from Exercise 1.11 that  $\log n$  is  $O(n^k)$  for any k > 0. Specifically, take k = 1/2; then there exists a > 0 such that for  $n > n_0$ :

$$an^{1/2} > \log n$$

and therefore squaring both sides:

$$a^2 n > \log n \log n = \log n^{\log n}$$

Now for any c>1, we can define  $a'=\frac{a^2}{\log c}>0$  and write:

$$na'\log c = \log c^{na'} \ge \log n^{\log n}$$

Exponentiating both sides preserves the inequality, and so:

$$c^{na'} = c^{a'}c^n > n^{\log n}$$

and so there exists a constant  $\frac{1}{c^{a'}}>0$  such that for  $n>n_0$ ,  $c^n\geq \frac{1}{c^{a'}}n^{\log n}$  and therefore  $c^n$  is  $\Omega(n^{\log n})$ . Now, let b>0 be some arbitrarily small constant. For sufficiently large n, we have that  $na'logc+\log b>0$ 

and so for sufficiently large n it further follows that:

$$\log b + na' \log c = \log \left( bc^{na'} \right) \ge \log n^{\log n}$$

where a' is defined as it was previously. Therefore exponentiating both sides:

$$bc^{na'} = bc^{a'}c^n > n^{\log n}$$

so for sufficiently large n, for any arbitrarily small constant b' it follows that  $b'c^n \ge n^{\log n}$  and so  $n^{\log n}$  is never  $\Omega(c^n)$ .

#### Exercise 3.14

Suppose e(n) is O(f(n)) and g(n) is O(h(n)). Show that e(n)g(n) is O(f(n)h(n)).

#### Solution

Concepts Involved: Asymptotic Notation.

By assumption, we have that  $e(n) \le c_1 f(n)$  for some  $c_1 > 0$  and for all  $n > n_1$  and that  $g(n) \le c_2 h(n)$  for some  $c_2 > 0$  and for all  $n > n_2$ . Let  $n_0 = \max n_1, n_2$ . We then have that for  $n > n_0$  that:

$$e(n)g(n) \le c_1 f(n)c_2 h(n) = (c_1 c_2)(f(n)h(n))$$

so therefore e(n)g(n) is O(f(n)h(n)).

#### Exercise 3.15: Lower bound for compare-and-swap based sorts

Suppose an n element list is sorted by applying some sequence of compare-and-swap operations to the list. There are n! possible initial orderings of the list. Show that after k of the compare-and-swap operations have been applied, at most  $2^k$  of the possible initial orderings will have been sorted into the correct order. Conclude that  $\Omega(n \log n)$  compare and swap operations are required to sort all possible initial orderings into the correct order.

#### Solution

Concepts Involved: Asymptotic Notation, Compare-and-Swap.

We prove the first statement by induction. After 0 steps, we have that  $1=2^0$  out of the n! possible orderings are already sorted. Let  $k\in\mathbb{N}, k\geq 0$  and suppose that after k swaps, at most  $2^k$  of the initial orderings have been sorted into the correct order. We now consider the state of the list after the k+1th swap. Each of the  $2^k$  initial orderings from the previous step are correctly sorted already (so the swap does nothing), and there are a further  $2^k$  initial orderings that are one swap away from the  $2^k$  from the previous step, and hence the k+1th swap will put  $2^k$  more initial orderings into the correct order. Therefore, after  $2^{k+1}$  compare and swaps, there are at most  $2^k+2^k=2^{k+1}$  possible initial orderings that are sorted into the correct order. This proves the claim.

Using the above fact, we have that in order to have all n! possible initial orderings correct after k

steps that  $2^k \ge n!$ . Taking logarithms on both sides, we have that  $\log(2^k) \ge \log(n!)$  and hence  $k \ge \log(n!)$ . Using Stirling's approximation for factorials (https://en.wikipedia.org/wiki/Stirling% 27s\_approximation), we have that:

$$k \ge n \log n - n \log e + O(\log n)$$

from which we conclude that k is  $\Omega(n \log n)$  and hence  $\Omega(n \log n)$  compare and swap operations are required to sort all possible initial orderings into the correct order.

#### Exercise 3.16: Hard-to-compute functions exist

Show there exist Boolean functions on n inputs which require at least  $2^n/\log n$  logic gates to compute.

#### Exercise 3.17

Prove that a polynomial-time algorithm for finding the factors of a number m exists if and ony if the factoring decision problem is in  $\mathbf{P}$ .

#### Exercise 3.18

Prove that if  $coNP \neq NP$  then  $P \neq NP$ .

#### Exercise 3.19

The Reachability problem is to determine whether there is a path between two specified vertices in a graph. Show that Reachability can be solved using O(n) operations if the graph has n vertices. Use the solution to Reachability to show that it is possible to decide whether a graph is connected in  $O(n^2)$  operations.

#### Exercise 3.20: Euler's theorem

Prove Euler's theorem. In particular, if each vertex has an even number of incident edges, give a constructive procedure for finding an Euler cycle.

#### **Exercise 3.21: Transitive property of reduction**

Show that if a language  $L_1$  is reducible to the language  $L_2$  and the language  $L_2$  is reducible to  $L_3$  then the language  $L_1$  is reducible to the language  $L_3$ .

#### Exercise 3.22

Suppose L is complete for a complexity class, and L' is another language in the complexity class such that L reduces to L'. Show that L' is complete for the complexity class.

## Exercise 3.23

Show that SAT is  $\mathbf{NP}$ -complete by first showing that the SAT is in  $\mathbf{NP}$ , and then showing that CSAT reduces to SAT.

#### Exercise 3.24: 2SAT has an efficient solution

Suppose  $\varphi$  is a Boolean formula in conjunctive normal form, in which each clause contains only two literals.

- (1) Construct a (directed) graph  $G(\varphi)$  with directed edges in the following way: the vertices of G correspond to variables  $x_k$  and their negations  $\neg x_j$  in  $\varphi$ . There is a (directed) edge  $(\alpha,\beta)$  in G if and only if the clause  $(\neg \alpha \lor \beta)$  or the clause  $(\beta \land \neg \alpha)$  is present in  $\varphi$ . Show that  $\varphi$  is not satisfiable if and only if there exists a variable x such that there are paths from x and x and from x to x in  $G(\varphi)$ .
- (2) Show that given a directed graph G containing n vertices it is possible to determine whether two vertices  $v_1$  and  $v_2$  are connected in polynomial time.
- (3) Find an efficient algorithm to solve 2SAT.

#### Exercise 3.25: $PSPACE \subseteq EXP$

The complexity class  $\mathbf{EXP}$  (for exponential time) contains all decision problems which may be decided by a Turing machine running in exponential time, that is time  $O(2^{n^k})$ , where k is any constant. Prove that  $\mathbf{PSPACE} \subseteq \mathbf{EXP}$ . (Hint: If a Turing machine has l internal states, an m letter alphabet, and uses space p(n), argue that the machine can exist in one of at most  $lm^{p(n)}$  different states, and that if the Turing machine is to avoid infinite loops then it must halt before revisiting a state.)

#### Exercise 3.26: $L \subseteq P$

The complexity class  $\mathbf L$  (for *logarithmic space*) contains all decision problems which may be decided by a Turing machine running in logarithmic space, that is, in space  $O(\log(n))$ . More precisely, the class  $\mathbf L$  is defined using a two-tape Turing machine. The first tape contains the problem instance, of size n, and is a read-only tape, in the sense that only program lines which don't change the contents of the first tape are allowed. The second tape is a working tape which initially contains only blanks. The logarithmic space requirement is imposed on the second, working tape only. Show that  $\mathbf L \subseteq \mathbf P$ .

## Exercise 3.27: Approximation algorithm for VERTEX COVER

Let G=(V,E) be an undirected graph. Prove that the following algorithm finds a vertex cover for G that is within a factor of two of being a minimial vertex cover.

```
\begin{array}{l} VC = \emptyset \\ E' = E \\ \text{while } E' = \emptyset \text{ do} \\ & | \text{ let } (\alpha,\beta) \text{ be any edge of } E' \\ & | VC = VC \cup \{\alpha,\beta\} \\ & | \text{ remove from } E' \text{ every edge incident on } \alpha \text{ or } \beta \\ \text{end} \\ \text{return } VC \end{array}
```

#### Exercise 3.28: Arbitrariness of the constant in the definition of BPP

Suppose k is a fixed constant,  $1/2 < k \le 1$ . Suppose L is a language such that there exists a Turing machine M with the property that whenever  $x \in L$ , M accepts x with probability at least k, and whenever  $x \notin L$ , M rejects x with probability at least k. Show that  $L \in \mathbf{BPP}$ .

#### Exercise 3.29: Fredkin gate is self-inverse

Show that applying two consecutive Fredkin gates gives the same outputs as inputs.

#### Solution

Concepts Involved: Fredkin Gates. Recall the input/output table of the Fredkin gate:

Inputs			Outputs		
a	b	c	a'	b'	c'
0	0	0	0	0	0
0	0	1	0	0	1
0	1	0	0	1	0
0	1	1	1	0	1
1	0	0	1	0	0
1	0	1	0	1	1
1	1	0	1	1	0
_1	1	1	1	1	1

We check for all possible 8 input states that applying the Fredkin gate returns the original input state.

$$F[F[(0,0,0)]] = F[(0,0,0)] = (0,0,0)$$

$$F[F[(0,0,1]] = F[(0,0,1)] = (0,0,1)$$

$$F[F[(0,1,0]] = F[(0,1,0)] = (0,1,0)$$

$$F[F[(0,1,1)]] = F[(1,0,1)] = (0,1,1)$$

$$F[F[(1,0,0)]] = F[(1,0,0)] = (1,0,0)$$

$$F[F[(1,0,1)]] = F[(0,1,1)] = (1,0,1)$$

$$F[F[(1,1,0)]] = F[(1,1,0)] = (1,1,0)$$

$$F[F[(1,1,1)]] = F[(1,1,1)] = (1,1,1)$$

We conclude that the Fredkin gate is self-inverse.

## Exercise 3.30

Verify that the billiard ball computer in Figure 3.14 computes the Fredkin gate.

#### Exercise 3.31: Reversible half-adder

Construct a reversible circuit which, when two bits x and y are input, outputs  $(x, y, c, x \oplus y)$ , where c is the carry bit when x and y are odd.

#### Exercise 3.32: From Fredkin to Toffoli and back again

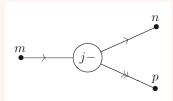
What is the smallest number of Fredkin gates needed to simulate a Toffoli gate? What is the smallest number of Toffoli gates needed to simulate a Fredkin gate?

## **Problem 3.1: Minsky machines**

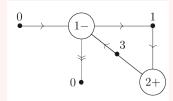
A *Minsky machine* consists of a finite set of *registers*,  $r_1, r_2, \ldots, r_k$ , each capable of holding an arbitrary non-negative integer, and a *program*, made up of *orders* of one of two types. The first type has the form:



The interpretation is that at point m in the program register rj is incremented by one, and execution proceeds to point n in the program. The second type of order has the form:



The interpretation is that at point m in the program, register  $r_j$  is decremented if it contains a positive integer, and execution proceeds to point n in the program. If register  $r_j$  is zero then execution simply proceeds to point p in the program. The *program* for the Minsky machine consists of a collection of such orders, of a form like:



The starting and all possible halting points for the program are conventionally labeled zero. This program takes the contents of register  $r_1$  and adds them to register  $r_2$ , while decrementing  $r_1$  to zero.

- (1) Prove that all (Turing) computable functions can be computed on a Minsky machine, in the sense that given a computable function  $f(\cdot)$  there is a Minsky machine program that when the registers start in the state  $(n,0,\ldots,0)$  gives as output  $(f(n),0,\ldots,0)$ .
- (2) Sketch a proof that any function which can be computed on a Minsky machine, in the sense just defined, can also be computed on a Turing machine.

#### **Problem 3.2: Vector games**

A vector game is specified by a finite list of vectors, all of the same dimension, and with integer coordinates. The game is to start with a vector x of non-negative integer co-ordinates and to add to x the first vector from the list which preserves the non-negativity of all the components, and to repeat this process until it is no longer possible. Prove that for any computable function  $f(\cdot)$  there is a vector game which when started with the vector  $(n,0,\ldots,0)$  reaches  $(f(n),0,\ldots,0)$  (Hint: Show that a vector game in k+2 dimensions can simulate a Minsky machine containing k registers.)

## Problem 3.3: Fractran

A Fractran program is defined by a list of positive rational numbers  $q_1,\ldots,q_n$ . It acts on a positive integer m by replacing it by  $q_im$  where i is the least number such that  $q_im$  is an integer. If there is ever a time when there is no i such that  $q_im$  is an integer, then execution stops. Prove that for any computable function  $f(\cdot)$  there is a Fractran program which when started with  $2^n$  reaches  $2^{f(n)}$  without going through any intermediate powers of 2. (Hint: use the previous problem.)

#### Problem 3.4: Undecidability of dynamical systems

A Fractran program is essentially just a very simple dynamical system taking positive integers to positive integers. Prove that there is no algorithm to decide whether such a dynamical system ever reaches 1.

#### Problem 3.5: Non-universality of two bit reversible logic

Suppose we are trying to build circuits using only one and two bit reversible logic gates, and ancilla bits. Prove that there are Boolean functions which cannot be computed in this fashion. Deduce that the Toffoli gate cannot be simulated using one and two bit reversible gates, even with the aid of ancilla bits.

#### Problem 3.6: Hardness of approximation of TSP

Let  $r\geq 1$  and suppose that there is an approximation algorithm for TSP which is guaranteed to find the shorted tour among n cities to within a factor r. Let G=(V,E) be any graph on n vertices. Define an instance of TSP by identifying cities with vertices in V, and deefinign the distance between cities i and j to be 1 if (i,j) is an edge of G, and to be  $\lceil r \rceil |V| + 1$  otherwise. Show that if the approximation algorithm is applied to this instance of TSP then it returns a Hamiltonian cycle for G if one exists, and otherwise returns a tour of length more than  $\lceil r \rceil |V|$ . From the  $\mathbf{NP}$ -completeness of HC it follows that no such approximation algorithm can exist unless  $\mathbf{P} = \mathbf{NP}$ .

#### **Problem 3.7: Reversible Turing machines**

- (1) Explain how to construct a reversible Turing machine that can compute the same class of functions as is computable on an ordinary Turing machine. (*Hint:* It may be helpful to use a multi-tape construction.)
- (2) Give general space and time bounds for the operation of your reversible Turing machine, in terms of the time t(x) and space s(x) required on an ordinary single-tape Turing machine to compute a function f(x).

## Problem 3.8: Find a hard-to-compute class of functions (Research)

Find a natural class of functions on n inputs which requires a super-polynomial number of Boolean gates to compute.

#### Problem 3.9: Reversible PSPACE = PSPACE

It can be shown that the problem 'quantified satisfiability', or QSAT, is **PSPACE**-complete. That is, every other language in **PSPACE** can be reduced to QSAT in polynomial time. The language QSAT is defined to consist of all Boolean formulae  $\varphi$  in n variables  $x_1, \ldots, x_n$ , and in conjunctive normal form, such that:

$$\exists_{x_1} \forall_{x_2} \exists_{x_3} \dots \forall_{x_n} \varphi$$
 if  $n$  is even;  $\exists_{x_1} \forall_{x_2} \exists_{x_3} \dots \exists_{x_n} \varphi$  if  $n$  is odd.

Prove that a reversible Turing machine operating in polynomial space can be used to solve QSAT. Thus, the class of languages decidable by a computer operating reversibly in polynomial space is equal to **PSPACE**.

#### Problem 3.10: Ancilla bits and efficiency of reversible computation

Let  $p_m$  be the mth prime number. Outline the construction of a reversible circuit which, upon the input of m and n such that n>m, outputs the product  $p_mp_n$ , that is  $(m,n)\mapsto (p_mp_n,g(m,n))$  where g(m,n) is the final state of the ancilla bits used by the circuit. Estimate the number of ancilla qubits your circuit requires. Prove that if a polynomial (in  $\log n$ ) size reversible circuit can be found that uses  $O(\log(\log n))$  ancilla bits then the problem of factoring a product of two prime numbers is in  $\mathbf{P}$ .

# 4 Quantum circuits

#### Exercise 4.1

In Exercise 2.11, which you should do now if you haven't already done it, you computed the eigenvectors of the Pauli matrices. Find the points on the Bloch sphere which correspond to the normalized eigenvectors of the different Pauli matrices.

#### Solution

Concepts Involved: Linear Algebra.

Recall that a single qubit in the state  $|\psi\rangle=a|0\rangle+b|1\rangle$  can be visualized as a point  $(\theta,\varphi)$  on the Bloch sphere, where  $a=\cos(\theta/2)$  and  $b=e^{i\varphi}\sin(\theta/2)$ .

We recall from 2.11 that Z (and I) has eigenvectors  $|0\rangle, |1\rangle$ , X has eigenvectors  $|+\rangle = \frac{|0\rangle + |1\rangle}{\sqrt{2}}, |-\rangle = \frac{|0\rangle - |1\rangle}{\sqrt{2}}$ , and Y has eigenvectors  $|y_+\rangle = \frac{|0\rangle + i|1\rangle}{\sqrt{2}}, |y_-\rangle = \frac{|0\rangle - i|1\rangle}{\sqrt{2}}$ . Expressing these vectors as points on the Bloch sphere (using spherical coordinates), we have:

$$|0\rangle \cong (0,0); |1\rangle \cong (\pi,0); |+\rangle \cong \left(\frac{\pi}{2},0\right);$$
$$|-\rangle \cong \left(\frac{\pi}{2},\pi\right); |y_{+}\rangle \cong \left(\frac{\pi}{2},\frac{\pi}{2}\right); |y_{-}\rangle \cong \left(\frac{\pi}{2},\frac{3\pi}{2}\right).$$

Exercise 4.2

Let x be a real number and A a matrix such that  $A^2 = I$ . Show that

$$\exp(iAx) = \cos(x)I + i\sin(x)A$$

Use this result to verify Equations (4.4) through (4.6).

#### Solution

Concepts Involved: Linear Algebra, Operator Functions.

Let  $|v\rangle$  be an eigenvector of A with eigenvalue  $\lambda$ . It then follows that  $A^2|v\rangle=\lambda^2|v\rangle$ , and furthermore we have that  $A^2|v\rangle=I|v\rangle=|v\rangle$  by assumption. We obtain that  $\lambda^2=1$  and therefore the only possible eigenvalues of A are  $\lambda=\pm 1$ . Let  $|v_1\rangle,\ldots,|v_k\rangle$  be the eigenvectors with eigenvalue 1 and  $|v_{k+1}\rangle,\ldots,|v_n\rangle$  be the eigenvectors with eigenvalue -1. By the spectral decomposition, we can write:

$$A = \sum_{i=1}^{k} |v_i\rangle\langle v_i| - \sum_{i=k+1}^{n} |v_i\rangle\langle v_i|$$

so by the definition of operator functions we have:

$$\exp(iAx) = \sum_{i=1}^{k} \exp(ix) |v_i\rangle\langle v_i| + \sum_{i=k+1}^{n} \exp(-ix) |v_i\rangle\langle v_i|.$$

By Euler's identity we have:

$$\exp(iAx) = \sum_{i=1}^{k} \left(\cos(x) + i\sin(x)\right) |v_i\rangle\langle v_i| + \sum_{i=k+1}^{n} \left(\cos(x) - i\sin(x)\right) |v_i\rangle\langle v_i|.$$

Grouping terms, we obtain:

$$\exp(iAx) = \cos(x) \sum_{i=1}^{n} |v_i\rangle\langle v_i| + i\sin(x) \left(\sum_{i=1}^{k} |v_i\rangle\langle v_i| - \sum_{i=k+1}^{n} |v_i\rangle\langle v_i|\right).$$

Using the spectral decomposition and definition of I, we therefore obtain the desired relation:

$$\exp(iAx) = \cos(x)I + i\sin(x)A.$$

Since all of the Pauli matrices satisfy  $A^2=I$  (Exercise 2.41), for  $\theta \in \mathbb{R}$  we can apply this obtained relation to obtain:

$$\exp(-i\theta X/2) = \cos\left(\frac{\theta}{2}\right)I - i\sin\left(\frac{\theta}{2}\right)X = \begin{bmatrix} \cos\frac{\theta}{2} & -i\sin\frac{\theta}{2} \\ -i\sin\frac{\theta}{2} & \cos\frac{\theta}{2} \end{bmatrix}$$

$$\exp(-i\theta Y/2) = \cos\left(\frac{\theta}{2}\right)I - i\sin\left(\frac{\theta}{2}\right)Y = \begin{bmatrix} \cos\frac{\theta}{2} & -\sin\frac{\theta}{2} \\ \sin\frac{\theta}{2} & \cos\frac{\theta}{2} \end{bmatrix}$$

$$\exp(-i\theta Z/2) = \cos\left(\frac{\theta}{2}\right)I - i\sin\left(\frac{\theta}{2}\right)Z = \begin{bmatrix} \cos\frac{\theta}{2} - i\sin\frac{\theta}{2} & 0 \\ 0 & \cos\frac{\theta}{2} - i\sin\frac{\theta}{2} \end{bmatrix} = \begin{bmatrix} e^{-i\theta/2} & 0 \\ 0 & e^{i\theta/2} \end{bmatrix}$$

which verifies equations (4.4)-(4.6).

#### Exercise 4.3

Show that, up to a global phase, the  $\pi/8$  gate satisfies  $T=R_z(\pi/4)$ 

## Solution

Concepts Involved: Linear Algebra, Quantum Gates.

Recall that the T gate is defined as:

$$T = \begin{bmatrix} 1 & 0 \\ 0 & \exp(i\pi/4) \end{bmatrix}$$

We observe that:

$$R_z(\pi/4) = \begin{bmatrix} e^{-i\pi/8} & 0\\ 0 & e^{i\pi/8} \end{bmatrix} = e^{-i\pi/8} \begin{bmatrix} 1 & 0\\ 0 & e^{i\pi/4} \end{bmatrix} = e^{-i\pi/8}T.$$

Exercise 4.4

Express the Hadamard gate H as a product of  $R_x$  and  $R_z$  rotations and  $e^{i\varphi}$  for some  $\varphi$ .

Solution

Concepts Involved: Linear algebra, Quantum Gates

We claim that  $H=R_z(\pi/2)R_x(\pi/2)R_z(\pi/2)$  up to a global phase of  $e^{-i\pi/2}$ . Doing a computation to verify this claim, we see that:

$$R_{z}(\pi/2)R_{x}(\pi/2)R_{z}(\pi/2) = \begin{bmatrix} e^{-i\pi/4} & 0 \\ 0 & e^{i\pi/4} \end{bmatrix} \begin{bmatrix} \cos\frac{\pi}{4} & -i\sin\frac{\pi}{4} \\ -i\sin\frac{\pi}{4} & \cos\frac{\pi}{4} \end{bmatrix} \begin{bmatrix} e^{-i\pi/4} & 0 \\ 0 & e^{i\pi/4} \end{bmatrix}$$

$$= \begin{bmatrix} e^{-i\pi/4} & 0 \\ 0 & e^{i\pi/4} \end{bmatrix} \begin{bmatrix} \frac{1}{\sqrt{2}} & -\frac{i}{\sqrt{2}} \\ -\frac{i}{\sqrt{2}} & \frac{1}{\sqrt{2}} \end{bmatrix} \begin{bmatrix} e^{-i\pi/4} & 0 \\ 0 & e^{i\pi/4} \end{bmatrix}$$

$$= \frac{1}{\sqrt{2}} \begin{bmatrix} e^{-i\pi/4} & 0 \\ 0 & e^{i\pi/4} \end{bmatrix} \begin{bmatrix} e^{-i\pi/4} & -ie^{i\pi/4} \\ -ie^{-i\pi/4} & e^{i\pi/4} \end{bmatrix}$$

$$= \frac{1}{\sqrt{2}} \begin{bmatrix} e^{-i\pi/2} & -i \\ -i & e^{i\pi/2} \end{bmatrix}$$

$$= \frac{1}{\sqrt{2}} \begin{bmatrix} e^{-i\pi/2} & -e^{-i\pi/2} \\ -e^{-i\pi/2} & e^{i\pi/2} \end{bmatrix}$$

$$= \frac{e^{-i\pi/2}}{\sqrt{2}} \begin{bmatrix} 1 & 1 \\ 1 & -1 \end{bmatrix}$$

$$= e^{-i\pi/2}H$$

Remark: If you are more algebraically minded, the following may appeal to you.

$$\begin{split} R_z(\pi/2)R_x(\pi/2)R_z(\pi/2) &= \frac{1}{2\sqrt{2}} \left(1 - iZ\right) \left(1 - iX\right) \left(1 - iZ\right) \\ &= \frac{1}{2\sqrt{2}} \left(1 - iZ - iX - ZX\right) \left(1 - iZ\right) \\ &= \frac{1}{2\sqrt{2}} \left(1 - iZ - iX - ZX - iZ - XZ - 1 + iZXZ\right) \\ &= \frac{1}{2\sqrt{2}} (-2iX - 2iZ) \quad \text{(using } ZXZ = -X\text{)} \\ &=: -iH \end{split}$$

#### Exercise 4.5

Prove that  $(\hat{\mathbf{n}}\cdot\boldsymbol{\sigma})^2=I$ , and use this to verify Equation (4.8)

### Solution

Concepts Involved: Linear Algebra

Expanding out the expression, we see that:

$$(\hat{\mathbf{n}} \cdot \boldsymbol{\sigma})^2 = (n_x X + n_y Y + n_z Z)^2$$

$$= n_x^2 X^2 + n_y^2 Y^2 + n_z^2 Z^2 + n_x n_y (XY + YX) + n_x n_z (XZ + ZX) + n_y n_z (YZ + ZY)$$

Using the result from Exercise 2.41 that  $\{\sigma_i, \sigma_i\} = 0$  if  $i \neq j$  and  $\sigma_i^2 = I$ , we have that:

$$(\hat{\mathbf{n}} \cdot \boldsymbol{\sigma})^2 = (n_x^2 + n_y^2 + n_z^2)I = I$$

where we use the fact that  $\hat{\mathbf{n}}$  is a vector of unit length. With this shown, we can use the result of Exercise 4.2 to conclude that:

$$\exp(-i\theta\hat{\mathbf{n}}\cdot\boldsymbol{\sigma}/2) = \cos(\frac{\theta}{2}) - i\sin(\frac{\theta}{2})(\hat{\mathbf{n}}\cdot\boldsymbol{\sigma})$$

which verifies equation (4.8).

#### **Exercise 4.6: Bloch sphere interpretation of rotations**

(\*) One reason why the  $R_{\hat{\mathbf{n}}}(\theta)$  operators are referred to as rotation operators is the following fact, which you are to prove. Suppose a single qubit has a state represented by the Bloch vector  $\lambda$ . Then, the effects of the rotation  $R_{\hat{\mathbf{n}}}(\theta)$  on the state is to rotate it by an angle  $\theta$  about the  $\hat{\mathbf{n}}$  axis of the Bloch sphere. This fact explains the rather mysterious looking factors of two in the definition of the rotation matrices.

#### Solution

Concepts Involved: Linear Algebra, Quantum Gates.

Let  $\lambda$  be an arbitrary Bloch vector. WLOG, we can express  $\lambda$  in a coordinate system such that  $\hat{\mathbf{n}}$  is aligned with the  $\hat{\mathbf{z}}$  axis, so it suffices to consider how the state behaves under application  $R_z(\theta)$ . Let  $\lambda = (\lambda_x, \lambda_y, \lambda_z)$  be the vector expressed in this coordinate system. By Exercise 2.72, the density operator corresponding to this Bloch vector is given by:

$$\rho = \frac{I + \lambda \cdot \sigma}{2}$$

We now observe how  $\rho$  transforms under conjugation by  $R_z(\theta)$ :

$$\begin{split} R_z(\theta)\rho R_z(\theta)^\dagger &= R_z(\theta)\rho R_z(-\theta) \\ &= R_z(\theta) \left(\frac{I + \lambda_x X + \lambda_y Y + \lambda_z Z}{2}\right) R_z(-\theta) \end{split}$$

Using that XZ = -ZX from Exercise 2.41, we make the observation that:

$$R_z(\theta)X = \left(\cos\left(\frac{\theta}{2}\right)I - i\sin\left(\frac{\theta}{2}\right)Z\right)X$$
$$= X\left(\cos\left(\frac{\theta}{2}\right)I + i\sin\left(\frac{\theta}{2}\right)Z\right)$$
$$= X\left(\cos\left(\frac{-\theta}{2}\right)I - i\sin\left(\frac{-\theta}{2}\right)Z\right)$$
$$= XR_z(-\theta)$$

Similarly, we find that  $R_z(\theta)Y=R_z(-\theta)Y$  (same anticommutation) and that  $R_z(\theta)Z=ZR_z(\theta)$  (all terms commute). With this, the expression for  $R_z(\theta)\rho R_z(\theta)^{\dagger}$  simplifies to:

$$R_{z}(\theta)\rho R_{z}(\theta)^{\dagger} = R_{z}(\theta) \left(\frac{I + \lambda_{x}X + \lambda_{y}Y + \lambda_{z}Z}{2}\right) R_{z}(-\theta)$$

$$= \left(\frac{IR_{z}(\theta) + \lambda_{x}XR_{z}(-\theta) + \lambda_{y}YR_{z}(-\theta) + \lambda_{z}ZR_{z}(\theta)}{2}\right) R_{z}(-\theta)$$

$$= \frac{I + \lambda_{x}XR_{z}(-2\theta) + \lambda_{y}YR_{z}(-2\theta) + \lambda_{z}Z}{2}$$

Calculating each of the terms in the above expression, we have:

$$XR_z(-2\theta) = X \left( \cos\left(\frac{-2\theta}{2}\right) - i\sin\left(\frac{-2\theta}{2}\right)Z \right)$$

$$= X \left( \cos(\theta) + i\sin(\theta)Z \right)$$

$$= \cos(\theta)X + i\sin(\theta)XZ$$

$$= \cos(\theta)X + i\sin(\theta)(-iY)$$

$$= \cos(\theta)X + \sin(\theta)Y$$

$$YR_z(-2\theta) = Y \left(\cos(\theta) + i\sin(\theta)Z\right)$$
$$= \cos(\theta)Y + i\sin(\theta)YZ$$
$$= \cos(\theta)Y + i\sin(\theta)(iX)$$
$$= \cos(\theta)Y - \sin(\theta)X.$$

Plugging these back into the expression for  $R_z(\theta)\rho R_z(\theta)^{\dagger}$  and collecting like terms, we have:

$$R_z(\theta)\rho R_z(\theta)^{\dagger} = \frac{I + (\lambda_x \cos(\theta) - \lambda_y \sin(\theta))X + (\lambda_x \sin(\theta) + \lambda_y \cos(\theta))Y + \lambda_z Z}{2}.$$

From this expression, we can read off the new Bloch vector  $\lambda'$  after conjugation by  $R_z(\theta)$  to be:

$$\lambda' = (\lambda_x \cos(\theta) - \lambda_y \sin(\theta), \lambda_x \sin(\theta) + \lambda_y \cos(\theta), \lambda_z).$$

Alternatively, suppose we apply the 3-dimensional rotation matrix  $A_z(\theta)$  to the original bloch vector  $\lambda$ . We have that:

$$A_z(\theta)\lambda = \begin{bmatrix} \cos \theta & -\sin \theta & 0 \\ \sin \theta & \cos \theta & 0 \\ 0 & 0 & 1 \end{bmatrix} \begin{bmatrix} \lambda_x \\ \lambda_y \\ \lambda_z \end{bmatrix} = \begin{bmatrix} \lambda_x \cos \theta - \lambda_y \sin \theta \\ \lambda_x \sin \theta + \lambda_y \cos \theta \\ \lambda_z \end{bmatrix}.$$

We see that we end up with the same resulting vector  $\lambda'$ . We conclude that the conjugation of  $\rho$  under  $R_z(\theta)$  has the equivalent effect to rotating the Bloch vector by  $\theta$  about the  $\hat{\mathbf{z}}$ -axis, and hence the effect of  $R_{\hat{\mathbf{n}}}(\theta)$  on a one qubit state is to rotate it by an angle  $\theta$  about  $\hat{\mathbf{n}}$ .

## Exercise 4.7

Show that XYX = -Y and use this to prove that  $XR_y(\theta)X = R_y(-\theta)$ .

#### Solution

Concepts Involved: Linear Algebra, Quantum Gates.

For the first claim, we use that XY = -YX and  $X^2 = I$  (Exercise 2.41) to obtain that:

$$XYX = -YXX = -YI = -Y.$$

Using this, we have that:

$$XR_{y}(\theta)X = X\left(\cos\left(\frac{\theta}{2}\right)I - i\sin\left(\frac{\theta}{2}\right)Y\right)X$$

$$= \cos\left(\frac{\theta}{2}\right)XIX - i\sin\left(\frac{\theta}{2}\right)XYX$$

$$= \cos\left(\frac{\theta}{2}\right)I + i\sin\left(\frac{\theta}{2}\right)Y$$

$$= \cos\left(-\frac{\theta}{2}\right)I - i\sin\left(-\frac{\theta}{2}\right)Y$$

$$= R_{y}(-\theta).$$

## Exercise 4.8

An arbitrary single qubit unitary operator can be written in the form

$$U = \exp(i\alpha)R_{\hat{\mathbf{n}}}(\theta)$$

for some real numbers  $\alpha$  and  $\theta$ , and a real three-dimensional unit vector  $\hat{\mathbf{n}}$ .

- 1. Prove this fact.
- 2. Find values for  $\alpha, \theta$ , and  $\hat{\mathbf{n}}$  giving the Hadamard gate H.
- 3. Find values for  $\alpha, \theta$ , and  $\hat{\mathbf{n}}$  giving the phase gate

$$S = \begin{bmatrix} 1 & 0 \\ 0 & i \end{bmatrix}$$

## Solution

Concepts Involved: Linear Algebra, Unitary Operators, Quantum Gates

- 1. By definition, for any unitary operator U we have that  $U^{\dagger}U = I$ , so for any state vector  $\langle \psi | \psi \rangle = \langle \psi | U^{\dagger}U | \psi \rangle$ . Therefore, all unitary Us are norm-preserving, and hence for a single qubit correspond to some reflection/rotation in 3-dimensional space (up to a global phase factor). Hence, we can write  $U = \exp(i\alpha) R_{\hat{\mathbf{n}}}(\theta)$  for some  $\hat{\mathbf{n}}$  (rotation axis),  $\theta$  (rotation angle) and  $\alpha$  (global phase).
- 2. Using the fact that  $H=\frac{X+Z}{\sqrt{2}}$ , and that modulo a factor of i that X/Z correspond to rotations

 $R_x(\pi)$  and  $R_z(\pi)$ , we find that:

$$H = \frac{iR_x(\pi) + iR_z(\pi)}{\sqrt{2}} = i\left(\frac{2\cos\left(\frac{\pi}{2}\right)I - i\sin\left(\frac{\pi}{2}\right)X - i\sin\left(\frac{\pi}{2}\right)Z}{\sqrt{2}}\right)$$
$$= i\left(\cos\left(\frac{\pi}{2}\right)I - i\sin\left(\frac{\pi}{2}\right)\left(\frac{1}{\sqrt{2}}X + 0Y + \frac{1}{\sqrt{2}}Z\right)\right)$$
$$= e^{i\pi/2}\left(\cos\left(\frac{\pi}{2}\right)I - i\sin\left(\frac{\pi}{2}\right)\left(\frac{1}{\sqrt{2}}X + 0Y + \frac{1}{\sqrt{2}}Z\right)\right)$$

Note that in the second last equality we use that  $\cos\left(\frac{\pi}{2}\right)=0$  and hence  $\frac{2}{\sqrt{2}}\cos\left(\frac{\pi}{2}\right)=\cos\left(\frac{\pi}{2}\right)$ . From the last expression, we can read off using the definition of  $R_{\hat{\mathbf{n}}}(\theta)$  that  $\hat{\mathbf{n}}=\left(\frac{1}{\sqrt{2}},0,\frac{1}{\sqrt{2}}\right)$ ,  $\theta=\pi$ , and  $\alpha=\frac{\pi}{2}$ .

3. We observe that:

$$R_z\left(\frac{\pi}{2}\right) = \begin{bmatrix} e^{-i\pi/4} & 0\\ 0 & e^{i\pi/4} \end{bmatrix} = e^{-i\pi/4} \begin{bmatrix} 1 & 0\\ 0 & i \end{bmatrix}$$

Hence:

$$S = e^{i\pi/4} R_z \left(\frac{\pi}{2}\right)$$

from which we obtain that  $\hat{\mathbf{n}} = \hat{\mathbf{z}} = (0,0,1)$ ,  $\theta = \frac{\pi}{2}$ , and  $\alpha = \frac{\pi}{4}$ .

**Remark:** For part (2), one just can use the definition

$$R_{\hat{n}}(\theta) \equiv \exp\left(-i\theta \hat{n} \cdot \vec{\sigma}/2\right) = \cos\left(\frac{\theta}{2}\right) I - i \sin\left(\frac{\theta}{2}\right) \left(n_x X + n_y Y + n_z Z\right),$$

and the fact  $H=\left(X+Z\right)/\sqrt{2}$ , to arrive at  $\cos\left(\frac{\theta}{2}\right)=0$ ,  $n_{x}=n_{z}=\frac{1}{\sqrt{2}}, n_{y}=0$ .

#### Exercise 4.9

Explain why any single qubit unitary operator may be written in the form (4.12).

#### Solution

**Concepts Involved:** Linear Algebra, Unitary Operators, Quantum Gates. Recall that (4.12) states that we can write any single qubit unitary U as:

$$U = \begin{bmatrix} e^{i(\alpha - \beta/2 - \delta/2)} \cos \frac{\gamma}{2} & -e^{i(\alpha - \beta/2 + \delta/2)} \sin \frac{\gamma}{2} \\ e^{i(\alpha + \beta/2 - \delta/2)} \sin \frac{\gamma}{2} & e^{i(\alpha + \beta/2 + \delta/2)} \cos \frac{\gamma}{2} \end{bmatrix}$$

where  $\alpha, \beta, \gamma, \delta \in \mathbb{R}$ .

Let U be a single qubit unitary operator. We then have that  $U^{\dagger}U = I$ , so identifying:

$$U = \begin{bmatrix} a & b \\ c & d \end{bmatrix} = \begin{bmatrix} \mathbf{v}_1 & \mathbf{v}_2 \end{bmatrix}$$

we obtain that:

$$\begin{bmatrix} \left| a \right|^2 + \left| c \right|^2 & a^*b + c^*d \\ ab^* + cd^* & \left| b \right|^2 + \left| d \right|^2 \end{bmatrix} = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix}$$

From the diagonal entries we obtain that  $|\mathbf{v}_1| = |\mathbf{v}_2| = 1$  and from the off diagonal entries we obtain that  $\langle \mathbf{v}_1, \mathbf{v}_2 \rangle = 0$  and hence the columns of U are orthonormal. From the fact that  $|v|_1$  is normalized, we can parameterize the magnitude of the entries with  $\gamma \in \mathbb{R}$  such that:

$$|a| = \cos \frac{\gamma}{2}, \quad |c| = \sin \frac{\gamma}{2}.$$

From the orthogonality, we further obtain that  $b=-c^*$  and  $d=a^*$ , from which we have that |b|=|c| and |d|=|a|. Furthermore, (also from the orthogonality) we can parameterize  $\arg(a)=-\frac{\beta}{2}-\frac{\delta}{2}$  and  $\arg(b)=\frac{\beta}{2}-\frac{\delta}{2}$  For  $\beta,\delta\in\mathbb{R}$ . Finally, multiplying U by a complex phase  $e^{i\alpha}$  for  $\alpha\in\mathbb{R}$  preserves the unitarity of U and the orthonormality of the colums. Combining these facts gives the form of (4.12) as desired.

## Exercise 4.10: X - Y decomposition of rotations

Give a decomposition analogous to Theorem 4.1 but using  $R_x$  instead of  $R_z$ .

#### Exercise 4.11

Suppose  $\hat{\mathbf{m}}$  and  $\hat{\mathbf{n}}$  are non-parallel real unit vectors in three dimensions. Use Theorem 4.1 to show that an arbitrary single qubit unitary U may be written

$$U = e^{i\alpha} R_{\hat{\mathbf{n}}}(\beta) R_{\hat{\mathbf{m}}}(\gamma) R_{\hat{\mathbf{n}}}(\delta)$$

#### Exercise 4.12

Give A, B, C, and  $\alpha$  for the Hadamard gate.

#### Solution

Concepts Involved: Linear Algebra, Decomposition of Rotations.

Recall that any single qubit unitary U can be written as  $U=e^{i\alpha}AXBXC$  where ABC=I and  $\alpha\in\mathbb{R}.$ 

First, observe that we can write:

$$H = \frac{1}{\sqrt{2}} \begin{bmatrix} 1 & 1 \\ 1 & -1 \end{bmatrix} = i \begin{bmatrix} -i & 0 \\ 0 & i \end{bmatrix} \begin{bmatrix} \frac{1}{\sqrt{2}} & \frac{1}{\sqrt{2}} \\ \frac{-1}{\sqrt{2}} & \frac{1}{\sqrt{2}} \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix} = e^{i\pi/2} R_z(\pi) R_y(-\pi/2) R_z(0)$$

so defining A,B,C according to the proof of Corollary 4.2, we have:

$$A = R_z(\pi)R_y(-\pi/4)$$
  

$$B = R_y(\pi/4)R_z(-\pi/2)$$
  

$$C = R_z(-\pi/2)$$

and  $\alpha = \frac{\pi}{2}$ .

## **Exercise 4.13: Circuit identities**

It is useful to be able to simplify circuits by inspection, using well-known identities. Prove the following three identities:

$$HXH = Z; \quad HYH = -Y; \quad HZH = X.$$

#### Solution

Concepts Involved: Linear Algebra, Quantum Gates.

By computation, we find:

$$\begin{split} HXH &= \frac{1}{2} \begin{bmatrix} 1 & 1 \\ 1 & -1 \end{bmatrix} \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} \begin{bmatrix} 1 & 1 \\ 1 & -1 \end{bmatrix} = \frac{1}{2} \begin{bmatrix} 2 & 0 \\ 0 & -2 \end{bmatrix} = \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} = Z \\ HYH &= \frac{1}{2} \begin{bmatrix} 1 & 1 \\ 1 & -1 \end{bmatrix} \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} \begin{bmatrix} 1 & 1 \\ 1 & -1 \end{bmatrix} = \frac{1}{2} \begin{bmatrix} 0 & -2i \\ 2i & 0 \end{bmatrix} = -\begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} = -Y \\ HZH &= \frac{1}{2} \begin{bmatrix} 1 & 1 \\ 1 & -1 \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} \begin{bmatrix} 1 & 1 \\ 1 & -1 \end{bmatrix} = \frac{1}{2} \begin{bmatrix} 0 & 2 \\ 2 & 0 \end{bmatrix} = \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} = X \end{split}$$

**Remark:** Notice once we have proved HXH=Z, we can directly say HZH=H(HXH)H=X as  $H^2=I$ . If one wants to prove everything algebraically, the following calculation suffices.

$$\begin{split} HXH &:= \frac{1}{2}\left(X+Z\right)X\left(X+Z\right) = \frac{1}{2}\left(I+ZX\right)\left(X+Z\right) = \frac{1}{2}(X+Z+Z+XZX) = Z \\ HYH &:= \frac{1}{2}\left(X+Z\right)Y\left(X+Z\right) = \frac{1}{2}\left(XY+ZY\right)\left(X+Z\right) = \frac{1}{2}(XYX+ZXY+ZYX+ZYZ) = -Y \end{split}$$

#### Exercise 4.14

Use the previous exercise to show that  $HTH = R_x(\pi/4)$ , up to a global phase.

#### Solution

Concepts Involved: Linear Algebra, Quantum Gates.

From Exercise 4.3, we know that  $T=R_z(\pi/4)$  up to a global phase  $e^{-i\pi/8}$ . We hence have that:

$$HTH = e^{-i\pi/8}HR_z(\pi/4)H$$

$$= e^{-i\pi/8}H\left(\cos\left(\frac{\pi}{8}\right)I - i\sin\left(\frac{\pi}{8}\right)Z\right)H$$

$$= e^{-i\pi/8}\left(\cos\left(\frac{\pi}{8}\right)I - i\sin\left(\frac{\pi}{8}\right)X\right)$$

$$= e^{-i\pi/8}R_x(\pi/4)$$

where in the second last equality we use the previous exercise, as well as the fact that  $HIH=H^2=I$  from Exercise 2.52.

## **Exercise 4.15: Composition of single qubit operations**

The Bloch representation gives a nice way to visualize the effect of composing two rotations.

(1) Prove that if a rotation through an angle  $\beta_1$  about the axis  $\hat{\mathbf{n}}_1$  is followed by a rotation through an angle  $\beta_2$  about an axis  $\hat{\mathbf{n}}_2$ , then the overall rotation is through an angle  $\beta_{12}$  about an axis  $\hat{\mathbf{n}}_{12}$  given by

$$c_{12} = c_1 c_2 - s_1 s_2 \hat{\mathbf{n}}_1 \cdot \hat{\mathbf{n}}_2$$
  

$$s_{12} \hat{\mathbf{n}}_{12} = s_1 c_2 \hat{\mathbf{n}}_1 + c_1 s_2 \hat{\mathbf{n}}_2 - s_1 s_2 \hat{\mathbf{n}}_2 \times \hat{\mathbf{n}}_1,$$

where  $c_i = \cos(\beta_i/2)$ ,  $s_i = \sin(\beta_i/2)$ ,  $c_{12} = \cos(\beta_{12}/2)$ , and  $s_{12} = \sin(\beta_{12}/2)$ .

(2) Show that if  $\beta_1=\beta_2$  and  $\mathbf{\hat{n}}_1=\mathbf{\hat{z}}$  these equations simplify to

$$c_{12} = c^2 - s^2 \hat{\mathbf{z}} \cdot \hat{\mathbf{n}}_2$$
  
$$s_{12} \hat{\mathbf{n}}_{12} = sc(\hat{\mathbf{z}} + \hat{\mathbf{n}}_2) - s^2 \hat{\mathbf{n}}_2 \times \hat{\mathbf{z}}$$

## Solution

Concepts Involved: Linear Algebra

(1) It suffices to show that  $R_{\mathbf{\hat{n}}_2}(\beta_2)R_{\mathbf{\hat{n}}_1}(\beta_1)$  is equivalent to  $R_{\mathbf{\hat{n}}_{12}}(\beta_{12})$ .

$$\begin{split} R_{\hat{\mathbf{n}}_2}(\beta_2)R_{\hat{\mathbf{n}}_1}(\beta_1) &= (c_2I - is_2\hat{\mathbf{n}}_2 \cdot \boldsymbol{\sigma}) \cdot (c_1I - is_1\hat{\mathbf{n}}_1 \cdot \hat{\boldsymbol{\sigma}}) \\ &= c_2c_1I - i\left(c_1s_2\hat{\mathbf{n}}_2 \cdot \boldsymbol{\sigma} + c_2s_1\hat{\mathbf{n}}_1 \cdot \boldsymbol{\sigma}\right) - s_2s_1\underbrace{\left(\hat{\mathbf{n}}_2 \cdot \boldsymbol{\sigma}\right) \cdot \left(\hat{\mathbf{n}}_1 \cdot \boldsymbol{\sigma}\right)}_{(\hat{\mathbf{n}}_2 \cdot \hat{\mathbf{n}}_1)I + i(\hat{\mathbf{n}}_2 \times \hat{\mathbf{n}}_1) \cdot \boldsymbol{\sigma}} \\ &= \left[c_2c_1 - s_2s_1\left(\hat{\mathbf{n}}_2 \cdot \hat{\mathbf{n}}_1\right)\right]I - i\left[c_1s_2\hat{\mathbf{n}}_2 + c_2s_1\hat{\mathbf{n}}_1 + s_2s_1\left(\hat{\mathbf{n}}_2 \times \hat{\mathbf{n}}_1\right)\right] \cdot \boldsymbol{\sigma} \end{split}$$

Identifying this operation to a single rotation  $R_{\hat{\mathbf{n}}_{12}}(\beta_{12}) \equiv c_{12}I - is_{12}\hat{\mathbf{n}}_{12} \cdot \boldsymbol{\sigma}$ , we arrive at the required relations (up to a presumable typesetting error)

$$c_{12} = c_2 c_1 - s_2 s_1 \left( \hat{\mathbf{n}}_2 \cdot \hat{\mathbf{n}}_1 \right)$$
  
$$s_{12} \hat{\mathbf{n}}_{12} = c_1 s_2 \hat{\mathbf{n}}_2 + c_2 s_1 \hat{\mathbf{n}}_1 + s_2 s_1 \left( \hat{\mathbf{n}}_2 \times \hat{\mathbf{n}}_1 \right)$$

(2) Setting  $\beta_1=\beta_2$  and  $\hat{\bf n}_1=\hat{\bf z}$  in the formulas proven above combined with the fact that  $c=c_1=\cos(\beta_1/2)=\cos(\beta_2/2)=c_2$  (and similarly  $s=s_1=s_2$ ), we have:

$$c_{12} = c^2 - s^2 \hat{\mathbf{z}} \cdot \hat{\mathbf{n}}_2$$
  

$$s_{12} \hat{\mathbf{n}}_{12} = sc\hat{\mathbf{z}} + cs\hat{\mathbf{n}}_2 - s^2 \hat{\mathbf{n}}_2 \times \hat{\mathbf{z}} = sc(\hat{\mathbf{z}} + \hat{\mathbf{n}}_2) - s^2 \hat{\mathbf{n}}_2 \times \hat{\mathbf{z}}.$$

**Remark:** For the sake of completeness, we provide a proof of the identity used in part 1 of the solution. First note the familiar Pauli matrix relation  $\sigma_i \sigma_j = \delta_{ij} I + i \epsilon_{ijk} \sigma_k$  (Exercise 2.43). Now massaging this equation gives

$$a_i \sigma_i b_j \sigma_j = a_i b_j \delta_{ij} + i \left( a_i b_j \epsilon_{ijk} \right) \sigma_k$$
$$= (\mathbf{a} \cdot \mathbf{b}) I + i \left( \mathbf{a} \times \mathbf{b} \right)_k \sigma_k,$$

where we have used standard Einstein index notation. Thus in matrix form, we have

$$(\mathbf{a} \cdot \boldsymbol{\sigma}) \cdot (\mathbf{b} \cdot \boldsymbol{\sigma}) = (\mathbf{a} \cdot \mathbf{b}) I + i (\mathbf{a} \times \mathbf{b}) \cdot \boldsymbol{\sigma}.$$

#### Exercise 4.16

What is the  $4 \times 4$  unitary matrix for the circuit

$$x_2 - H$$

 $x_1$  ———

in the computational basis? What is the unitary matrix for the circuit

$$x_1$$
  $H$ 

#### Solution

Concepts Involved: Linear Algebra, Quantum Gates, Tensor Products.

The unitary matrix for the first circuit is given by:

$$I_1 \otimes H_2 = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix} \otimes \frac{1}{\sqrt{2}} \begin{bmatrix} 1 & 1 \\ 1 & -1 \end{bmatrix} = \frac{1}{\sqrt{2}} \begin{bmatrix} 1 & 1 & 0 & 0 \\ 1 & -1 & 0 & 0 \\ 0 & 0 & 1 & 1 \\ 0 & 0 & 1 & -1 \end{bmatrix}.$$

The unitary matrix for the second circuit is given by:

$$H_1\otimes I_2=rac{1}{\sqrt{2}}egin{bmatrix}1&1\1&-1\end{bmatrix}\otimesegin{bmatrix}1&0\0&1\end{bmatrix}=rac{1}{\sqrt{2}}egin{bmatrix}1&0&1&0\0&1&0&1\1&0&-1&0\0&1&0&-1\end{bmatrix}$$

# Exercise 4.17: Building a CNOT from controlled- $\!Z$ gates

$$\begin{bmatrix} 1 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & -1 \end{bmatrix}$$

#### Solution

Concepts Involved: Linear Algebra, Quantum Gates, Controlled Operations.

We showed in Exercise 4.13 that HZH=X. Hence, to obtain a CNOT gate from a single controlled Z gate, we can conjugate the target qubit with Hadamard gates:

We can verify this via matrix multiplication, using the result from the previous exercise:

$$(I_1 \otimes H_2)(CZ_{1,2})(I_1 \otimes H_2) = \frac{1}{2} \begin{bmatrix} 1 & 1 & 0 & 0 \\ 1 & -1 & 0 & 0 \\ 0 & 0 & 1 & 1 \\ 0 & 0 & 1 & -1 \end{bmatrix} \begin{bmatrix} 1 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & -1 \end{bmatrix} \begin{bmatrix} 1 & 1 & 0 & 0 \\ 1 & -1 & 0 & 0 \\ 0 & 0 & 1 & 1 \\ 0 & 0 & 1 & -1 \end{bmatrix}$$

$$= \frac{1}{2} \begin{bmatrix} 2 & 0 & 0 & 0 \\ 0 & 2 & 0 & 0 \\ 0 & 0 & 0 & 2 \\ 0 & 0 & 2 & 0 \end{bmatrix}$$

$$= CX_{1,2}$$

Remark:

$$CX_{1,2} := |0\rangle\langle 0| \otimes I + |0\rangle\langle 0| \otimes X$$
  
=  $|0\rangle\langle 0| \otimes HH + |0\rangle\langle 0| \otimes HZH$   
=:  $(I \otimes H)(CZ_{1,2})(I \otimes H)$ .

## Exercise 4.18

Show that

## Solution

Concepts Involved: Linear Algebra, Quantum Gates, Controlled Operations.

It suffices to verify that the two gates have the same effect on the 2-qubit computational basis states (as it will then follow by linearity that they will have the same effect on any such superposition of the basis states). Checking the 8 necessary cases, we then have that:

$$\begin{split} CZ_{1,2}(|0\rangle_1 \otimes |0\rangle_2) &= |0\rangle_1 \otimes |0\rangle_2 \\ CZ_{2,1}(|0\rangle_1 \otimes |0\rangle_2) &= |0\rangle_1 \otimes |0\rangle_2 \\ CZ_{1,2}(|1\rangle_1 \otimes |0\rangle_2) &= |1\rangle_1 \otimes Z|0\rangle_2 = |1\rangle_1 \otimes |0\rangle_2 \\ CZ_{2,1}(|1\rangle_1 \otimes |0\rangle_2) &= |1\rangle_1 \otimes |0\rangle_2 \\ CZ_{2,1}(|1\rangle_1 \otimes |0\rangle_2) &= |1\rangle_1 \otimes |0\rangle_2 \\ CZ_{1,2}(|0\rangle_1 \otimes |1\rangle_2) &= |0\rangle_1 \otimes |1\rangle_2 \\ CZ_{2,1}(|0\rangle_1 \otimes |1\rangle_2) &= Z|0\rangle_1 \otimes |1\rangle_2 = |0\rangle_1 \otimes |1\rangle_2 \\ CZ_{2,1}(|1\rangle_1 \otimes |1\rangle_2) &= Z|0\rangle_1 \otimes |1\rangle_2 = |1\rangle_1 \otimes -|1\rangle_2 = -(|1\rangle_1 \otimes |1\rangle_1) \\ CZ_{2,1}(|1\rangle_1 \otimes |1\rangle_2) &= Z|1\rangle_1 \otimes |1\rangle_2 = -|1\rangle_1 \otimes |1\rangle_2 = -(|1\rangle_1 \otimes |1\rangle_1) \end{split}$$

from which we observe equality for each. The claim follows.

**Remark:** More compactly, we have  $CZ_{1,2}|b_1b_2\rangle=|b_1\rangle\otimes Z^{b_1}|b_2\rangle=(-1)^{b_1.b_2}|b_1b_2\rangle$  for computational basis states  $b_1,b_2\in\{0,1\}$ .

Using this form we can write

$$CZ_{1,2}|b_1b_2\rangle = (-1)^{b_1.b_2}|b_1b_2\rangle$$

$$= (-1)^{b_2.b_1}|b_1b_2\rangle$$

$$= Z^{b_2}|b_1\rangle \otimes |b_2\rangle$$

$$=: CZ_{2,1}|b_1b_2\rangle.$$

#### Exercise 4.19: CNOT action on unitary matrices

The CNOT gate is a simple permutation whose action on a density matrix  $\rho$  is to rearrange the elements in the matrix. Write out this action explicitly in the computational basis.

#### Solution

Concepts Involved: Linear Algebra, Quantum Gates, Controlled Operations, Density Operators

Let  $\rho$  be an arbitrary density matrix corresponding to a 2-qubit state. In the computational basis, we can write  $\rho$  as:

$$\rho \cong \begin{bmatrix} a_{11} & a_{12} & a_{13} & a_{14} \\ a_{21} & a_{22} & a_{23} & a_{24} \\ a_{31} & a_{32} & a_{33} & a_{34} \\ a_{41} & a_{42} & a_{43} & a_{44} \end{bmatrix}.$$

Studying the action of the CNOT gate on this density matrix, we calculate:

$$CX_{1,2} \rho CX_{1,2} = \begin{bmatrix} 1 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & 0 & 1 \\ 0 & 0 & 1 & 0 \end{bmatrix} \begin{bmatrix} a_{11} & a_{12} & a_{13} & a_{14} \\ a_{21} & a_{22} & a_{23} & a_{24} \\ a_{31} & a_{32} & a_{33} & a_{34} \\ a_{41} & a_{42} & a_{43} & a_{44} \end{bmatrix} \begin{bmatrix} 1 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & 0 & 1 \\ 0 & 0 & 1 & 0 \end{bmatrix}$$

$$= \begin{bmatrix} a_{11} & a_{12} & a_{13} & a_{14} \\ a_{21} & a_{22} & a_{23} & a_{24} \\ a_{41} & a_{42} & a_{43} & a_{34} \\ a_{31} & a_{32} & a_{33} & a_{34} \end{bmatrix} \begin{bmatrix} 1 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & 0 & 1 \\ 0 & 0 & 1 & 0 \end{bmatrix}$$

$$= \begin{bmatrix} a_{11} & a_{12} & a_{14} & a_{13} \\ a_{21} & a_{22} & a_{24} & a_{23} \\ a_{41} & a_{42} & a_{44} & a_{33} \\ a_{31} & a_{32} & a_{34} & a_{34} \end{bmatrix}$$

#### Exercise 4.20: CNOT basis transformations

Unlike ideal classical gates, ideal quantum gates do not have (as electrical engineers say) 'high-impedance' inputs. In fact, the role of 'control' and 'target' are arbitrary — they depend on what basis you think of a device as operating in. We have described how the CNOT behaves with respect to the computational basis, and in this description the state of the control qubit is not changed. However, if we work in a different basis then the control qubit does change: we will show that its phase is flipped depending on the state of the 'target' qubit! Show that

Introducing basis states  $|\pm\rangle\equiv(|0\rangle\pm|1\rangle)/\sqrt{2}$ , use this circuit identity to show that the effect of a CNOT with the first qubit as control and the second qubit as target is as follows:

$$\begin{aligned} |+\rangle|+\rangle &\mapsto |+\rangle|+\rangle \\ |-\rangle|+\rangle &\mapsto |-\rangle|+\rangle \\ |+\rangle|-\rangle &\mapsto |-\rangle|-\rangle \\ |-\rangle|-\rangle &\mapsto |+\rangle|-\rangle. \end{aligned}$$

Thus, with respect to this new basis, the state of the target qubit is not changed, while the state of the control qubit is flipped if the target starts as  $|-\rangle$ , otherwise it is left alone. That is, in this basis, the target and control have essentially interchanged roles!

#### Solution

Concepts Involved: Linear Algebra, Quantum Gates, Controlled Operations.

First, we have that:

Now conjugating CNOT<sub>1,2</sub> under  $H_1 \otimes H_2$ , we have:

which proves the circuit identity. We know already that:

$$CX_{2,1}|0\rangle|0\rangle = |0\rangle|0\rangle$$

$$CX_{2,1}|1\rangle|0\rangle = |1\rangle|0\rangle$$

$$CX_{2,1}|0\rangle|1\rangle = |1\rangle|1\rangle$$

$$CX_{2,1}|1\rangle|1\rangle = |0\rangle|1\rangle$$

so using the proven circuit identity and the fact that  $H|0\rangle=|+\rangle, H|1\rangle=|-\rangle$ , we obtain the map:

$$\begin{aligned} |+\rangle|+\rangle &\mapsto |+\rangle|+\rangle \\ |-\rangle|+\rangle &\mapsto |-\rangle|+\rangle \\ |+\rangle|-\rangle &\mapsto |-\rangle|-\rangle \\ |-\rangle|-\rangle &\mapsto |+\rangle|-\rangle \end{aligned}$$

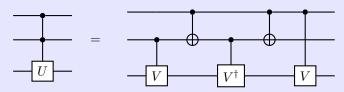
which is exactly what we wanted to prove.

Remark: Algebraically,

$$(H \otimes H)CX_{1,2}(H \otimes H) = (H \otimes H)(I \otimes H)(CZ_{1,2})(I \otimes H)(H \otimes H)$$
$$= (H \otimes I)(CZ_{1,2})(H \otimes I)$$
$$= (H \otimes I)(CZ_{2,1})(H \otimes I)$$
$$= CX_{2,1}.$$

#### Exercise 4.21

Verify that Figure 4.8 implements the  $C^2(U)$  operation.



#### Solution

#### **Concepts Involved:**

Let us analyse the second circuit. First, note that circuit in symmetric in the top two control registers. If  $c_1=0$ , the transformation applied to the final register is either I or  $VV^\dagger=I$ . The same is the case if  $c_2=0$ . When  $c_1=c_2=1$ , the unitary applied to the final qubit is VV=U.

**Remark:** Note that the  $V^{\dagger}$  in not applied for  $c_1=c_2=1$  because of the first CNOT gate. This is the clever part of the construction.

# Exercise 4.22

Prove that a  $C^2(U)$  gate (for any single qubit unitary U) can be constructed using at most eight one-qubit gates, and six controlled-NOTs.

#### Exercise 4.23

Construct a  $C^1(U)$  gate for  $U=R_x(\theta)$  and  $U=R_y(\theta)$ , using only CNOT and single qubit gates. Can you reduce the number of single qubit gates needed in the construction from three to two?

## Exercise 4.24

Verify that Figure 4.9 implements the Toffoli gate.

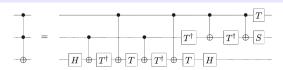


Figure 4.9. Implementation of the Toffoli gate using Hadamard, phase, controlled-NOT and  $\pi/8$  gates

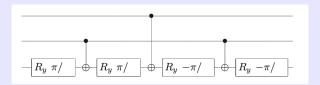
#### **Exercise 4.25: Fredkin gate construction**

Recall that the Fredkin (controlled-swap) gate performs the transform

- (1) Give a quantum circuit which uses three Toffoli gates to construct the Fredkin gate (*Hint:* think of the swap gate construction you can control each gate, one at a time).
- (2) Show that the first and last Toffoli gates can be replaced by CNOT gates.
- (3) Now replace the middle Toffoli gate with the circuit in Figure 4.8 to obtain a Fredkin gate construction using only six two-qubit gates.
- (4) Can you come up with an even simpler construction, with only five two-qubit gates?

#### Exercise 4.26

Show that the circuit:



differs by a Toffoli gate only by relative phases. That is, the circuit that takes  $|c_1,c_2,t\rangle$  to  $e^{i\theta(c_1,c_2,t)}\,|c_1,c_2,t\oplus c_1\cdot c_2\rangle$ , where  $e^{i\theta(c_1,c_2,t)}$  is some relative phase factor. Such gates can be sometimes be useful in experimental implementations, where it may be much easier to implement a gate that is the same as the Toffoli gate up to relative phases than it is to do the Toffoli directly.

#### Exercise 4.27

Using just CNOTs and Toffoli gates, construct a quantum circuit to perform the transformation

This kind of partial cyclic permutation operation will be useful later, in Chapter 7.

## Exercise 4.28

For  $U=V^2$  with V unitary, construct a  $C^5(U)$  gate analogous to that in Figure 4.10, but using no work qubits. You may use controlled-V and controlled- $V^\dagger$  gates.

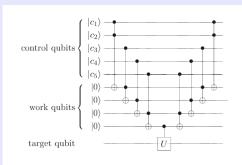


Figure 4.10. Network implementing the  $C^n(U)$  operation, for the case n=5

## Exercise 4.29

Find a circuit containing  $O(n^2)$  Toffoli, CNOT and single qubit gates which implements a  $C^n(X)$  gate (for n>3), using no work qubits.

# Exercise 4.30

Suppose U is a single qubit unitary operation. Find a circuit containing  $O(n^2)$  Toffoli, CNOT and single qubit gates which implements a  $C^n(U)$  gate (for n>3), using no work qubits.

#### **Exercise 4.31: More circuit identities**

Let subscripts denote which qubit an operator acts on, and let C be a CNOT with qubit 1 the control qubit and qubit 2 the target qubit. Prove the following identities:

$$CX_1C = X_1X_2$$
  
 $CY_1C = Y_1X_2$   
 $CZ_1C = Z_1$   
 $CX_2C = X_2$   
 $CY_2C = Z_1Y_2$   
 $CZ_2C = Z_1Z_2$   
 $R_{z,1}(\theta)C = CR_{z,1}(\theta)$   
 $R_{x,2}(\theta)C = CR_{x,2}(\theta)$ .

## Solution

#### Concepts Involved:

• In the computation basis, we have

$$\begin{split} CX_1C|q_1,q_2\rangle &= CX_1|q_1,q_1\oplus q_2\rangle\\ &= C|\widetilde{q}_1,q_1\oplus q_2\rangle\\ &= |\widetilde{q}_1,\widetilde{q}_1\oplus q_1\oplus q_2\rangle\\ &= |\widetilde{q}_1,\widetilde{q}_2\rangle\\ &= X_1X_2|q_1,q_2\rangle. \end{split}$$

• Similarly,

$$CZ_{1}C|q_{1}, q_{2}\rangle = CZ_{1}|q_{1}, q_{1} \oplus q_{2}\rangle$$

$$= C(-1)^{q_{1}}|q_{1}, q_{1} \oplus q_{2}\rangle$$

$$= (-1)^{q_{1}}|q_{1}, q_{1} \oplus q_{1} \oplus q_{2}\rangle$$

$$= (-1)^{q_{1}}|q_{1}, q_{2}\rangle$$

$$= Z_{1}|q_{1}, q_{2}\rangle.$$

•

$$CY_1C = iCX_1Z_1C = iCX_1(CC)Z_1C = i(CX_1C)(CZ_1C) = iX_1X_2Z_1 = Y_1X_2.$$

•

$$CR_{z,1}(\theta)C = C(aI_1 + bZ_1)C = aI_1 + bCZ_1C = aI_1 + bZ_1 = R_{z,1}(\theta)$$

#### Exercise 4.32

Let  $\rho$  be the density matrix describing a two qubit system. Suppose we perform a projective measurement in the computational basis of the second qubit. Let  $P_0=|0\rangle\langle 0|$  and  $P_1=|1\rangle\langle 1|$  be the projectors onto the  $|0\rangle$  and the  $|1\rangle$  states of the second qubit, respectively. Let  $\rho'$  be the density matrix which would be assigned to the system after the measurement by an observer who did not learn the measurement result. Show that

$$\rho' = P_0 \rho P_0 + P_1 \rho P_1$$

Also show that the reduced density matrix for the first qubit is not affected by the measurement, that is  $\operatorname{tr}_2(\rho) = \operatorname{tr}_2(\rho')$ .

#### Solution

#### Concepts Involved:

The density matrix for the system after the measurement is given by

$$\rho' = p_0 \rho_0 + p_1 \rho_1 \tag{2}$$

$$=\operatorname{Tr}(\rho P_0)\frac{P_0\rho P_0}{\operatorname{Tr}(\rho P_0)}+\operatorname{Tr}(\rho P_1)\frac{P_1\rho P_1}{\operatorname{Tr}(\rho P_1)}\tag{3}$$

$$=P_0\rho P_0 + P_1\rho P_1 \tag{4}$$

Finally, we have

$$\operatorname{Tr}_{2}(\rho') = \operatorname{Tr}_{2}(P_{0}\rho P_{0} + P_{1}\rho P_{1})$$

$$= \operatorname{Tr}_{2}(\rho P_{0}^{2}) + \operatorname{Tr}_{2}(\rho P_{1}^{2})$$

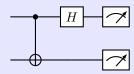
$$= \operatorname{Tr}_{2}(\rho P_{0}) + \operatorname{Tr}_{2}(\rho P_{1})$$

$$= \operatorname{Tr}_{2}(\rho(P_{0} + P_{1}))$$

$$= \operatorname{Tr}_{2}(\rho)$$

#### Exercise 4.33: Measurement in the Bell basis

The measurement model we have specified for the quantum circuit model is that measurements are performed only in the computational basis. However, often we want to perform a measurement in some other basis, defined by a complete set of orthonormal states. To perform this measurement, simply unitarily transform from the basis we wish to perform the measurement in to the computational basis, then measure. For example, show that the circuit



performs a measurement in the basis of the Bell states. More precisely, show that this circuit results in a measurement being performed with corresponding POVM elements the four projectors onto the Bell states. What are the corresponding measurement operators?

#### Solution

#### Concepts Involved:

The full action of the circuit is equivalent to the following four operators

$$\begin{split} M'_{00} &= |00\rangle \left( \langle 00| \left( H_1 \otimes I_2 \right) \text{CNOT} \right) = \frac{1}{\sqrt{2}} |00\rangle (\langle 00| + \langle 11|) = |00\rangle \langle \beta_{00}| \\ M'_{01} &= |01\rangle \left( \langle 01| \left( H_1 \otimes I_2 \right) \text{CNOT} \right) = \frac{1}{\sqrt{2}} |01\rangle (\langle 01| + \langle 10|) = |01\rangle \langle \beta_{01}| \\ M'_{10} &= |10\rangle \left( \langle 10| \left( H_1 \otimes I_2 \right) \text{CNOT} \right) = \frac{1}{\sqrt{2}} |10\rangle (\langle 00| - \langle 11|) = |10\rangle \langle \beta_{10}| \\ M'_{11} &= |11\rangle \left( \langle 11| \left( H_1 \otimes I_2 \right) \text{CNOT} \right) = \frac{1}{\sqrt{2}} |11\rangle (\langle 01| - \langle 10|) = |11\rangle \langle \beta_{11}| \\ \end{split}$$

Therefore, it is easy to see that  $\left(M_k'\right)^\dagger M_k'$  are projections onto the Bell states

$$(M'_{00})^{\dagger} M_{00} = |\beta_{00}\rangle\langle\beta_{00}|$$

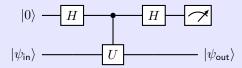
$$(M'_{01})^{\dagger} M_{01} = |\beta_{01}\rangle\langle\beta_{01}|$$

$$(M'_{10})^{\dagger} M_{10} = |\beta_{10}\rangle\langle\beta_{10}|$$

$$(M'_{11})^{\dagger} M_{11} = |\beta_{11}\rangle\langle\beta_{11}|$$

### Exercise 4.34: Measuring an operator

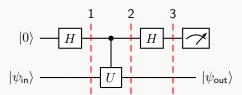
Suppose we have a single qubit operator U with eigenvalues  $\pm 1$ , so that U is both Hermitian and unitary, so it can be regarded as both an observable and a quantum gate. Suppose we wish to measure the observable U. That is, we desire to obtain a measurement result indicating one of the two eigenvalues, and leaving a post-measurement state which is the corresponding eigenvector. How can this be implemented by a quantum circuit? Show that the following circuit implements a measurement of U:



#### Solution

### Concepts Involved:

First, let us determine the action of this circuit on the eigenstates of U,  $|\psi_{in}\rangle = |u_{\pm}\rangle$ .



$$|\Psi(t_1)\rangle = \frac{1}{\sqrt{2}} (|0\rangle + |1\rangle) |u_{\pm}\rangle$$
.

$$\begin{split} |\Psi(t_2)\rangle &= \frac{1}{\sqrt{2}} \left( |0\rangle |u_{\pm}\rangle + |1\rangle U |u_{\pm}\rangle \right) \\ &= \frac{1}{\sqrt{2}} \left( |0\rangle |u_{\pm}\rangle \pm |1\rangle U |u_{\pm}\rangle \right) \\ &= |\pm\rangle |u_{\pm}\rangle \end{split}$$

$$|\Psi(t_3)\rangle = |0/1\rangle|u_+\rangle$$

This shows that the above circuit implements a measurement of U for input states,  $|\psi_{\rm in}\rangle=|u_\pm\rangle$ . Now, as U is also Hermitian,  $|u_\pm\rangle$  forms a full basis and thereby, we can extend the result to arbitrary input states simply via linearity.

# Exercise 4.35: Measurement commutes with controls

A consequence of the principle of deferred measurement is that measurements commute with quantum gates when the qubit being measured is a control qubit, that is:

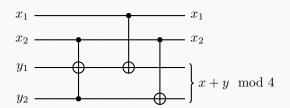
(Recall that the double lines represent classical bits in this diagram.) Prove the first equality. The rightmost circuit is simply a convenient notation to depict the use of a measurement result to classically control a quantum gate.

# Exercise 4.36

Construct a quantum circuit to add two two-bit numbers x and y modulo 4. That is, the circuit should perform the transformation  $|x,y\rangle\mapsto |x,x+y\operatorname{mod} 4\rangle$ .

#### Solution

#### **Concepts Involved:**



#### Exercise 4.37

Provide a decomposition of the transform

$$\frac{1}{2} \begin{bmatrix} 1 & 1 & 1 & 1 \\ 1 & i & -1 & -i \\ 1 & -1 & 1 & 1 \\ 1 & -i & -1 & i \end{bmatrix}$$

into a product of two-level unitaries. This is a special case of the quantum Fourier transform, which we study in more detail in the next chapter.

# Solution

#### **Concepts Involved:**

We provide the standard QFT decomposition

# Exercise 4.38

Prove that there exist a  $d \times d$  unitary matrix U which cannot be decomposed as a product of fewer than d-1 two-level unitary matrices.

# Exercise 4.39

Find a quantum circuit using single qubit operations and CNOTs to implement the transformation

$$\begin{bmatrix} 1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\ 0 & 0 & a & 0 & 0 & 0 & 0 & c \\ 0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 & 1 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 & 0 & 1 & 0 & 0 \\ 0 & 0 & 0 & 0 & 0 & 0 & 1 & 0 \\ 0 & 0 & b & 0 & 0 & 0 & 0 & d \end{bmatrix}$$

where  $\tilde{U} = \begin{bmatrix} a & c \\ b & d \end{bmatrix}$  is an arbitrary  $2 \times 2$  unitary matrix.

# Exercise 4.40

For arbitrary  $\alpha$  and  $\beta$  show that

$$E(R_n(\alpha), R_n(\alpha + \beta)) = \left| 1 - \exp(i\beta/2) \right|,$$

and use this to justify (4.76).

# Solution

#### Concepts Involved:

We start with the rotation operators

$$R_n(\alpha) = e^{-i\frac{\alpha}{2}\mathbf{n}\cdot\sigma}, R_n(\alpha+\beta) = e^{-i\frac{\alpha+\beta}{2}\mathbf{n}\cdot\sigma}.$$

The error between them is

$$E(R_n(\alpha), R_n(\alpha + \beta)) = ||R_n(\alpha) - R_n(\alpha + \beta)||_{\text{op}}.$$

$$= ||R_n(\alpha) - R_n(\alpha)R_n(\beta)||_{\text{op}}$$

$$= ||I - R_n(\beta)||_{\text{op}}$$

where we have used  $R_n(\alpha + \beta) = R_n(\alpha)R_n(\beta)$  and the fact that operator norms are invariant under unitary operators.

We now notice that  $I-R_n(\beta)$  is normal, i.e.  $(I-R_n(\beta))(I-R_n(\beta))^\dagger=(I-R_n(\beta))^\dagger(I-R_n(\beta))$ . For normal operators, the operator norm coincides with the largest eigenvalue.[Ref] The eigenvalues of  $I-R_n(\beta)$  are given by

$$1 - e^{i\frac{\beta}{2}}$$
 and  $1 - e^{-i\frac{\beta}{2}}$ .

Thus, the operator norm is

$$||I - R_n(\beta)||_{\text{op}} = \max\left(|1 - e^{i\frac{\beta}{2}}|, |1 - e^{-i\frac{\beta}{2}}|\right)$$
  
=  $|1 - e^{i\frac{\beta}{2}}|$ 

since both are equal. Thus, the error

$$E(R_n(\alpha), R_n(\alpha + \beta)) = |1 - e^{i\frac{\beta}{2}}|.$$

# Exercise 4.41

This and the next two exercises develop a construction showing that the Hadamard, phase, controlled-NOT and Toffoli gates are universal. Show that the circuit in Figure 4.17 applies the operation  $R_z(\theta)$  to the third (target) qubit if the measurement outcomes are both 0, where  $\cos\theta=3/5$ , and otherwise applies Z to the target qubit. Show that the probability of both measurement outcomes being 0 is 5/8, and explain how repeated use of this circuit and  $Z=S^2$  gates may be used to apply a  $R_z(\theta)$  gate with probability approaching 1.

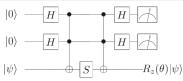


Figure 4.17. Provided both measurement outcomes are 0 this circuit applies  $R_z(\theta)$  to the target, where  $\cos \theta = 3/5$ . If some other measurement outcome occurs then the circuit applies Z to the target.

# Exercise 4.42: Irrationality of $\theta$

Suppose  $\cos \theta = 3/5$ . We give a proof by contradiction that  $\theta$  is an irrational multiple of  $2\pi$ .

- (1) Using the fact that  $e^{i\theta}=(3+4i)/5$ , show that if  $\theta$  is rational, then there must exist a positive integer m such that  $(3+4i)^m=5^m$ .
- (2) Show that  $(3+4i)^m = 3+4i \pmod 5$  for all m > 0, and conclude that no m such that  $(3+4i)^m = 5^m$  can exist.

#### Solution

# Concepts Involved:

If  $\theta$  is a rational multiple of  $2\pi$ , then there exist a positive integer m for which

$$e^{im\theta} = 1$$

$$\Rightarrow \frac{(3+4i)^m}{5^m} = 1$$

$$\Rightarrow (3+4i)^m = 5^m.$$

We want to prove that  $(3+4i)^m \equiv 3+4i \pmod 5$  for all m>0. Base Case (m=1):

$$(3+4i)^1 = 3+4i \pmod{5}$$
.

**Inductive Step:** Assume  $(3+4i)^k \equiv 3+4i \pmod 5$  for some k. We show this holds for k+1:

$$(3+4i)^{k+1} = (3+4i)^k \cdot (3+4i).$$

Using the inductive hypothesis:

$$(3+4i)^{k+1} \equiv (3+4i) \cdot (3+4i) \pmod{5}$$
.

Now, compute:

$$(3+4i)^2 = 9 + 24i + 16i^2 = -7 + 24i \equiv 3 + 4i \pmod{5}.$$

Thus,  $(3+4i)^{k+1} \equiv 3+4i \pmod{5}$ .

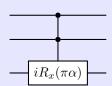
By induction,  $(3+4i)^m \equiv 3+4i \pmod{5}$  for all m>0.

Exercise 4.43

Use the results of the previous two exercises to show that the Hadamard, phase, controlled-NOT and Toffoli gates are universal for quantum computation.

Exercise 4.44

Show that the three qubit gate G defined by the circuit:



is universal for quantum computation whenever  $\alpha$  is irrational.

# Exercise 4.45

Suppose U is a unitary transform implemented by an n qubit quantum circuit constructed from H,S, CNOT and Toffoli gates. Show that U is of the form  $2^{-k/2}M$  for some integer k, where M is a  $2^n \times 2^n$  matrix with only complex integer entries. Repeat this exercise with the Toffoli gate replaced by the  $\pi/8$  gate.

# Exercise 4.46: Exponential complexity growth of quantum systems

Let  $\rho$  be a density matrix describing the state of n qubits. Show that describing  $\rho$  requires  $4^n-1$  independent real numbers.

# Solution

# **Concepts Involved:**

To show that a density matrix  $\rho$  for n qubits requires  $4^n-1$  independent real parameters, we proceed as follows:

# 1. Dimension of the Hilbert Space

The state space for n qubits is  $\mathbb{C}^{2^n}$ , which has dimension  $2^n$ .

# 2. Structure of the Density Matrix

A density matrix  $\rho$  is:

Hermitian:  $\rho = \rho^{\dagger}$ ,

Trace One:  $Tr(\rho) = 1$ ,

Positive Semidefinite:  $\langle v|\rho|v\rangle \geq 0$  for all vectors  $|v\rangle$ .

# 3. Counting Parameters

A. Hermitian Condition: A  $2^n \times 2^n$  Hermitian matrix has:

 $2^n$  real parameters for the diagonal,

$$\frac{(2^n)(2^n-1)}{2}$$
 complex parameters for the off-diagonal elements, contributing

 $(2^n)(2^n-1)$  real parameters.

Total:

$$2^n + (2^n)(2^n - 1) = 2^{2n}.$$

B. Trace Condition: The trace condition reduces the number of independent parameters by 1:

$$2^{2n} - 1$$
.

Thus, a density matrix  $\rho$  for n qubits indeed requires  $4^n-1$  independent real numbers.

# Exercise 4.47

For  $H=\sum_{k=1}^L H_k$ , prove that  $e^{-iHt}=e^{-iH_1t}e^{-iH_2t}\dots e^{-iH_Lt}$  for all i if  $[H_j,H_k]=0$ , for all j,k.

### Solution

### **Concepts Involved:**

We proceed by induction. For L=2 we have:

$$e^{-i(H_1+H_2)t} = \sum_{n=0}^{\infty} \frac{([-it]^n (H_1+H_2)^n}{n!}$$

Using the fact that  $[H_1, H_2] = 0$ , we can then write:

$$e^{-i(H_1+H_2)t} = \sum_{n=0}^{\infty} \frac{1}{n!} \sum_{k=0}^{n} \binom{n}{k} H_1^k H_2^{n-k} = \sum_{n=0}^{\infty} (-it)^n \sum_{k=0}^{n} \frac{H_1^k H_2^{n-k}}{k!(n-k)!} = \sum_{n=0}^{\infty} \frac{(-itH_1)^n}{n!} \sum_{m=0}^{\infty} \frac{(-itH_2)^m}{m!}$$

and so:

$$e^{-i(H_1+H_2)t} = e^{-iH_1t}e^{-iH_2t}$$

Suppose the claim holds for L-1. Then letting  $\sum_{k=1}^{L-1} H_k$  take the place of  $H_1$  and  $H_L$  take the place of  $H_2$  in the above argument (since  $[\sum_{k=1}^{L-1} H_k, H_L] = 0$ ), we find:

$$e^{-i\sum_{k=1}^{L}H_{k}t} = e^{-i(\sum_{k=1}^{L-1}H_{k}+H_{L})t} = e^{-i\sum_{k=1}^{L-1}H_{k}t}e^{-iH_{L}t} = \prod_{k=1}^{L-1}e^{-iH_{k}t}e^{-iH_{L}t} = \prod_{k=1}^{L}e^{-iH_{k}t}$$
(5)

where in the second-to-last equality we use the inductive hypothesis. Thus the L case holds and the claim is proven via induction.

# Exercise 4.48

Show that the restriction of  $H_k$  to at most c particles *implies* that in the sum (4.97), L is upper bounded by a polynomial in n.

# Exercise 4.49: Baker-Campbell-Hausdorf formula

Prove that

$$e^{(A+B)\Delta t} = e^{A\Delta t}e^{B\Delta t}e^{-\frac{1}{2}[A,B]\Delta t^2} + O(\Delta t^3)$$

and also prove Equations (4.103) and (4.104).

### Solution

#### **Concepts Involved:**

To prove that

$$e^{(A+B)\Delta t} = e^{A\Delta t}e^{B\Delta t}e^{-\frac{1}{2}[A,B]\Delta t^2} + O(\Delta t^3),$$

we use the Taylor expansion of the exponential

$$e^{(A+B)\Delta t} = I + (A+B)\Delta t + \frac{(A+B)^2(\Delta t)^2}{2} + O(\Delta t^3).$$

Calculating

$$(A + B)^2 = A^2 + AB + BA + B^2 = A^2 + B^2 + 2AB - [A, B].$$

Thus, we have

$$\begin{split} e^{(A+B)\Delta t} &= I + (A+B)\Delta t + \left(\frac{A^2 + B^2}{2} + AB - \frac{[A,B]}{2}\right)\Delta t^2 + O(\Delta t^3) \\ &= \left(I + A\Delta t + \frac{A^2 (\Delta t)^2}{2}\right)\left(I + B\Delta t + \frac{B^2 (\Delta t)^2}{2}\right)\left(I - \frac{1}{2}[A,B]\Delta t^2\right) + O(\Delta t^3) \\ &= e^{A\Delta t}e^{B\Delta t}e^{-\frac{1}{2}[A,B]\Delta t^2} + O(\Delta t^3). \end{split}$$

This completes the proof.

Exercise 4.50

Let  $H = \sum_{k}^{L} H_k$ , and define

$$U_{\Delta t} = \left[ e^{-iH_1\Delta t} e^{-iH_2\Delta t} \dots e^{-iH_L\Delta t} \right] \left[ e^{-iH_L\Delta t} e^{-iH_{L-1}\Delta t} \dots e^{-iH_1\Delta t} \right]$$

- (a) Prove that  $U_{\Delta t} = e^{-2iH\Delta t} + O(\Delta t^3)$
- (b) Use the results in Box 4.1 to prove that for a positive integer m,

$$E(U_{\Delta t}^m, e^{-2miH\Delta t}) \le m\alpha \Delta t^3,$$

for some constant  $\alpha$ .

# Exercise 4.51

Construct a quantum circuit to simulate the Hamiltonian

$$H = X_1 \otimes Y_2 \otimes Z_3$$

performing the unitary transform  $e^{-i\Delta tH}$  for any  $\Delta t$ .

# Problem 4.1: Computable phase shifts

Let m and n be positive integers. Suppose  $f: \{0, \dots, 2^m - 1\} \mapsto \{0, \dots, 2^n - 1\}$  is a classical function from m to n bits which may be computed reversibly using T Toffoli gates, as described in Section 3.2.5.

That is, the function  $(x,y)\mapsto (x,y\oplus f(x))$  may be implemented using T Toffoli gates. Give a quantum circuit using 2T+n (or fewer) one, two and three qubit gates to implement the unitary operation defined by

$$|x\rangle \mapsto \exp\left(\frac{-2i\pi f(x)}{2^n}\right)|x\rangle$$

### Problem 4.2

Find a depth  $O(\log n)$  construction for the  $C^n(X)$  gate. (Comment: The depth of a circuit is the number of distinct timesteps at which gates are applied; the point of this problem is that it is possible to parallelize the  $C^n(X)$  construction by applying many gates in parllel during the same timestep.)

#### Problem 4.3: Alternate universality construction

Suppose U is a unitary matrix on n qubits. Define  $H \equiv i \ln(U)$ . Show that

- (1) H is Hermitian, with eigenvalues in the range 0 to  $2\pi$ .
- (2) H can be written

$$H = \sum_{g} h_g g,$$

where  $h_g$  are real numbers and the sum is over all n-fold tensor products g of the Pauli matrices  $\{I, X, Y, Z\}$ .

- (3) Let  $\Delta=1/k$ , for some positive integer k. Explain how the unitary operation  $\exp(-ih_g g\Delta)$  may be implemented using O(n) one and two qubit operations.
- (4) Show that

$$\exp(-iH\Delta) = \prod_g \exp(-ih_g g\Delta) + O(4^n \Delta^2)$$

where the product is taken with respect to any fixed ordering of the n-fold tensor products of Pauli matrices, g.

(5) Show that

$$U = \left[\prod_{g} \exp(-ih_g h\Delta)\right]^k + O(4^n \Delta)$$

(6) Explain how to approximate U to within a distance  $\epsilon > 0$  using  $O(n16^n/\epsilon)$  one and two qubit unitary operations.

### Problem 4.4: Minimal Toffoli construction (Research)

The following problems concern constructions of the Toffoli with some minimal number of other gates. '

- (1) What is the smallest number of two qubit gates that can be used to implement the Toffoli gate?
- (2) What is the smallest number of one qubit gates and CNOT gates that can be used to implement the Toffoli gate?
- (3) What is the smallest number of one qubit gates and controlled-Z gates that can be used to implement the Toffoli gate?

# Problem 4.5: (Research)

Construct a family of Hamiltonians,  $\{H_n\}$ , on n qubits, such that simulating  $H_n$  requires a number of operations super-polynomial in n. (Comment: This problem seems to be quite difficult.)

# Problem 4.6: Universality with prior entanglement

Controlled-NOT gates and single qubit gates form a universal set of quantum logic gates. Show that an alternative universal set of resources is comprised of single qubit unitaries, the ability to perform measurements of pairs of qubits in the Bell basis, and the ability to prepare arbitrary four qubit entangled states.

# 8 Quantum noise and quantum operations

# Exercise 8.1: Unitary evolution as a quantum operation

Pure states evolve under unitary transforms as  $|\psi\rangle \mapsto U|\psi\rangle$ . Show that, equivalently, we may write  $\rho \mapsto \mathcal{E}(\rho) \equiv U\rho U^{\dagger}$ , for  $\rho = |\psi\rangle\langle\psi|$ .

#### Solution

Concepts Involved: Unitary Transformations, Density Operators, Quantum Operations

Under unitary evolution, a pure state maps to another pure state via  $|\psi\rangle\mapsto U\,|\psi\rangle$ . The density operator corresponding to the final state is then  $\rho_{U\psi}=U\,|\psi\rangle\langle\psi|\,U^\dagger=U\rho U^\dagger$ . We can see that this final state is equivalent to if we started in the density operator formalism  $\rho=|\psi\rangle\langle\psi|$  and acted on it with the quantum operation  $\mathcal{E}(\rho)=U\rho U^\dagger$ .

# Exercise 8.2: Measurement as a quantum operation

Recall from Section 2.2.3 (on page 84) that a quantum measurement with outcomes labeled by m is described by a set of measurement operators  $M_m$  such that  $\sum_m M_m^\dagger M_m = I$ . Let the state of the system immediately before the measurement be  $\rho$ . Show that for  $\mathcal{E}_m(\rho) \equiv M_m \rho M_m^\dagger$ , the state of the system immediately after the measurement is

$$\frac{\mathcal{E}_m(\rho)}{\operatorname{tr}\left(\mathcal{E}_m(\rho)\right)}$$

Also show that the probability of obtaining this measurement result is  $p(m) = \operatorname{tr} \left( \mathcal{E}_m(\rho) \right)$ .

# Solution

Concepts Involved: Measurement, Density Operators, Quantum Operations

It suffices to prove the claim for pure states  $\rho = |\psi\rangle\langle\psi|$ , as mixed states are a linear combination of pure states and  $\mathcal{E}_m(\cdot)$  and the trace are both linear in their arguments.

The postulate of quantum measurement tells us that the probability of measuring outcome m given state  $|\psi\rangle$  is:

$$p(m) = \langle \psi | M_m^{\dagger} M_m | \psi \rangle$$

For pure  $\rho = |\psi\rangle\langle\psi|$  we have:

$$\operatorname{tr}\left(\mathcal{E}_{m}(\rho)\right) = \operatorname{tr}\left(M_{m}\rho M_{m}^{\dagger}\right) = \operatorname{tr}\left(M_{m}\left|\psi\right\rangle\!\langle\psi\right|M_{m}^{\dagger}\right) = \left\langle\psi\right|M_{m}^{\dagger}M_{m}\left|\psi\right\rangle = p(m)$$

where we have used that  $\operatorname{tr}(|\alpha\rangle\langle\beta|) = \langle\alpha|\beta\rangle$ . The postulate of measurement tells us that the post-

measurement state (when measuring a pure state) is:

$$|\psi\rangle \to \frac{M_m |\psi\rangle}{\sqrt{\langle \psi | M_m^{\dagger} M_m |\psi\rangle}}.$$

The density operator corresponding to the post measurement state is:

$$\frac{M_{m}\left|\psi\right\rangle}{\sqrt{\left\langle\psi\right|M_{m}^{\dagger}M_{m}\left|\psi\right\rangle}}\left(\frac{\left\langle\psi\right|M_{m}^{\dagger}}{\sqrt{\left\langle\psi\right|M_{m}^{\dagger}M_{m}\left|\psi\right\rangle}}\right) = \frac{M_{m}\left|\psi\right\rangle\!\left\langle\psi\right|M_{m}^{\dagger}}{\left\langle\psi\right|M_{m}^{\dagger}M_{m}\left|\psi\right\rangle} = \frac{M_{m}\rho M_{m}^{\dagger}}{p(m)} = \frac{\mathcal{E}_{m}(\rho)}{\operatorname{tr}\left(\mathcal{E}_{m}(\rho)\right)}$$

which is what we wished to show.

# Exercise 8.3

Our derivation of the operator-sum representation implicitly assumed that the input and output spaces for the operation were the same. Suppose a composite system AB initially in an unknown quantum state  $\rho$  is brought into contact with a composite system CD initially in some standard state  $|0\rangle$ , and the two systems interact according to a unitary interaction U. After the interaction we discard systems A and D, leaving a state  $\rho'$  of system BC. Show that the map  $\mathcal{E}(\rho) = \rho'$  satisfies

$$\mathcal{E}(\rho) = \sum_{k} E_k \rho E_k^{\dagger}$$

for some set of linear operators  $E_k$  from the state space of system AB to the state space of system BC, and such that  $\sum_k E_k^{\dagger} E_k = I$ .

# Solution

**Concepts Involved:** Stinespring Representation, Partial Trace

This problem involves deriving the operator-sum representation for the map  $\mathcal{E}(\rho) = \rho'$ , where  $\rho'$  is the state of subsystems BC after discarding subsystems AD.

Step 1: State evolution via unitary U

Initially, the total state of the system ABCD is

$$\rho_{AB} \otimes |0\rangle_C \langle 0| \otimes |0\rangle_D \langle 0|.$$

After applying the unitary U, the state becomes

$$U(\rho_{AB} \otimes |0\rangle_C \langle 0| \otimes |0\rangle_D \langle 0|) U^{\dagger}.$$

Step 2: Partial trace over AD

After the interaction, systems A and D are discarded. To describe the remaining state of systems BC, we need to trace out the subsystems AD. The state of BC after tracing out AD is

$$\rho_{BC}' = \operatorname{Tr}_{AD} \left( U(\rho_{AB} \otimes |0\rangle_C \langle 0| \otimes |0\rangle_D \langle 0|) U^{\dagger} \right).$$

# Step 3: Kraus decomposition

Now, we can express this map  ${\mathcal E}$  in terms of Kraus operators. To this end, we express

$$\mathcal{E}(\rho_{AB}) = \sum_{k} E_{k} \rho_{AB} E_{k}^{\dagger},$$

where the operators  $E_k$  are linear operators from the state space of system AB to the state space of system BC.

Specifically, the Kraus operators  $E_k$  are given by:

$$E_k = \langle e_k | U | 0 \rangle_D$$

where  $\{|e_k\rangle\}$  is some orthonormal basis for the combined system AD.

Step 4: Trace-preserving condition

Finally, the requirement that  ${\mathcal E}$  is trace-preserving implies that the Kraus operators must satisfy

$$\sum_{k} E_{k}^{\dagger} E_{k} = I.$$

This follows simply from the cyclicity of the trace.

# **Exercise 8.4: Measurement**

Suppose we have a single qubit principal system, interacting with a single qubit environment through the transform

$$U = P_0 \otimes I + P_1 \otimes X$$

where X is the usual Pauli matrix (acting on the environment), and  $P_0 \equiv |0\rangle\langle 0|$ ,  $P_1 \equiv |1\rangle\langle 1|$  are projectors (acting on the system). Give the quantum operation for this process, in the operator-sum representation, assuming the environment starts in the state  $|0\rangle$ .

# Solution

Concepts Involved: Kraus Representation, Measurements, Partial Trace

Simply using the natural form of the Kraus operators  $E_k \equiv \langle e_k | U | e_0 \rangle$ , we have:

$$E_0 = \langle 0_E | U | 0_E \rangle = P_0,$$

$$E_1 = \langle 1_E | U | 0_E \rangle = P_1.$$

# Exercise 8.5: Spin flips

Just as in the previous exercise, but now let

$$U = \frac{X}{\sqrt{2}} \otimes I + \frac{Y}{\sqrt{2}} \otimes X$$

Give the quantum operation for this process, in the operator sum representation.

# Solution

Concepts Involved: Kraus Represenations, Partial Trace

Using the same idea,

$$E_0 = \langle 0_E | U | 0_E \rangle = \frac{X}{\sqrt{2}},$$

$$E_1 = \langle 1_E | U | 0_E \rangle = \frac{Y}{\sqrt{2}}.$$

Exercise 8.6: Composition of quantum operations

Suppose  $\mathcal E$  and  $\mathcal F$  are quantum operations on the same quantum system. Show that the composition  $\mathcal F\circ\mathcal E$  is a quantum operation, in the sense that is has an operator-sum representation. State and prove an extension of this result to the case where  $\mathcal E$  and  $\mathcal F$  do not necessarily have the same input and output spaces.

# Solution

Concepts Involved: Krauss Representation, Linear Maps, Compositions

# Step 1: Composition of Two Quantum Operations

For an input state  $\rho$ , the action of the composition is:

$$(\mathcal{F} \circ \mathcal{E})(\rho) = \mathcal{F}(\mathcal{E}(\rho)).$$

First, apply  $\mathcal{E}(\rho)$  using its Kraus decomposition:

$$\mathcal{E}(\rho) = \sum_{k} E_k \rho E_k^{\dagger}.$$

Next, apply  $\mathcal{F}$  to the result of  $\mathcal{E}(\rho)$ :

$$\mathcal{F}(\mathcal{E}(
ho)) = \mathcal{F}\left(\sum_k E_k 
ho E_k^{\dagger}\right) = \sum_k \mathcal{F}(E_k 
ho E_k^{\dagger}).$$

Since  $\mathcal F$  is linear, we can apply it to each term in the sum separately:  $\mathcal F(E_k \rho E_k^\dagger) = \sum_j F_j E_k \rho E_k^\dagger F_j^\dagger$ . Thus, the composition  $\mathcal F \circ \mathcal E$  can be written as:

$$(\mathcal{F} \circ \mathcal{E})(\rho) = \sum_{j,k} F_j E_k \rho E_k^{\dagger} F_j^{\dagger}.$$

This expression is again in the form of an operator-sum representation, where the Kraus operators for the composed operation  $\mathcal{F} \circ \mathcal{E}$  are given by the products  $F_j E_k$ . Hence, we have shown that the composition of two quantum operations is itself a quantum operation, with Kraus operators  $F_j E_k$ .

# Step 2: Trace-Preserving Condition

To check that the map  $\mathcal{F} \circ \mathcal{E}$  is trace-preserving, we compute:

$$\sum_{j,k} (F_j E_k)^{\dagger} (F_j E_k) = \sum_{j,k} E_k^{\dagger} F_j^{\dagger} F_j E_k.$$

Since  $\mathcal F$  is trace-preserving, we know that  $\sum_i F_i^\dagger F_j = I$ , so this simplifies to:

$$\sum_{k} E_k^{\dagger} I E_k = \sum_{k} E_k^{\dagger} E_k = I,$$

# Step 3: Extension to Different Input and Output Spaces

Suppose that:

- ullet  $\mathcal E$  maps states from the Hilbert space  $\mathcal H_A$  to  $\mathcal H_B$ ,
- $\mathcal{F}$  maps states from  $\mathcal{H}_B$  to  $\mathcal{H}_C$ .

We can still write the composition  $\mathcal{F} \circ \mathcal{E}$ , but now  $\mathcal{E}$  will have Kraus operators  $E_k : \mathcal{H}_A \to \mathcal{H}_B$  and  $\mathcal{F}$  will have Kraus operators  $F_i : \mathcal{H}_B \to \mathcal{H}_C$ .

The composition will act as:

$$(\mathcal{F} \circ \mathcal{E})(\rho) = \sum_{j,k} F_j E_k \rho E_k^{\dagger} F_j^{\dagger},$$

where  $\rho$  is a state on  $\mathcal{H}_A$ , and the result will be a state on  $\mathcal{H}_C$ . The Kraus operators for the composed map are still given by the products  $F_j E_k$ , which now act from  $\mathcal{H}_A$  to  $\mathcal{H}_C$ .

#### Exercise 8.7

Suppose that instead of doing a projective measurement on the combined principal system and environment we had performed a general measurement described by measurement operators  $\{M_m\}$ . Find operator-sum representations for the corresponding quantum operations  $\mathcal{E}_m$  on the principal system, and show that the respective measurement probabilities are  $\operatorname{tr}[\mathcal{E}(\rho)]$ .

# Exercise 8.8: Non-trace-preserving quantum operations

Explain how to construct a unitary operator for a system–environment model of a non-trace-preserving quantum operation, by introducing an extra operator,  $E_{\infty}$ , into the set of operation elements  $E_k$ , chosen so that when summing over the complete set of k, including  $k=\infty$ , one obtains  $\sum_k E_k^{\dagger} E_k = I$ .

# Exercise 8.9: Measurement model

If we are given a set of quantum operations  $\{\mathcal{E}_m\}$  such that  $\sum_m \mathcal{E}_m$  is trace-preserving, then it is possible to construct a *measurement model* giving rise to this set of quantum operations. For each m, let  $E_{mk}$  be a set of operation elements for  $\mathcal{E}_m$ . Introduce an environmental system, E, with an orthonormal basis  $|m,k\rangle$  in one-to-one correspondence with the set of indices for the operation elements. Analogously to the earlier construction, define an operator U such that

$$U|\psi\rangle|e_0\rangle = \sum_{m,k} E_{mk}|\psi\rangle|m,k\rangle.$$

Next, define projectors  $P_m \equiv \sum_k |m,k\rangle\langle m,k|$  on the environmental system, E. Show that performing U on  $\rho \otimes |e_0\rangle\langle e_0|$ , then measuring  $P_m$  gives m with probability  $\operatorname{tr}(\mathcal{E}_m(\rho))$ , and the corresponding post-measurement state of the principal system is  $\mathcal{E}_m(\rho)/\operatorname{tr}(\mathcal{E}_m(\rho))$ .

# Exercise 8.10

Give a proof of Theorem 8.3 based on the freedom in the operator-sum representation, as follows. Let  $\left\{E_j\right\}$  be a set of operation elements for  $\mathcal{E}$ . Define a matrix  $W_{jk} \equiv \operatorname{tr}\left(E_j^{\dagger}E_k\right)$ . Show that the matrix W is Hermitian and of rank at most  $d^2$ , and thus there is a unitary matrix u such that  $uWu^{\dagger}$  is diagonal with at most  $d^2$  non-zero entries. Use u to define a new set of at most  $d^2$  non-zero operation elements  $\left\{F_i\right\}$  for  $\mathcal{E}$ .

#### Exercise 8.11

Suppose  $\mathcal{E}$  is a quantum operation mapping a d-dimensional input space to a d'-dimensional output space. Show that  $\mathcal{E}$  can be described using a set of at most dd' operation elements  $\{E_k\}$ .

# Exercise 8.12

Why can we assume that O has determinant 1 in the decomposition (8.93)?

# Exercise 8.13

Show that unitary transformations correspond to rotations of the Bloch sphere.

#### Exercise 8.14

Show that det(S) need not be positive.

# Exercise 8.15

Suppose a projective measurement is performed on a single qubit in the basis  $|+\rangle\,,\,|-\rangle$ , where  $|\pm\rangle\equiv(|0\rangle\pm|1\rangle)/\sqrt{2}$ . In the event that we are ignorant of the result of the measurement, the density matrix evolves according to the equation

$$\rho \mapsto \mathcal{E}(\rho) = |+\rangle \langle +|\rho| + |+\rangle \langle +|+|-\rangle \langle -|\rho| - |+\rangle \langle -|$$

Illustrate this transformation on the Bloch sphere.

# Exercise 8.16

The graphical method for understanding single qubit quantum operations was derived for trace-preserving quantum operations. Find an explicit example of a non-trace-preserving quantum operation which cannot be described as a deformation of the Bloch sphere, followed by a rotation and a displacement.

#### Exercise 8.17

Verify (8.101) as follows. Define

$$\mathcal{E}(A) \equiv \frac{A + XAX + YAY + ZAZ}{4}$$

and show that

$$\mathcal{E}(I) = I$$
;  $\mathcal{E}(X) = \mathcal{E}(Y) = \mathcal{E}(Z) = 0$ 

Now use the Bloch sphere representation for single qubit density matrices to verify (8.101).

# Solution

#### **Concepts Involved:**

It is straightforward to check

$$\mathcal{E}(X) \equiv \frac{X + XXX + YXY + ZXZ}{4} = \frac{X + X + YXY + ZXZ}{4} = \frac{X + X - X - X}{4} = 0$$

Similarly, it follows  $\mathcal{E}(Y) = \mathcal{E}(Z) = 0$ . Now, we simply employ linearity and the single qubit density matrix representation to arrive at

$$\mathcal{E}(\rho) = \mathcal{E}\left(\frac{I + r.\sigma}{2}\right) = \mathcal{E}\left(\frac{I}{2}\right) + 0$$

$$\implies \frac{\rho + X\rho X + Y\rho Y + Z\rho Z}{4} = \frac{I}{2}$$

which we set out to prove.

# Exercise 8.18

For  $k \ge 1$  show that  $\operatorname{tr}(\rho^k)$  is never increased by the action of the depolarizing channel.

# Solution

Concepts Involved: Trace, Depolarizing, Binomial expansion

$$\mathcal{E}_p(\rho) = (1-p)\rho + p\frac{I}{d}$$

$$[\mathcal{E}_p(\rho)]^k = \left((1-p)\rho + p\frac{I}{d}\right)^k$$

Now taking the trace, we have

$$\begin{split} \operatorname{Tr}\left([\mathcal{E}_p(\rho)]^k\right) &= \sum_{m=0}^k \binom{k}{m} (1-p)^{k-m} p^m \operatorname{Tr}\left(\rho^{k-m}\right) \left(\frac{I}{d}\right)^m \\ &\leq \operatorname{Tr}\left(\rho^k\right) \sum_{m=0}^k \binom{k}{m} (1-p)^{k-m} p^m \left(\frac{1}{d}\right)^m \qquad \text{as } \operatorname{Tr}(\rho^k) \geq \operatorname{Tr}(\rho^{k-m}) \ \forall m \geq 0 \\ &= \operatorname{Tr}\left(\rho^k\right) (1-p+\frac{p}{d})^k \\ &\leq \operatorname{Tr}\left(\rho^k\right) (1-p+p)^k \qquad \qquad \forall \, d \geq 1 \\ &= \operatorname{Tr}\left(\rho^k\right) \end{split}$$

# Exercise 8.19

Find an operator-sum representation for a generalized depolarizing channel acting on a d-dimensional Hilbert space.

# Solution

#### **Concepts Involved:**

The generalized depolarizing channel  $\mathcal{E}$  on a d-dimensional Hilbert space is defined as:

$$\mathcal{E}(\rho) = (1 - p)\rho + \frac{p}{d}I,$$

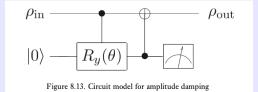
where  $0 \le p \le 1$  is the depolarizing probability and I is the identity operator in the d-dimensional space.

The operator-sum (Kraus) representation for this channel is:

$$E_0 = \sqrt{1-p} I$$
 
$$E_{i,j} = \sqrt{\frac{p}{d}} |i\rangle\langle j| \qquad i, j = 1, \dots, d.$$

# Exercise 8.20

Show that the circuit in Figure 8.13 models the amplitude damping quantum operation, with  $\sin^2(\theta/2) = \gamma$ .



# Solution

# Concepts Involved:

It is instructive to write the unitary interaction as

$$U = (I \otimes P_0 + X \otimes P_1)(P_0 \otimes I + P_1 \otimes R_y(\theta))$$
  
=  $P_0 \otimes P_0 + P_1 \otimes P_0 R_y(\theta) + X P_0 \otimes P_1 + X P_1 \otimes P_1 R_y(\theta)$ 

Thus, the Kraus operators can be calculated using

$$E_0 \equiv \langle 0_E | U | 0_E \rangle = P_0 \langle 0_E | P_0 | 0_E \rangle + P_1 \langle 0_E | P_0 R_y(\theta) | 0_E \rangle = P_0 + P_1 \cdot \cos(\theta/2)$$
  

$$E_1 \equiv \langle 1_E | U | 0_E \rangle = X P_1 \langle 1_E | P_1 R_y(\theta) | 0_E \rangle = X P_1 \cdot \sin(\theta/2)$$

To match this with the matrix representation in Eq. (8.108), we should have

$$\cos(\theta/2) = \sqrt{1 - \gamma}$$

$$\implies \gamma = \sin^2(\theta/2).$$

### Exercise 8.21: Amplitude damping of a harmonic oscillator

Suppose that our principal system, a harmonic oscillator, interacts with an environment, modeled as another harmonic oscillator, through the Hamiltonian

$$H = \chi(a^{\dagger}b + b^{\dagger}a)$$

where a and b are the annihilation operators for the respective harmonic oscillators, as defined in Section 7.3.

(1) Using  $U=\exp(-iH\Delta t)$ , denoting the eigenstates of  $b^{\dagger}b$  at  $|k_b\rangle$ , and selecting the vacuum state  $|0_b\rangle$  as the initial state of the environment, show that the operation elements  $E_k=\langle k_b|U|0_b\rangle$  are found to be

$$E_k = \sum_{n} \sqrt{\binom{n}{k}} \sqrt{(1-\gamma)^{n-k} \gamma^k} |n-k\rangle\langle n|$$

where  $\gamma=1-\cos^2(\chi\Delta t)$  is the probability of losing a single quantum of energy, and states such as  $|n\rangle$  are eigenstates of  $a^\dagger a$ .

(2) Show that the operation elements  $E_k$  define a trace-preserving quantum operation.

# Exercise 8.22: Amplitude damping of single qubit density matrix

For the general single qubit state

$$\rho = \begin{bmatrix} a & b \\ b^* & c \end{bmatrix}$$

show that the amplitude damping leads to

$$\mathcal{E}_{AD}(\rho) = \begin{bmatrix} 1 - (1 - \gamma)(1 - a) & b\sqrt{1 - \gamma} \\ b^*\sqrt{1 - \gamma} & c(1 - \gamma) \end{bmatrix}$$

#### Solution

### Concepts Involved:

For the general single qubit state

$$\rho = \begin{pmatrix} a & b \\ b^* & c \end{pmatrix}, \quad a + c = 1,$$

the amplitude damping channel is

$$\mathcal{E}_{AD}(\rho) = K_0 \rho K_0^{\dagger} + K_1 \rho K_1^{\dagger},$$

where the Kraus operators are

$$K_0 = \begin{pmatrix} 1 & 0 \\ 0 & \sqrt{1-\gamma} \end{pmatrix}, \quad K_1 = \begin{pmatrix} 0 & \sqrt{\gamma} \\ 0 & 0 \end{pmatrix}.$$

Now, applying the Kraus operators

$$K_0 \rho K_0^{\dagger} = \begin{pmatrix} a & b\sqrt{1-\gamma} \\ b^*\sqrt{1-\gamma} & c(1-\gamma) \end{pmatrix},$$
  
$$K_1 \rho K_1^{\dagger} = \begin{pmatrix} \gamma c & 0 \\ 0 & 0 \end{pmatrix}.$$

Summing the two results

$$\mathcal{E}_{AD}(\rho) = \begin{pmatrix} a + \gamma(1-a) & b\sqrt{1-\gamma} \\ b^*\sqrt{1-\gamma} & c(1-\gamma) \end{pmatrix}.$$

Since a+c=1, we arrive at the desired form

$$\mathcal{E}_{AD}(\rho) = \begin{pmatrix} 1 - (1 - \gamma)(1 - a) & b\sqrt{1 - \gamma} \\ b^*\sqrt{1 - \gamma} & c(1 - \gamma) \end{pmatrix}.$$

Exercise 8.23: Amplitude damping of dual-rail qubits

Suppose that a single qubit state is represented by using two qubits, as

$$|\psi\rangle = a |01\rangle + b |10\rangle$$
.

Show that  $\mathcal{E}_{AD}\otimes\mathcal{E}_{AD}$  applied to this state gives a process which can be described by the operation elements

$$\begin{split} E_0^{\rm dr} &= \sqrt{1-\gamma} \; I \\ E_1^{\rm dr} &= \sqrt{\gamma} \left[ |00\rangle\langle 01| + |00\rangle\langle 10| \right] \end{split}$$

that is, either nothing  $(E_0^{\mathrm{dr}})$  happens to the qubit, or the qubit is transformed  $(E_1^{\mathrm{dr}})$  into the state  $|00\rangle$ , which is orthogonal to  $|\psi\rangle$ . This is a simple error-detection code, and is also the basis for the robustness of the 'dual-rail' qubit discussed in Section 7.4.

Solution

**Concepts Involved:** 

The amplitude damping channel  $\mathcal{E}_{AD}$  is described by the Kraus operators

$$K_0 = \begin{pmatrix} 1 & 0 \\ 0 & \sqrt{1-\gamma} \end{pmatrix}, \qquad K_1 = \begin{pmatrix} 0 & \sqrt{\gamma} \\ 0 & 0 \end{pmatrix}.$$

The action of  $\mathcal{E}_{AD}\otimes\mathcal{E}_{AD}$  on  $|\psi\rangle$  results in

$$K_0 \otimes K_0 |\psi\rangle = \sqrt{1 - \gamma} |\psi\rangle,$$
  

$$K_1 \otimes K_0 |\psi\rangle = a\sqrt{\gamma} |00\rangle,$$
  

$$K_0 \otimes K_1 |\psi\rangle = b\sqrt{\gamma} |00\rangle,$$
  

$$K_1 \otimes K_1 |\psi\rangle = 0.$$

Thus, the total channel action is given by

$$\mathcal{E}_{AD} \otimes \mathcal{E}_{AD}(|\psi\rangle\langle\psi|) = (1 - \gamma) |\psi\rangle\langle\psi| + |a|^2 \gamma |00\rangle\langle00| + |b|^2 \gamma |00\rangle\langle00|$$
$$= (1 - \gamma) |\psi\rangle\langle\psi| + \gamma |00\rangle\langle00|$$

This process can be described by the operation elements

$$\begin{split} E_0^{\mathrm{dr}} &= \sqrt{1 - \gamma} I, \\ E_1^{\mathrm{dr}} &= \sqrt{\gamma} \left( \left| 00 \right\rangle \left\langle 01 \right| + \left| 00 \right\rangle \left\langle 10 \right| \right). \end{split}$$

Hence,  $\mathcal{E}_{AD}\otimes\mathcal{E}_{AD}$  either applies no change to the qubits  $(E_0^{\mathrm{dr}})$  or transforms the state into  $|00\rangle$   $(E_1^{\mathrm{dr}})$ .

# Exercise 8.24: Spontaneous emission is amplitude damping

A single atom coupled to a single mode of electromagnetic radiation undergoes spontaneous emission, as was described in Section 7.6.1. To see that this process is just amplitude damping, take the unitary operation resulting from the Jaynes–Cummings interaction, Equation (7.77), with detuning  $\delta=0$ , and give the quantum operation resulting from tracing over the field.

# Exercise 8.25

If we define the temperature T of a qubit by assuming that in equilibrium the probabilities of being in the  $|0\rangle$  or  $|1\rangle$  states satisfy a Boltzmann distribution, that is  $r_0=e^{-E_0/k_Bt}/\mathcal{Z}$  and  $p_1=e^{-E_1/k_BT}/\mathcal{Z}$ , where  $E_0$  is the energy of the state  $|0\rangle$ ,  $E_1$  the energy of the state  $|1\rangle$ , and  $\mathcal{Z}=e^{-E_0/k_Bt}+e^{-E_1/k_BT}$ , what temperature described the state  $\rho_\infty$ ?

# Exercise 8.26: Circuit model for phase damping

Show that the circuit in Figure 8.15 can be used to model the phase damping quantum operation, provided  $\theta$  is chosen appropriately.

### Exercise 8.27: Phase damping = phase flip channel

Give the unitary transformation which related the operation elements of (8.127)-(8.128) to those of (8.129)-(8.130); that is, find u such that  $\tilde{E}_k = \sum_j u_{kj} E_j$ .

# Exercise 8.28: One CNOT phase damping model circuit

Show that a single controlled-CNOT gate can be used as a model for phase damping, if we let the initial state of the environment be a mixed state, where the amount of damping is determined by the probability of the states in the mixture.

# **Exercise 8.29: Unitality**

A quantum process  $\mathcal{E}$  is *unital* if  $\mathcal{E}(I) = I$ . Show that the depolarizing and phase damping channels are unital, while amplitude damping is not.

# **Exercise 8.30:** $T_2 \le T_1/2$

The  $T_2$  phase coherence relaxation rate is just the exponential decay rate of the off-diagonal elements in the qubit density matrix, while  $T_1$  is the decay rate of the diagonal elements (see Equation (7.144)). Amplitude damping has both nonzero  $T_1$  and  $T_2$  rates; show that for amplitude damping  $T_2 = T_1/2$ . Also show that if amplitude and phase damping are both applied then  $T_2 \leq T_1/2$ .

# Exercise 8.31: Exponential sensitivity to phase damping

Using (8.126), show that the element  $\rho_{nm}=\langle n|\rho|m\rangle$  in the density matrix of the harmonic oscillator decays exponentially as  $e^{-\lambda(n-m)^2}$  under the effect of phase damping, for some constant  $\lambda$ .

### Exercise 8.32

Explain how to extend quantum process tomography to the case of non-trace-preserving quantum operations, such as arise in the study of measurement.

#### Exercise 8.33: Specifying a quantum process

Suppose that one wished to completely specify an arbitrary single qubit operation  $\mathcal{E}$  by describing how a set of points on the Bloch sphere  $\{\mathbf{r}_k\}$  transform under  $\mathcal{E}$ . Prove that the set must contain at least four points.

### Exercise 8.34: Process tomography for two qubits

Show that the  $\chi_2$  describing the black box operations on two qubits can be expressed as

$$\chi_2 = \Lambda_2 \bar{\rho}' \Lambda_2$$

where  $\Lambda_2 = \Lambda \otimes \Lambda$ ,  $\Lambda$  is as defined in Box 8.5, and  $\bar{\rho}'$  is a block matrix of 16 measured density matrices,

$$\bar{\rho}' = P^T \begin{bmatrix} \rho'_{11} & \rho'_{12} & \rho'_{13} & \rho'_{14} \\ \rho'_{21} & \rho'_{22} & \rho'_{23} & \rho'_{24} \\ \rho'_{31} & \rho'_{32} & \rho'_{33} & \rho'_{34} \\ \rho'_{41} & \rho'_{42} & \rho'_{43} & \rho'_{44} \end{bmatrix} P$$

where  $\rho'_{nm} = \mathcal{E}\left(\rho_{nm}\right)$ ,  $\rho_{nm} = T_n \left|00\right\rangle\!\left\langle00\right| T_m$ ,  $T_1 = I \otimes I$ ,  $T_2 = I \otimes X$ ,  $T_3 = X \otimes I$ ,  $T_4 = X \otimes X$ , and  $P = I \otimes \left[\left(\rho_{00} + \rho_{12} + \rho_{21} + \rho_{33}\right) \otimes I\right]$  is a permutation matrix.

# Exercise 8.35: Process tomography example

Consider a one qubit black box of unknown dynamics  $\mathcal{E}_1$ . Suppose that the following four density matrices are obtained from experimental measurements, performed according to Equations (8.173)–(8.176):

$$\rho_1' = \begin{bmatrix} 1 & 0 \\ 0 & 0 \end{bmatrix}$$

$$\rho_2' = \begin{bmatrix} 0 & \sqrt{1-\gamma} \\ 0 & 0 \end{bmatrix}$$

$$\rho_3' = \begin{bmatrix} 0 & 0 \\ \sqrt{1-\gamma} & 0 \end{bmatrix}$$

$$\rho_4' = \begin{bmatrix} \gamma & 0 \\ 0 & 1-\gamma \end{bmatrix}$$

where  $\gamma$  is a numerical parameter. From an independent study of each of these input-output relations, one could make several important observations: the ground state  $|0\rangle$  is left invariant by  $\mathcal{E}_1$ . the excited state  $|1\rangle$  partially decays to the ground state, and superposition states are damped. Determine the  $\chi$  matrix for this process.

#### Problem 8.1: Lindblad form to quantum operation

In the notation of Section 8.4.1, explicitly work through the steps to solve the differential equation

$$\dot{\rho} = -\frac{\lambda}{2}(\sigma_{+}\sigma_{-}\rho + \rho\sigma_{+}\sigma_{-} - 2\sigma_{-}\rho\sigma_{+})$$

for  $\rho(t)$ . Express the map  $\rho(0)\mapsto \rho(t)$  as  $\rho(t)=\sum_k E_k(t)\rho(0)E_k^\dagger(t)$ .

#### Solution

### Concepts Involved: Differential Equations, Kraus Represenations

Let us use the master equation to obtain the differential equations for the diagonal and off-diagonal elements: - \*\*Population (diagonal) elements\*\*:

$$\dot{\rho}_{11}(t) = -\lambda \rho_{11}(t), \quad \dot{\rho}_{00}(t) = \lambda \rho_{11}(t),$$

with solutions:

$$\rho_{11}(t) = \rho_{11}(0)e^{-\lambda t}, \quad \rho_{00}(t) = \rho_{00}(0) + \rho_{11}(0)(1 - e^{-\lambda t}).$$

- \*\*Coherence (off-diagonal) elements\*\*:

$$\dot{\rho}_{10}(t) = -\frac{\lambda}{2}\rho_{10}(t), \quad \dot{\rho}_{01}(t) = -\frac{\lambda}{2}\rho_{01}(t),$$

with solutions:

$$\rho_{10}(t) = \rho_{10}(0)e^{-\lambda t/2}, \quad \rho_{01}(t) = \rho_{01}(0)e^{-\lambda t/2}.$$

The time evolution of the density matrix can be expressed using Kraus operators for amplitude damping:

$$E_0(t) = \begin{pmatrix} 1 & 0 \\ 0 & \sqrt{e^{-\lambda t}} \end{pmatrix}, \quad E_1(t) = \begin{pmatrix} 0 & \sqrt{1 - e^{-\lambda t}} \\ 0 & 0 \end{pmatrix}.$$

These Kraus operators satisfy the completeness relation:

$$E_0^{\dagger}(t)E_0(t) + E_1^{\dagger}(t)E_1(t) = I,$$

and the solution to the master equation is:

$$\rho(t) = E_0(t)\rho(0)E_0^{\dagger}(t) + E_1(t)\rho(0)E_1^{\dagger}(t).$$

#### Problem 8.2: Teleportation as a quantum operation

Suppose Alice is in possession of a single qubit, denoted as system 1, which she wishes to teleport to Bob. Unfortunately, she and Bob only share an imperfectly entangled pair of qubits. Alice's half of this pair is denoted system 2, and Bob's half is denoted system 3. Suppose Alice performs a measurement described by a set of quantum operations  $\mathcal{E}_m$  with result m on systems 1 and 2. Show that this induces an operation  $\tilde{E}_m$  relating the initial stat of the system 1 to the final state of system 3, and that teleportation is accomplished if Bob can reverse this operation using a trace-preserving quantum operation  $\mathcal{R}_m$ , to

obtain

$$\mathcal{R}_m \left( \frac{\tilde{\mathcal{E}}_m(\rho)}{\operatorname{tr} \left[ \tilde{\mathcal{E}}_m(\rho) \right]} \right) = \rho,$$

where  $\rho$  is the initial state of system 1.

#### Solution

#### Concepts Involved:

# Teleportation Protocol

Let us consider a scenario where Alice holds qubit 1 (state  $\rho$ ) to be teleported, and shares an imperfectly entangled pair of qubits with Bob, where Alice holds qubit 2 and Bob holds qubit 3. The initial state of the combined system (qubits 1, 2, and 3) is:

$$\rho_{123} = \rho_1 \otimes |\Psi\rangle_{23} \langle \Psi|,$$

where  $|\Psi\rangle_{23}$  represents the (imperfectly) entangled state between qubits 2 and 3.

# Step 1: Alice's Measurement on Systems 1 and 2

Alice performs a measurement on systems 1 and 2, described by the quantum operation  $\mathcal{E}_m$ , corresponding to Kraus operators  $E_m^{(12)}$ , and records the measurement outcome m. After this measurement, the state of the combined system becomes:

$$\tilde{\rho}_{123}^{(m)} = \frac{(E_m^{(12)} \otimes I_3) \rho_{123} (E_m^{(12)} \otimes I_3)^\dagger}{\mathrm{Tr}[(E_m^{(12)} \otimes I_3) \rho_{123} (E_m^{(12)} \otimes I_3)^\dagger]}.$$

# Step 2: Induced Operation on Bob's Qubit

The measurement on systems 1 and 2 induces an operation  $\tilde{\mathcal{E}}_m$  on Bob's qubit (system 3), which can be expressed by tracing out systems 1 and 2:

$$\tilde{\mathcal{E}}_m(\rho_1) = \mathrm{Tr}_{12} \left[ (E_m^{(12)} \otimes I_3) (\rho_1 \otimes |\Psi\rangle_{23} \, \langle \Psi|) (E_m^{(12)} \otimes I_3)^\dagger \right].$$

This describes how Alice's measurement modifies the state of Bob's qubit, based on the initial state  $\rho_1$  of system 1.

# Step 3: Bob's Recovery Operation

Once Alice communicates the measurement outcome m to Bob, he applies a trace-preserving quantum operation  $\mathcal{R}_m$  to reverse the effect of the induced operation  $\tilde{\mathcal{E}}_m$ . The condition for successful teleportation

is:

$$\mathcal{R}_m\left(rac{ ilde{\mathcal{E}}_m(
ho_1)}{\mathsf{Tr}[ ilde{\mathcal{E}}_m(
ho_1)]}
ight)=
ho_1,$$

where the division by  ${\rm Tr}[\tilde{\mathcal{E}}_m(
ho_1)]$  ensures normalization.

# **Problem 8.3: Random unitary channels**

It is tempting to believe that all unital channels, that is, those for which  $\mathcal{E}(I)=I$ , result from averaging over random unitary operations, that is,  $\mathcal{E}(\rho)=\sum_k p_k\;U_k\;\rho\,U_k^\dagger$ , where  $U_k$  are unitary operations and the  $p_k$  form a probability distribution. Show that while this is true for single qubits, it is untrue for larger systems.

#### Solution

#### **Concepts Involved:**

Let  $\mathcal{X} = \mathbb{C}^3$  and define the (Werner-Holevo) map

$$\mathcal{E}(X) = \frac{1}{2} \mathrm{Tr}(X) I - \frac{1}{2} X^T$$

for all  $X \in L(\mathcal{X})$ . It can be confirmed that  $\mathcal{E}$  is a channel and unital, but it is not a mixed-unitary channel. To demonstrate that  $\mathcal{E}$  is not a mixed-unitary channel, observe the decomposition

$$\mathcal{E}(X) = A_1 X A_1^{\dagger} + A_2 X A_2^{\dagger} + A_3 X A_3^{\dagger}$$

for all  $X \in L(\mathcal{X})$ , where the matrices are given as

$$A_1 = \begin{pmatrix} 0 & 0 & 0 \\ 0 & 0 & \frac{1}{\sqrt{2}} \\ 0 & -\frac{1}{\sqrt{2}} & 0 \end{pmatrix}, \quad A_2 = \begin{pmatrix} 0 & 0 & \frac{1}{\sqrt{2}} \\ 0 & 0 & 0 \\ -\frac{1}{\sqrt{2}} & 0 & 0 \end{pmatrix}, \quad A_3 = \begin{pmatrix} 0 & \frac{1}{\sqrt{2}} & 0 \\ -\frac{1}{\sqrt{2}} & 0 & 0 \\ 0 & 0 & 0 \end{pmatrix}.$$

The validity of this expression for all  $X \in L(\mathcal{X})$  can be verified through the Choi representation of the map, which matches the right-hand side.

Recalling the Choi theorem.

Let  $\mathcal{X}$  and  $\mathcal{Y}$  be complex Euclidean spaces. Suppose  $\mathcal{E} \in C(\mathcal{X}, \mathcal{Y})$  is a quantum channel, and let  $\{A_a : a \in \Sigma\} \subset L(\mathcal{X}, \mathcal{Y})$  be a linearly independent set of operators. Then

$$\mathcal{E}(X) = \sum_{a \in \Sigma} A_a X A_a^{\dagger}$$

for all  $X \in L(\mathcal{X})$ . The channel  $\mathcal{E}$  is an extreme point of the set of channels  $C(\mathcal{X}, \mathcal{Y})$  if and only if the collection

$$\{A_b^{\dagger}A_a:(a,b)\in\Sigma\times\Sigma\}\subset L(\mathcal{X})$$

is linearly independent.

For us, the set  $\{A_i^{\dagger}A_k: 1 \leq j, k \leq 3\}$  contains the following operators:

$$\begin{split} A_1^\dagger A_1 &= \begin{pmatrix} 0 & 0 & 0 \\ 0 & \frac{1}{2} & 0 \\ 0 & 0 & \frac{1}{2} \end{pmatrix}, \quad A_1^\dagger A_2 = \begin{pmatrix} 0 & 0 & 0 \\ \frac{1}{2} & 0 & 0 \\ 0 & 0 & 0 \end{pmatrix}, \quad A_1^\dagger A_3 = \begin{pmatrix} 0 & 0 & 0 \\ 0 & 0 & 0 \\ 0 & -\frac{1}{2} & 0 \end{pmatrix}, \\ A_2^\dagger A_1 &= \begin{pmatrix} 0 & \frac{1}{2} & 0 \\ 0 & 0 & 0 \\ 0 & 0 & 0 \end{pmatrix}, \quad A_2^\dagger A_2 = \begin{pmatrix} 0 & 0 & 0 \\ 0 & 0 & 0 \\ 0 & 0 & \frac{1}{2} \end{pmatrix}, \quad A_2^\dagger A_3 = \begin{pmatrix} 0 & 0 & 0 \\ 0 & 0 & \frac{1}{2} \\ 0 & 0 & 0 \end{pmatrix}, \\ A_3^\dagger A_1 &= \begin{pmatrix} 0 & 0 & -\frac{1}{2} \\ 0 & 0 & 0 \\ 0 & 0 & 0 \end{pmatrix}, \quad A_3^\dagger A_2 = \begin{pmatrix} 0 & 0 & 0 \\ 0 & 0 & \frac{1}{2} \\ 0 & 0 & 0 \end{pmatrix}, \quad A_3^\dagger A_3 = \begin{pmatrix} \frac{1}{2} & 0 & 0 \\ 0 & 0 & 0 \\ 0 & 0 & 0 \end{pmatrix}. \end{split}$$

This collection forms a linearly independent set and thus  $\mathcal{E}$  is an extreme point of the set of channels. Since  $\mathcal{E}$  is not a unitary channel, it follows that it cannot be written as a convex combination of unitary channels.

# 10 Quantum error-correction

#### Exercise 10.1

Verify that the encoding circuit in Figure 10.2 works as claimed.

$$\begin{array}{c|c} |\psi\rangle & & \\ |0\rangle & & \\ \\ |0\rangle & & \\ \end{array}$$

that is, it encodes  $|\psi\rangle = a|0\rangle + b|1\rangle$  to  $a|000\rangle + b|111\rangle$ .

#### Solution

Concepts Involved: Linear Algebra, Quantum Gates, Controlled Operations.

We calculate:

$$\begin{aligned} |\psi\rangle\otimes|0\rangle\otimes|0\rangle&\mapsto CX_{1,3}CX_{1,2}\,|\psi\rangle\otimes|0\rangle\otimes|0\rangle\\ &=aCX_{1,3}CX_{1,2}\,|000\rangle+bCX_{1,3}CX_{1,2}\,|100\rangle\\ &=a\,|000\rangle+b\,|111\rangle \end{aligned}$$

where we use the definition of the controlled-X gate  $CX_{1,2} |0\rangle \otimes |\varphi\rangle = |0\rangle \otimes |\varphi\rangle$  and  $CX_{1,2} |1\rangle \otimes |\varphi\rangle = |1\rangle \otimes X_2 |\varphi\rangle$  in the last line.

# Exercise 10.2

The action of the bit flip channel can be described by the quantum operation  $\mathcal{E}(\rho)=(1-p)\rho+pX\rho X$ . Show that this may be given an alternative operator-sum representation, as  $\mathcal{E}(\rho)=(1-2p)\rho+2pP_+\rho P_++2pP_-\rho P_-$  where  $P_+$  and  $P_-$  are projectors onto the +1 and -1 eigenstates of X,  $(|0\rangle+|1\rangle)/\sqrt{2}$  and  $(|0\rangle-|1\rangle)/\sqrt{2}$ , respectively. This latter representation can be understood as a model in which the qubit is left alone with probability 1-2p, and is 'measured' by the environment in the  $|+\rangle$ ,  $|-\rangle$  basis with probability 2p.

#### Solution

Concepts Involved: Linear Algebra, Projectors, Spectral Decomposition.

First, note that we may write:

$$P_{+} = |+\rangle\langle +| = \frac{|+\rangle\langle +| +|-\rangle\langle -| +|+\rangle\langle +| -|-\rangle\langle -|}{2} = \frac{I+X}{2}$$

where we have used  $X=|+\rangle\langle+|-|-\rangle\langle-|$  as the spectral decomposition of X and  $I=|+\rangle\langle+|+|-\rangle\langle-|$  the resolution of the identity. We can obtain the analogous relation for  $P_-$ :

$$P_{-} = \frac{I - X}{2}$$

Therefore we can expand:

$$\begin{split} \mathcal{E}(\rho) &= (1-2p)\rho + 2pP_{+}\rho P_{+} + 2pP_{-}\rho P_{-} \\ &= (1-2p)\rho + 2p\frac{I+X}{2}\rho\frac{I+X}{2} + 2p\frac{I-X}{2}\rho\frac{I-X}{2} \\ &= (1-2p)\rho + \frac{p}{2}(\rho + X\rho + \rho X + X\rho X) + \frac{p}{2}(\rho - X\rho - \rho X + X\rho X) \\ &= (1-p)\rho + pX\rho X \end{split}$$

which proves the claim.

# Exercise 10.3

Show by explicit calculation that measuring  $Z_1Z_2$  followed by  $Z_2Z_3$  is equivalent, up to labeling of the measurement outcomes, to measuring the four projectors defined by (10.5)–(10.8), in the sense that both procedures result in the same measurement statistics and post-measurement states.

# Solution

Concepts Involved: Linear Algebra, Projectors, Spectral Decomposition.

The four projection operators corresponding to the four error syndromes of the three-qubit repetition code are given by:

$$P_0 \equiv |000\rangle\langle000| + |111\rangle\langle111|$$

$$P_1 \equiv |100\rangle\langle100| + |011\rangle\langle011|$$

$$P_2 \equiv |010\rangle\langle010| + |101\rangle\langle101|$$

$$P_3 \equiv |001\rangle\langle001| + |110\rangle\langle110|$$

It suffices to show that the composition of projectors corresponding to measurements of  $Z_1Z_2$  and  $Z_2Z_3$  yield the same four projectors as the above.  $Z_1Z_2$  has spectral decomposition:

$$Z_1 Z_2 = (|00\rangle\langle 00| + |11\rangle\langle 11|) \otimes I - (|01\rangle\langle 01| + |10\rangle\langle 10|) \otimes I$$

which corresponds to a projective measurement with projectors:

$$P_{Z_1Z_2=+1} = (|00\rangle\langle 00| + |11\rangle\langle 11|) \otimes I$$
  
 $P_{Z_1Z_2=-1} = (|01\rangle\langle 01| + |10\rangle\langle 10|) \otimes I$ 

analogously,  $Z_2Z_3$  has spectral decomposition:

$$Z_2Z_3 = I \otimes (|00\rangle\langle 00| + |11\rangle\langle 11|) - I \otimes (|01\rangle\langle 01| + |10\rangle\langle 10|)$$

which corresponds to a projective measurement with projectors:

$$P_{Z_2 Z_3 = +1} = I \otimes (|00\rangle\langle 00| + |11\rangle\langle 11|)$$
  
$$P_{Z_2 Z_3 = -1} = I \otimes (|01\rangle\langle 01| + |10\rangle\langle 10|).$$

Now, we observe that:

$$P_{Z_1Z_2=-1}P_{Z_2Z_3=+1} = [(|01\rangle\langle 01| + |10\rangle\langle 10|) \otimes I][I \otimes (|00\rangle\langle 00| + |11\rangle\langle 11|)] = |011\rangle\langle 011| + |100\rangle\langle 100| = P_1 \otimes I_1 \otimes I_2 \otimes I_2 \otimes I_3 \otimes I_4 \otimes$$

$$P_{Z_1Z_2=+1}P_{Z_2Z_3=-1} = [(|00\rangle\langle 00| + |11\rangle\langle 11|) \otimes I][I \otimes (|01\rangle\langle 01| + |10\rangle\langle 10|)] = |001\rangle\langle 001| + |110\rangle\langle 110| = P_3$$

so the claim is proven.

#### Exercise 10.4

Consider the three qubit bit flip code. Suppose we had performed the error syndrome measurement by measuring the eight orthogonal projectors corresponding to projections onto the eight computational basis states.

- (1) Write out the projectors corresponding to this measurement, and explain how the measurement result can be used to diagnose the error syndrome: either no bits flipped or bit number j flipped, where j is in the range one to three.
- (2) Show that the recovery procedure works only for computational basis states.
- (3) What is the minimum fidelity for the error-correction procedure?

#### Solution

Concepts Involved: Linear Algebra, Projectors, Error Syndrome and Recovery, Fidelity

(1) The projectors are:

$$|000\rangle\langle000|$$
,  $|001\rangle\langle001|$ ,  $|010\rangle\langle010|$ ,  $|100\rangle\langle100|$ ,  $|011\rangle\langle011|$ ,  $|101\rangle\langle101|$ ,  $|110\rangle\langle110|$ ,  $|111\rangle\langle111|$ .

Since the encoded state is  $|\psi\rangle=a\,|000\rangle+b\,|111\rangle$ , if we measure one of  $|000\rangle\,, |111\rangle$  we would conclude no bits have been flipped. If we measure  $|001\rangle$  or  $|110\rangle$ , we would conclude that bit 1 was flipped (and so the state had become  $a\,|001\rangle+b\,|110\rangle$ . If we measure  $|010\rangle$  or  $|101\rangle$ , then we would conclude that bit 2 was flipped (and so the state had become  $a\,|010\rangle+b\,|101\rangle$ ). Finally, if we measure  $|100\rangle$  or  $|011\rangle$ , we would conclude that bit 3 was flipped (and so the state had become  $a\,|100\rangle+b\,|011\rangle$ ).

(2) The recovery procedure involves doing nothing if the error syndrome tells us that no bits are flipped, or flipping the jth bit if the jth bit was flipped. In the no bit flip case, after recovery we have  $|000\rangle \rightarrow |000\rangle, |111\rangle \rightarrow |111\rangle$ . In the case we conclude the first bit was flipped, after recovery we have  $|001\rangle \rightarrow |000\rangle, |110\rangle \rightarrow |111\rangle$ . In the case we conclude the second bit was flipped, after recovery we have  $|010\rangle \rightarrow |000\rangle, |101\rangle \rightarrow |111\rangle$ . Finally if we conclude that the third bit was flipped,

after recovery we have  $|100\rangle \to |000\rangle$ ,  $|011\rangle \to |111\rangle$ . In all cases, the post-recovery state is one of the computational basis states  $|000\rangle / |111\rangle$ , so the recovery only succeeds if the initial state was one of the computational basis states.

(3) Supposing we use the three qubit error-correcting code to protect  $|\psi\rangle=a\,|0\rangle+b\,|1\rangle$ . The encoded state is  $|\Psi\rangle=a\,|000\rangle+b\,|111\rangle$ . After applying the noise channel (i.e. each bit flips with probability p), the state we have is:

$$\mathcal{E}(\rho_{\Psi}) = \left(|a|^2 (1-p)^3 + |b|^2 p^3\right) |000\rangle\langle 000|$$

$$+ \left(|a|^2 p (1-p)^2 + |b|^2 p^2 (1-p)\right) \left(|001\rangle\langle 001| + |010\rangle\langle 010| + |100\rangle\langle 100|\right)$$

$$+ \left(|b|^2 p (1-p)^2 + |a|^2 p^2 (1-p)\right) \left(|110\rangle\langle 110| + |101\rangle\langle 101| + |011\rangle\langle 011|\right)$$

$$+ \left(|b|^2 (1-p)^3 + |a|^2 p^3\right) |111\rangle\langle 111|$$

The recovery procedure maps any state with  $\geq 2$  zeros to  $|000\rangle\langle000|$  (and hence back to  $|0\rangle\langle0|$  after decoding) and any state with  $\geq 2$  ones to  $|111\rangle\langle111|$  (and hence back to  $|1\rangle\langle1|$  after decoding), so the final state is:

$$\rho := \mathcal{R}(\mathcal{E}(\rho_{\Psi})) = \left( |a|^2 \left[ (1-p)^3 + 3p(1-p)^2 \right] + |b|^2 \left[ p^3 + 3p^2(1-p) \right] \right) |0\rangle\langle 0|$$
$$+ \left( |b|^2 \left[ (1-p)^3 + 3p(1-p)^2 \right] + |a|^2 \left[ p^3 + 3p^2(1-p) \right] \right) |1\rangle\langle 1|$$

To find the minimum fidelity, it suffices to minimize  $\langle \psi | \rho | \psi \rangle$ , which we compute to be:

$$\langle \psi | \rho | \psi \rangle = |a|^2 \left( |a|^2 \left[ (1-p)^3 + 3p(1-p)^2 \right] + |b|^2 \left[ p^3 + 3p^2 (1-p) \right] \right)$$
$$+ |b|^2 \left( |b|^2 \left[ (1-p)^3 + 3p(1-p)^2 \right] + |a|^2 \left[ p^3 + 3p^2 (1-p) \right] \right)$$

With the normalization constraint on the initial state of  $|a|^2 + |b|^2 = 1$ , we can rewrite the above in terms of of  $|a|^2$  alone:

$$\langle \psi | \rho | \psi \rangle = \left( 2(|a|^2)^2 - 2|a|^2 + 1 \right) \left[ (1-p)^3 + 3p(1-p)^2 \right] + \left( 2|a|^2 - 2(|a|^2)^2 \right) \left[ p^3 + 3p^2(1-p) \right]$$

We take the derivative of the above w.r.t.  $|a|^2$  and set it to zero to find the minimzing value of  $|a|^2$ :

$$\frac{\partial}{\partial |a|^2} \langle \psi | \rho | \psi \rangle = (4|a|^2 - 2) \left[ (1-p)^3 + 3p(1-p)^2 - p^3 + 3p^2(1-p) \right] = 0$$

Which for any value of p is satisfied when  $|a|^2=\frac{1}{2}$ , i.e.  $|a|=\frac{1}{\sqrt{2}}$  (and so  $|b|=\frac{1}{\sqrt{2}}$  as well). So, the fidelity is minimized for  $|\psi\rangle=\frac{1}{\sqrt{2}}(e^{i\varphi_0}\,|0\rangle+e^{i\varphi_1}\,|1\rangle)$  (which is what we might have expected given that the error correction succeeds only for the computational basis states, we should have the worst fidelity for states which are the furthest from both). For these states, the fidelity is (plugging

into our former general expression):

$$F_{\min} = \sqrt{|\psi|\rho\,|\psi\rangle} = \sqrt{\frac{(1-p)^3 + 3p(1-p)^2 + 3p^2(1-p) + p^3}{2}} = \frac{1}{\sqrt{2}}$$

# Exercise 10.5

Show that the syndrome measurement for detecting phase flip errors in the Shor code corresponds to measuring the observables  $X_1X_2X_3X_4X_5X_6$  and  $X_4X_5X_6X_7X_8X_9$ .

#### Solution

Concepts Involved: Linear Algebra, Eigenvalues, Eigenvectors, Error Syndrome

We have the codewords:

$$|0_L\rangle = \frac{(|000\rangle + |111\rangle)(|000\rangle + |111\rangle)(|000\rangle + |111\rangle)}{2\sqrt{2}}$$
  
$$|1_L\rangle = \frac{(|000\rangle - |111\rangle)(|000\rangle - |111\rangle)(|000\rangle - |111\rangle)}{2\sqrt{2}}$$

A phase flip error on a given block amounts to flipping the phase of the block:

$$Z_i(|000\rangle \pm |111\rangle) = |000\rangle \mp |111\rangle$$

Measuring  $X_{1-2} = X_1 X_2 X_3 X_4 X_5 X_6$  compares the phase of the first and second blocks (with +1 if they are the same, -1 if they are different) - we can see this from the eigenvalue relations:

$$\begin{split} X_1 X_2 X_3 X_4 X_5 X_6 (|000\rangle + |111\rangle) (|000\rangle + |111\rangle) &= +1 (|000\rangle + |111\rangle) (|000\rangle + |111\rangle) \\ X_1 X_2 X_3 X_4 X_5 X_6 (|000\rangle - |111\rangle) (|000\rangle + |111\rangle) &= +1 (|111\rangle - |000\rangle) (|000\rangle + |111\rangle) \\ &= -1 (|000\rangle - |111\rangle) (|000\rangle + |111\rangle) \\ X_1 X_2 X_3 X_4 X_5 X_6 (|000\rangle + |111\rangle) (|000\rangle - |111\rangle) &= +1 (|000\rangle + |111\rangle) (|111\rangle - |000\rangle) \\ &= -1 (|000\rangle + |111\rangle) (|000\rangle - |111\rangle) \\ X_1 X_2 X_3 X_4 X_5 X_6 (|000\rangle - |111\rangle) (|000\rangle - |111\rangle) &= +1 (|111\rangle - |000\rangle) (|111\rangle - |000\rangle) \\ &= +1 (|000\rangle - |111\rangle) (|000\rangle - |111\rangle) \end{split}$$

analogously, measuring  $X_{2-3}=X_4X_5X_6X_7X_8X_9$  compares the phase of the second and third blocks. The codewords have all three blocks with the same phase, so if we measure  $X_{1-2}=X_{2-3}=+1$  then we conclude no phase flip error occured. If we measure  $X_{1-2}=+1$  and  $X_{2-3}=-1$  then we conclude that a phase flip must have occured in the third block (one of qubits 7/8/9). If we measure  $X_{1-2}=-1$  and  $X_{2-3}=+1$  then we conclude that a phase flip occured on the first block (one of qubits 1/2/3). If we measure  $X_{1-2}=X_{2-3}=-1$  then we conclude that a phase flip error occured on the second block (one of qubits 1/2/3). We thus conclude that measuring these two operators yields the syndrome for detecting phase flip errors.

#### Exercise 10.6

Show that recovery from a phase flip on any of the first three qubits may be accomplished by applying the operator  $Z_1Z_2Z_3$ .

#### Solution

Concepts Involved: Linear Algebra, Error Recovery

We have the encoded state:

$$\begin{aligned} |\psi_L\rangle &= \alpha \, |0_L\rangle + \beta \, |1_L\rangle \\ &= \alpha \frac{(|000\rangle + |111\rangle)(|000\rangle + |111\rangle)(|000\rangle + |111\rangle)}{2\sqrt{2}} + \beta \frac{(|000\rangle - |111\rangle)(|000\rangle - |111\rangle)(|000\rangle - |111\rangle)}{2\sqrt{2}} \end{aligned}$$

We saw that  $Z_i(|000\rangle \pm |111\rangle) = |000\rangle \mp |111\rangle$  for  $i \in \{1, 2, 3\}$ , so if a phase flip error occurs on the first qubit, we have:

$$|\psi_L\rangle \xrightarrow{\mathcal{E}} \alpha \frac{(|000\rangle - |111\rangle)(|000\rangle + |111\rangle)(|000\rangle + |111\rangle)}{2\sqrt{2}} + \beta \frac{(|000\rangle + |111\rangle)(|000\rangle - |111\rangle)(|000\rangle - |111\rangle)}{2\sqrt{2}}$$

Now since  $Z_1 Z_2 Z_3 (|000\rangle \pm |111\rangle) = |000\rangle \pm (-1)^3 |111\rangle = |000\rangle \mp |111\rangle$ , we have:

$$Z_{1}Z_{2}Z_{3}\left(\alpha\frac{(|000\rangle - |111\rangle)(|000\rangle + |111\rangle)(|000\rangle + |111\rangle)}{2\sqrt{2}} + \beta\frac{(|000\rangle + |111\rangle)(|000\rangle - |111\rangle)(|000\rangle - |111\rangle)}{2\sqrt{2}}$$

$$= \alpha\frac{(|000\rangle + |111\rangle)(|000\rangle + |111\rangle)(|000\rangle + |111\rangle)}{2\sqrt{2}} + \beta\frac{(|000\rangle - |111\rangle)(|000\rangle - |111\rangle)(|000\rangle - |111\rangle)}{2\sqrt{2}}$$

$$= |\psi_{L}\rangle$$

so the error correction is accomplished.

# Exercise 10.7

Consider the three qubit bit flip code of Section 10.1.1, with corresponding projector  $P=|000\rangle\langle000|+|111\rangle\langle111|$ . The noise process this code protects against has operation elements  $\left\{\sqrt{(1-p)^3}I,\sqrt{p(1-p)^2}X_1,\sqrt{p(1-p)^2}X_2,\sqrt{p(1-p)^2}X_3\right\}$  where p is the probability that the bit flips. Note that this quantum operation is not trace-preserving, since we have omitted operation elements corresponding to bit flips on two and three qubits. Verify the quantum error-correction conditions for this code and noise process.

#### Solution

Concepts Involved: Linear Algebra, Projectors, Error Correction Conditions.

We calculate  $PE_i^{\dagger}E_jP$  for each of the errors  $E_i$ . Note that in this case the errors are all Hermitian, so this reduces to calculation of  $PE_iE_jP$  for all combinations of errors. Furthermore, note that the set of errors  $\left\{\sqrt{(1-p)^3}I,\sqrt{p(1-p)^2}X_1,\sqrt{p(1-p)^2}X_2,\sqrt{p(1-p)^2}X_3\right\}$  is mutually commuting, so

 $PE_iE_jP = PE_jE_iP$  so we only need to check all combinations (and the order does not effect the result).

$$P\sqrt{(1-p)^3}I\sqrt{(1-p)^3}IP = (1-p)^3P^2 = (1-p)^3P$$

$$P\sqrt{p(1-p)^2}X_1\sqrt{p(1-p)^2}X_1P = p(1-p)^2PIP = p(1-p)^2P^2 = p(1-p)^2P$$

$$P\sqrt{p(1-p)^2}X_2\sqrt{p(1-p)^2}X_2P = p(1-p)^2PIP = p(1-p)^2P^2 = p(1-p)^2P$$

$$P\sqrt{p(1-p)^2}X_3\sqrt{p(1-p)^2}X_3P = p(1-p)^2PIP = p(1-p)^2P^2 = p(1-p)^2P$$

where we have used that Paulis square to the identity and projectors are idempotent. For the other combinations we find:

$$P\sqrt{(1-p)^3}I\sqrt{p(1-p)^2}X_1P = \sqrt{p(1-p)^5}(|000\rangle\langle 000| + |111\rangle\langle 111|)(|100\rangle\langle 000| + |011\rangle\langle 111|) = 0 = 0P$$

$$P\sqrt{(1-p)^3}I\sqrt{p(1-p)^2}X_2P = \sqrt{p(1-p)^5}(|000\rangle\langle 000| + |111\rangle\langle 111|)(|010\rangle\langle 000| + |101\rangle\langle 111|) = 0 = 0P$$

$$P\sqrt{(1-p)^3}I\sqrt{p(1-p)^2}X_3P = \sqrt{p(1-p)^5}(|000\rangle\langle 000| + |111\rangle\langle 111|)(|001\rangle\langle 000| + |110\rangle\langle 111|) = 0 = 0P$$

$$P\sqrt{p(1-p)^2}X_1\sqrt{p(1-p)^2}X_2P = p(1-p)^2(|000\rangle\langle 100| + |111\rangle\langle 011|)(|010\rangle\langle 000| + |101\rangle\langle 111|) = 0 = 0P$$

$$P\sqrt{p(1-p)^2}X_1\sqrt{p(1-p)^2}X_3P = p(1-p)^2(|000\rangle\langle 100| + |111\rangle\langle 011|)(|001\rangle\langle 000| + |110\rangle\langle 111|) = 0 = 0P$$

$$P\sqrt{p(1-p)^2}X_2\sqrt{p(1-p)^2}X_3P = p(1-p)^2(|000\rangle\langle 010| + |111\rangle\langle 101|)(|001\rangle\langle 000| + |110\rangle\langle 111|) = 0 = 0P$$

So we therefore find:

$$PE_i^{\dagger}E_jP = \alpha_{ij}P$$

where:

$$\alpha = \begin{bmatrix} (1-p)^3 & 0 & 0 & 0\\ 0 & p(1-p)^2 & 0 & 0\\ 0 & 0 & p(1-p)^2 & 0\\ 0 & 0 & 0 & p(1-p)^2 \end{bmatrix}$$

is Hermitian. The error-correction conditions are therefore verified.

#### Exercise 10.8

Verify that the three qubit phase flip code  $|0_L\rangle = |+++\rangle$ ,  $|1_L\rangle = |---\rangle$  satisfies the quantum error-correction conditions for the set of error operators  $\{I, Z_1, Z_2, Z_3\}$ .

# Solution

Concepts Involved: Linear Algebra, Projectors, Error Correction Conditions.

First, note the projector onto the code C in this case is:

$$P = |+ + +\rangle + + + |+ |- - -\rangle - - -|$$

We calculate  $PE_i^\dagger E_j P$  for each of the errors  $E_i \in \{I, Z_1, Z_2, Z_3\}$ . The calculation is completely analogous to that in Ex. 10.7, simply replacing  $X \to Z$ ,  $|0\rangle \to |+\rangle$ ,  $|1\rangle \to |-\rangle$ , and setting all of the pre-factors  $\sqrt{(1-p)^3}$  and  $\sqrt{p(1-p)^2}$  to one. The result we find is:

$$PE_i^{\dagger}E_jP = \alpha_{ij}P$$

where:

$$\alpha = \begin{bmatrix} 1 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \end{bmatrix}$$

which is of course Hermitian, and the QEC conditions are thus satisfied.

# Exercise 10.9

Again, consider the three qubit phase flip code. Let  $P_i$  and  $Q_i$  be the projectors into the  $|0\rangle$  and  $|1\rangle$  states, respectively, of the ith qubit. Prove that the three qubit phase flip code protects against the error set  $\{I, P_1, Q_1, P_2, Q_2, P_3, Q_3\}$ .

#### Solution

Concepts Involved: Linear Algebra, Error Correction Conditions, Discretization of Errors.

By Theorem 10.2 in the text, we know that if C is a quantum code and  $\mathcal R$  is the error-correction procedure constructed via the error correction conditions that corrects for a noise process  $\mathcal E$  with operation elements  $\{E_i\}$ , then  $\mathcal R$  also corrects for arbitrary complex linear combinations of the  $E_i$ . We then note that all errors in  $\{I, P_1, Q_1, P_2, Q_2, P_3, Q_3\}$  can be written as linear combinations of errors in  $\{I, Z_1, Z_2, Z_3\}$ :

$$I = I$$
,  $P_i = \frac{I + Z_i}{2}$ ,  $Q_i = \frac{I - Z_i}{2}$ 

therefore by Theorem 10.2 and the result of Ex. 10.8, the phase flip code can protect against  $\{IP_1, Q_1, P_2, Q_2, P_3, Q_3\}$ .

Explicitly verify the quantum error-correction conditions for the Shor code, for the error set containing I and the error operators  $X_i, Y_i, Z_j$  for j = 1 through 9.

#### Solution

Concepts Involved: Linear Algebra, Projectors, Error Correction Conditions.

We have:

$$P = |0_L\rangle\langle 0_L| + |1_L\rangle\langle 1_L|$$

where:

$$\begin{split} |0_L\rangle &= \frac{(|000\rangle + |111\rangle)(|000\rangle + |111\rangle)(|000\rangle + |111\rangle)}{2\sqrt{2}} \\ |1_L\rangle &= \frac{(|000\rangle - |111\rangle)(|000\rangle - |111\rangle)(|000\rangle - |111\rangle)}{2\sqrt{2}} \end{split}$$

All  $E_i$  we consider in the set of errors are Hermitian, so  $E_i^{\dagger} = E_i$ . Firstly, since all Pauli operators square to the identity, we find for all  $E_i$  that:

$$PE_i^{\dagger}E_iP = PE_i^2P = PIP = P^2 = P$$

Next, we observe:

$$PIX_kP = PIY_kP = PIZ_kP = PX_kZ_kP = PZ_kX_kP = PX_kY_kP = PY_kX_kP = PY_kZ_kP = PZ_kY_kP$$
$$= 0 = 0P$$

as all possible bit/phase flips on a single qubit map the codewords to states orthogonal to both codewords. Next, we find that for  $k \neq l$  that:

$$PX_kX_lP = PY_kY_lP = PX_kY_lP = PY_kX_lP = PX_kZ_lP = PZ_kX_lP = PY_kZ_lP = PZ_kY_lP$$
$$= 0 = 0P$$

as in each case we have a bit flip on one or two qubits which maps the codewords to state orthogonal to both codwords.

Finally, we find that:

$$PZ_kZ_lP = \begin{cases} P & \text{if } k,l \text{ are in the same block} \\ 0 & \text{if } k,l \text{ are in different blocks} \end{cases}$$

as in the former case the two phase flips cancel out (and thus P is preserved), and in the latter case we have phase flips on two different blocks and the codewords are mapped to states orthogonal to both codewords.

We thus conclude that:

$$PE_i^{\dagger}E_jP = \alpha_{ij}P$$

where  $\alpha_{ij}$  is Hermitian (as it has 1s on the diagonal, 1s on entries for  $Z_iZ_j$  with i,j in the same block (thus symmetric across the diagonal), and zero elsewhere). We have thus verified the quantum error correction conditions.

### Exercise 10.11

Construct operation elements for a single qubit quantum operation  $\mathcal E$  that upon input of any state  $\rho$  replaces it with the completely randomized state I/2. It is amazing that even such noise models as this may be corrected by codes such as the Shor code!

#### Solution

Concepts Involved: Linear Algebra, Density Operators, Quantum Operations.

We wish to find operation elements  $\{E_k\}$  such that:

$$\mathcal{E}(\rho) = \sum_{k} E_{k} \rho E_{k}^{\dagger} = \frac{I}{2}$$

for any input single-qubit state  $\rho$ . We claim that  $\left\{\frac{1}{2}I, \frac{1}{2}X, \frac{1}{2}Y, \frac{1}{2}Z\right\}$  are the operation elements with this desired property. First, we verify that they satisfy the completeness relation:

$$\sum_{k} E_{k}^{\dagger} E_{k} = \frac{1}{4} (I^{2} + X^{2} + Y^{2} + Z^{2}) = \frac{1}{4} (4I) = I$$

Next, we verify that they have the claimed property of sending every initial qubit state to the maximally mixed state. From Ex. 2.72, we can write a single-qubit density operator as:

$$\rho = \frac{I + r_x X + r_y Y + r_z Z}{2}$$

where  $\mathbf{r}=(r_x,r_y,r_z)\in\mathbb{R}^3$  and  $\|\mathbf{r}\|=1$ . Now calculating  $\mathcal{E}(\rho)$ , we have:

$$\begin{split} \mathcal{E}(\rho) &= \frac{1}{4}(\rho + X\rho X + Y\rho Y + Z\rho Z) \\ &= \frac{1}{8}\left((I + X^2 + Y^2 + Z^2) + r_x(X + X^3 + YXY + ZXZ) \right. \\ &+ r_y(Y + XYX + Y^3 + ZYZ) + r_z(Z + XZX + YZY + Z^3)\right) \\ &= \frac{1}{8}\left(4I + r_x(2X + 2iYZ) + r_y(2Y + 2iZX) + r_z(2Z + 2iXY)\right) \\ &= \frac{1}{8}\left(4I + r_x(2X + 2i(iX) + r_y(2Y + 2i(iY)) + r_z(2Z + 2i(iZ))\right) \\ &= \frac{1}{8}(4I) \\ &= \frac{I}{2} \end{split}$$

where we use that XY=iZ,YZ=iX, and XZ=iY. The claim is thus proven.

#### Remark:

The described channel is of course the single-qubit depolarizing channel of full strength.

#### Exercise 10.12

Show that the fidelity between the state  $|0\rangle$  and  $\mathcal{E}(|0\rangle\langle 0|)$  is  $\sqrt{1-2p/3}$ , and use this to argue that the minimum fidelity for the depolarizing channel is  $\sqrt{1-2p/3}$ .

### Solution

Concepts Involved: Linear Algebra, Density Operators, Quantum Operations, Fidelity.

The density operator corresponding to  $|0\rangle$  is  $|0\rangle\langle 0|$ , and sending this through the depolarizing channel we have:

$$\begin{split} \mathcal{E}(|0\rangle\!\langle 0|) &= (1-p)\,|0\rangle\!\langle 0| + \frac{p}{3}(X\,|0\rangle\!\langle 0|\,X + Y\,|0\rangle\!\langle 0|\,Y + Z\,|0\rangle\!\langle 0|\,Z) \\ &= (1-p)\,|0\rangle\!\langle 0| + \frac{p}{3}(|1\rangle\!\langle 1| + |1\rangle\!\langle 1| + |0\rangle\!\langle 0|) \\ &= (1-\frac{2p}{3})\,|0\rangle\!\langle 0| + \frac{2p}{3}\,|1\rangle\!\langle 1| \end{split}$$

and so:

$$F(|0\rangle, \mathcal{E}(|0\rangle\langle 0|)) = \sqrt{\langle 0|\left(\left(1 - \frac{2p}{3}\right)|0\rangle\langle 0| + \frac{2p}{3}|1\rangle\langle 1|\right)|0\rangle} = \sqrt{1 - \frac{2p}{3}}$$
 (6)

as claimed. Because the depolarizing channel is symmetric in X/Y/Z, it is therefore symmetric in possible input states and so for any input state  $|\psi\rangle$  we would find that  $F(|0\rangle,\mathcal{E}(|0\rangle\langle 0|))=\sqrt{1-2p/3}$ . As such, this is the minimum fidelity.

For a more rigorous argument; by from Ex. 2.72, we can write a single-qubit density operator as:

$$\rho = \frac{I + r_x X + r_y Y + r_z Z}{2}$$

where  $\mathbf{r}=(r_x,r_y,r_z)\in\mathbb{R}^3$  and  $\|\mathbf{r}\|=1$  when  $\rho=|\psi\rangle\!\langle\psi|$  a pure state. We then have that:

$$\langle \psi | X | \psi \rangle = \text{Tr}(\rho_{\psi} X) = \text{Tr}\left(\frac{X + r_x I + r_y Y X + r_z Z X}{2}\right) = r_x$$

and analogously for Y/Z. We then have:

$$F(|\psi\rangle, \mathcal{E}(|\psi\rangle\langle\psi|)) = \sqrt{\langle\psi|\left((1-p)|\psi\rangle\langle\psi| + \frac{p}{3}(X|\psi\rangle\langle\psi|X + Y|\psi\rangle\langle\psi|Y + Z|\psi\rangle\langle\psi|Z)\right)}$$

$$= \sqrt{(1-p)|\psi\rangle\langle\psi|^2 + \frac{p}{3}(\langle\psi|X|\psi\rangle^2 + \langle\psi|Y|\psi\rangle^2 + \langle\psi|Z|\psi\rangle^2}$$

$$= \sqrt{(1-p) + \frac{p}{3}(r_x^2 + r_y^2 + r_z^2)}$$

$$= \sqrt{(1-p) + \frac{p}{3}}$$

$$= \sqrt{1 - \frac{2p}{3}}$$

where we use that  $\|\mathbf{r}\| = 1$  in the fourth equality. Since this is true for all pure states, it must be the minimum fidelity.

### Exercise 10.13

Show that the minimum fidelity  $F(|\psi\rangle, \mathcal{E}(|\psi\rangle\langle\psi|))$  when  $\mathcal{E}$  is the amplitude damping channel when  $\mathcal{E}$  is the amplitude damping channel with parameter  $\gamma$  is  $\sqrt{1-\gamma}$ .

#### Solution

Concepts Involved: Linear Algebra, Fidelity.

Let  $|\psi\rangle=a\,|0\rangle+b\,|1\rangle$  with  $|a|^2+|b|^2=1.$  Then we have:

$$|\psi\rangle\langle\psi| = \begin{bmatrix} |a|^2 & ab^* \\ a^*b & |b|^2 \end{bmatrix}$$

and after applying the amplitude damping channel we have (from Ex. 8.22):

$$\mathcal{E}(|\psi\rangle\!\langle\psi|) = \begin{bmatrix} 1 - (1-\gamma)(1-|a|^2) & ab^*\sqrt{1-\gamma} \\ a^*b\sqrt{1-\gamma} & |b|^2(1-\gamma) \end{bmatrix}$$

Calculating the fidelity, we then have:

$$\begin{split} F(|\psi\rangle\,,\mathcal{E}(|\psi\rangle\langle\psi|)) &= \sqrt{\langle\psi|\,\mathcal{E}(|\psi\rangle\langle\psi|)\,|\psi\rangle} \\ &= \sqrt{\begin{bmatrix} a^* & b^* \end{bmatrix} \begin{bmatrix} 1 - (1-\gamma)(1-|a|^2) & ab^*\sqrt{1-\gamma} \\ a^*b\sqrt{1-\gamma} & |b|^2(1-\gamma) \end{bmatrix} \begin{bmatrix} a \\ b \end{bmatrix}} \\ &= \sqrt{|a|^2(1-(1-\gamma)(1-|a|^2)) + 2|a|^2|b|^2\sqrt{1-\gamma} + |b|^4(1-\gamma)} \end{split}$$

Using the normalization condition,  $|b|^2 = 1 - |a|^2$  so we can write the above in terms of  $|a|^2$  alone:

$$F(|\psi\rangle, \mathcal{E}(|\psi\rangle\langle\psi|)) = \sqrt{|a|^2(1 - (1 - \gamma)(1 - |a|^2)) + 2|a|^2(1 - |a|^2)\sqrt{1 - \gamma} + (1 - |a|^2)^2(1 - \gamma)}$$
$$= \sqrt{2(1 - \sqrt{1 - \gamma} - \gamma)(|a|^2)^2 + (-2 + 2\sqrt{1 - \gamma} + 3\gamma)|a|^2 + 1 - \gamma}$$

This is minimized when the expression under the square root is minimized. Taking the derivative w.r.t  $\left|a\right|^2$ , we find:

$$\frac{\partial}{\partial |a|^2} (F(|\psi\rangle, \mathcal{E}(|\psi\rangle\langle\psi|)))^2 = 4(1 - \sqrt{1 - \gamma} - \gamma)|a|^2 - 2 + 2\sqrt{1 - \gamma} + 3\gamma$$

which is non-negative for  $|a|^2 \in [0,1]$  (which is the domain over which it is defined). Therefore  $F(|\psi\rangle, \mathcal{E}(|\psi\rangle\langle\psi|))$  is an increasing function of  $|a|^2$  over [0,1], and hence  $F(|\psi\rangle, \mathcal{E}(|\psi\rangle\langle\psi|))$  is minimized when a=0 and therefore  $|\psi\rangle=|1\rangle$ . In this case, the fidelity is:

$$F(|\psi\rangle, \mathcal{E}(|\psi\rangle\langle\psi|))_{\min} = F(|1\rangle, \mathcal{E}(|1\rangle\langle1|)) = \sqrt{1-\gamma}$$

as claimed.

### Exercise 10.14

Write an expression for a generator matrix encoding k bits using r repetitions for each bit. This is an [rk, k] linear code, and should have an  $rk \times k$  generator matrix.

#### Solution

Concepts Involved: Linear Algebra, Generator Matrices.

The claimed generator matrix G is an  $rk \times k$  matrix such that:

$$G_{ij} = egin{cases} 1 & r(j-1) < i \leq rj \\ 0 & \text{otherwise}. \end{cases}$$

By matrix multiplication, we can see that:

$$G(x_1,x_2,\ldots,x_k) = (\overbrace{x_1,\ldots,x_1}^{r \text{ times}},\overbrace{x_2,\ldots,x_2}^{r \text{ times}},\ldots,\overbrace{x_k,\ldots,x_k}^{r \text{ times}})$$

Show that adding one column of G to another results in a generator matrix generating the same code.

#### Solution

Concepts Involved: Linear Algebra, Generator Matrices.

The set of possible codewords of a code corresponds to the vector space spanned by the columns of G. So if  $G = [\mathbf{v}_1, \mathbf{v}_2, \dots \mathbf{v}_k]$  then the possible codewords are:

$$\mathbf{v} = c_1 \mathbf{v}_1 + c_2 \mathbf{v}_2 + \ldots + c_k \mathbf{v}_k$$

where  $c_i \in \mathbb{Z}_2$  and addition is done modulo 2. WLOG suppose we add column 2 to column 1 (permuting the columns of G clearly preserves the codespace, as it just amounts to permuting labels in the equation above), so we have  $G' = [\mathbf{v}_1 + \mathbf{v}_2, \mathbf{v}_2, \dots, \mathbf{v}_k]$ , so the possible codewords are:

$$\mathbf{v}' = c_1'(\mathbf{v}_1 + \mathbf{v}_2) + c_2'\mathbf{v}_2 + \ldots + c_k'\mathbf{v}_k$$

where  $c_i' \in \mathbb{Z}_2$ . By defining  $c_1' = c_1, c_2' = c_1 + c_2$ , and  $c_j' = c_j$  for  $j \ge 2$  we can see that the codewords  $\mathbf{v}'$  are the same as the codewords  $\mathbf{v}$ , and thus G and G' generate the same code.

### Exercise 10.16

Show that adding one row of the parity check matrix to another does not change the code. Using Gaussian elimination and swapping of bits it is therefore possible to assume that the parity check matrix has the standard form  $[A|I_{n-k}]$  where A is an  $(n-k)\times k$  matrix.

#### Solution

Concepts Involved: Linear Algebra, Parity Check Matrices.

In the parity check matrix formulation, an [n,k] code is all  $\mathbf{x} \in \mathbb{Z}_2^n$  such that Hx = 0 where  $H \in \mathbb{Z}_2^{(n-k) \times n}$  is the parity check matrix. We can write:

$$H = \begin{bmatrix} \mathbf{v}_1^T \\ \mathbf{v}_2^T \\ \vdots \\ \mathbf{v}_n^T \end{bmatrix}$$

with  $\mathbf{v}_i^T$  the rows of H. By the definition of matrix multiplication, it follows from Hx=0 that:

$$\mathbf{v}_i \cdot \mathbf{x} = 0$$

for each i and for all codewords x. WLOG suppose we add row 2 to row 1 (permuting the rows of H

clearly preserves the codespace, as the above condition is unchanged). We then have:

$$H' = egin{bmatrix} \mathbf{v}_1^T + \mathbf{v}_2^T \ \mathbf{v}_2^T \ dots \ \mathbf{v}_n^T \end{bmatrix}$$

It then follows that if Hx = 0, then H'x = 0 as:

$$(\mathbf{v}_1 + \mathbf{v}_2) \cdot \mathbf{x} = \mathbf{v}_1 \cdot \mathbf{x} + \mathbf{v}_2 \cdot \mathbf{x} = 0 + 0 = 0$$

and  $\mathbf{v}_i \cdot \mathbf{x} = 0$  for  $i \geq 2$ . Furthermore, if H'x = 0 then Hx = 0 as:

$$\mathbf{v}_1 \cdot \mathbf{x} = (\mathbf{v}_1 + \mathbf{v}_2 + \mathbf{v}_2) \cdot \mathbf{x} = (\mathbf{v}_1 \cdot \mathbf{v}_2) \cdot \mathbf{x} + \mathbf{v}_2 \cdot \mathbf{x} = 0 + 0 = 0$$

and  $\mathbf{v}_i \cdot \mathbf{x} = 0$  for  $i \geq 2$ . Therefore, H, H' correspond to the same code.

Since Gaussian elimination only involves swapping rows (does nothing to H), swapping columns (changes the labels of the qubits) and adding rows to each other (does nothing as shown above), we can thus always assume that the parity check matrix can be brought to standard form.

### Exercise 10.17

Find a parity check matrix for the [6,2] repetition code defined by the generator matrix in (10.54).

### Solution

Concepts Involved: Linear Algebra, Generator Matrices, Parity Check Matrices.

The [6, 2] repetition code has generator matrix:

$$G = \begin{bmatrix} 1 & 0 \\ 1 & 0 \\ 1 & 0 \\ 0 & 1 \\ 0 & 1 \\ 0 & 1 \end{bmatrix}$$

To construct H, we pick out 6-2=4 linearly independent vectors orthogonal to the columns of G. Four such vectors are:

Therefore one such parity check matrix is:

$$\begin{bmatrix} 1 & 1 & 0 & 0 & 0 & 0 \\ 0 & 1 & 1 & 0 & 0 & 0 \\ 0 & 0 & 0 & 1 & 1 & 0 \\ 0 & 0 & 0 & 0 & 1 & 1 \end{bmatrix}$$

### Exercise 10.18

Show that the parity check matrix H and generator matrix G for the same linear code satisfy HG=0

#### Solution

Concepts Involved: Linear Algebra, Genrator Matrices, Parity Check Matrices.

For a given [n, k] code with parity check matrix H and generator matrix G, we can write:

$$H = egin{bmatrix} \mathbf{v}_1^T \ \mathbf{v}_2^T \ dots \ \mathbf{v}_{n-k}^T \end{bmatrix}, \quad G = egin{bmatrix} \mathbf{y}_1 & \mathbf{y}_2 & \cdots & \mathbf{y}_k \end{bmatrix}$$

where  $\mathbf{v}_i, \mathbf{y}_j$  are each n-dimensional vectors which are orthogonal to one another. By the definition of matrix multiplication, HG is a  $n-k \times k$  matrix with entries:

$$(HG)_{ij} = \mathbf{v}_i \cdot \mathbf{y}_j = 0$$

where the last equality follows by orthogonality, thus proving the claim.

### Exercise 10.19

Suppose an [n,k] linear code C has a parity check matrix of the form  $H=[A|I_{n-k}]$ , for some  $(n-k)\times k$  matrix A. Show that the corresponding generator matrix is

$$G = \left[\frac{I_k}{-A}\right].$$

#### Solution

Concepts Involved: Linear Algebra, Genrator Matrices, Parity Check Matrices.

To get a generator matrix from a parity check matrix, we pick k linearly independent vectors  $\mathbf{y}_1, \dots \mathbf{y}_k$  spanning the kernel of H, and set G to have columns  $\mathbf{y}_1$  through  $\mathbf{y}_{y_k}$ . In this case, H has standard form

 $H = [A|I_{n-k}]$  and so we may write:

$$H = egin{bmatrix} \mathbf{v}_1^T & \mathbf{e}_{1,n-k}^T \ \mathbf{v}_2^T & \mathbf{e}_{2,n-k}^T \ dots \ \mathbf{v}_{n-k}^T & \mathbf{e}_{n-k,n-k}^T \end{bmatrix}$$

where  $\mathbf{v}_i^T$  are the rows of A and  $\mathbf{e}_{i,n-k}$  is a n-k length vector with 1 in the ith position and 0s elsewhere. We want to find vectors in the kernel of H, i.e. vectors  $\mathbf{y}$  that satisfy:

$$\mathbf{v}_i \cdot \mathbf{y}_{1,\dots k} + \mathbf{e}_{i,n-k} \cdot \mathbf{y}_{k+1,\dots n}$$

for every  $i \in \{1, ..., n-k\}$ . A clear choice that satisfies this relation is  $\mathbf{y}_{1,...,n-k} = \mathbf{e}_{n,k}$  and  $\mathbf{y}_{k+1,...,n} = -\mathbf{w}_n$  where  $\mathbf{w}_n$  is the nth column of A; this choice satisfies the above as:

$$\mathbf{v}_i \cdot \mathbf{e}_{n,k} + \mathbf{e}_{i,n-k} \cdot (-\mathbf{w}_n) = a_{in} - a_{in} = 0$$

This yields  $y_i, \dots y_k$ , one for each column of A. We therefore construct G as:

$$G = \begin{bmatrix} \mathbf{e}_{1,k} & \mathbf{e}_{2,k} & \dots & \mathbf{e}_{k,k} \\ -\mathbf{w}_1 & -\mathbf{w}_2 & \dots & -\mathbf{w}_n \end{bmatrix} = \begin{bmatrix} I_k \\ -A \end{bmatrix}$$

which was what we wished to show.

#### Exercise 10.20

Let H be a parity check matrix such that any d-1 columns are linearly independent, but there exists a set of d linearly dependent columns. Show that the code defined by H has distance d.

### Solution

Concepts Involved: Linear Algebra, Parity Check Matrices, Code Distance.

Given a parity check matrix H, a codeword  $\mathbf{x} = (x_1, x_2, \dots x_n)$  satisfies  $H\mathbf{x} = 0$ , which implies:

$$x_1\mathbf{h}_1 + x_2\mathbf{h}_2 + \ldots + x_n\mathbf{h}_n = 0$$

where  $\mathbf{h}_i$  are the columns of H. Since any d-1 columns are linearly independent, it follows that any sum of d-1 (or less) columns of H is nonzero, and therefore there are no codewords of weight d-1 or lower. However, there exists a set of d linearly dependent columns, and therefore there exists a codeword that is 1 for each of those columns and 0 elsewhere, and is therefore of weight d. There are no codewords of smaller weight, and therefore we conclude that the code has distance d.

### **Exercise 10.21: Singleton bound**

Show than an [n, k, d] code must satisfy  $n - k \ge d - 1$ .

Concepts Involved: Linear Algebra, Parity Check Matrices, Code Distance.

An [n,k,d] code has an  $n-k \times n$  parity check matrix H. It therefore has at most n-k linearly independent columns. From the solution to the previous exercise, if a code has distance d then the parity check matrix must have all sets of d-1 columns be linearly independent (else there would exist a codeword of weight d-1 and hence the code would have distance < d, a contradiction). Combining these two facts, it follows that:

$$n-k \ge d-1$$

as claimed.

# Exercise 10.22

Show that all Hamming codes have distance 3, and thus can correct an error on a single bit. The Hamming codes are therefore  $[2^r - 1, 2^r - r - 1, 3]$  codes.

### Solution

Concepts Involved: Linear Algebra, Parity Check Matrices, Code Distance.

Recall that the  $[2^r - 1, 2^r - r - 1]$  Hamming code (for  $r \ge 2$ ) is a linear code having parity check matrix whose columns are all  $2^r - 1$  bit strings of length r which are not zero.

Any two columns of H are different and therefore linearly independent. Furthermore, H has the 3 columns:

$$\begin{bmatrix} 1 \\ 0 \\ 0 \\ 0 \\ \vdots \\ 0 \end{bmatrix}, \quad \begin{bmatrix} 0 \\ 1 \\ 0 \\ \vdots \\ 0 \end{bmatrix}, \quad \begin{bmatrix} 1 \\ 1 \\ 0 \\ \vdots \\ 0 \end{bmatrix}$$

which are linearly dependent as they add together to make zero. By Ex. 10.20, any Hamming code has distance 3. Since a distance  $d \geq 2t+1$  code can correct errors on t bits, all Hamming codes can correct errors on one bit.

### Exercise 10.23

(\*) Prove the Gilbert-Varshamov bound.

# Exercise 10.24

Show that a code with generator matrix G is weakly self-dual if and only if  $G^TG=0$ .

Concepts Involved: Linear Algebra, Generator Matrices, Dual codes.

Recall that if we have a [n,k] code C with generator matrix G and parity check matrix H, then the dual code  $C^{\perp}$  is the code consisting of all codewords y such that y is orthogonal to all codewords in C. Furthermore, recall that a code is weakly self-dual if  $C \subseteq C^{\perp}$ .

 $\Longrightarrow$ : Suppose that G is weakly self dual. Then, it follows that every codeword  $y=Gx\in C$  is contained in  $C^{\perp}$ . Since the parity check matrix applied to a codeword of a code gives zero, for  $C^{\perp}$  we have that  $G^Ty=0$  and thus  $G^Ty=G^TGx=0$  for all length k binary vectors x, but this is only possible if  $G^TG=0$ .

 $\sqsubseteq$  : Suppose that  $G^TG=0$ . Suppose that y=Gx,y'=Gx' are codewords of C. We then see that  $y'\cdot y=x'^TG^TGx=x'^T0x=0$ . Thus, all codewords of C are orthogonal to each other, and thus all codewords y of C are codewords of  $C^\perp$  by definition. Thus  $C\subseteq C^\perp$ .

### Exercise 10.25

Let C be a linear code. Show that if  $x \in C^{\perp}$  then  $\sum_{y \in C} (-1)^{x \cdot y} = |C|$ , while if  $x \notin C^{\perp}$  then  $\sum_{y \in C} (-1)^{x \cdot y} = 0$ .

#### Solution

Concepts Involved: Linear Algebra, Dual Codes

Suppose  $x\in C^\perp$ . Then it follows that  $y\cdot x=0$  for every codeword  $y\in C$ . Thus,  $\sum_{y\in C}(-1)^{x\cdot y}=\sum_{y\in C}(-1)^0=\sum_{y\in C}1=|C|$ . Suppose instead that  $x\notin C^\perp$ . Then we claim that half the codewords of C are orthogonal to x with  $x\cdot y=0$ , and the other half are not orthogonal with  $x\cdot y=1$ . As proof of this fact, suppose there are n basis codewords  $y_1,\ldots,y_n$  of C; since  $x\notin C^\perp$ , it follows that at least one of these basis codewords is not orthogonal to x. Let  $y_1,\ldots,y_k$  be the non-orthogonal basis codewords and  $y_{k+1},\ldots,y_n$  be the orthogonal codewords. To form an arbitrary codeword of C, we take a linear combination of the basis codewords. If an even number fo  $y_1,\ldots,y_k$  are included in the sum then the codeword is orthogonal to x and otherwise the codeword is not orthogonal to x. There are  $2^{k-1}$  to choose an even number of elements out of a set of k elements and k0 ways to choose an odd number, and therefore there are the same number of codewords in k0 that are orthogonal to k1 as there are those that are not orthogonal. Hence, the contributions of these two evenly weighted parts cancel in the sum to yield k1 to k2 the contributions of these two evenly weighted parts cancel in the sum to yield k3 there are those that are not orthogonal.

### Exercise 10.26

Suppose H is a parity check matrix. Explain how to compute the transformation  $|x\rangle|0\rangle\mapsto|x\rangle|Hx\rangle$  using a circuit composed entirely of controlled-NOTs.

#### Solution

Concepts Involved: Linear Algebra, Parity Check Matrices, Controlled Operations

By the definition of matrix multiplication:

$$(Hx)_i = \sum_j H_{ij} x_j \tag{7}$$

So, if we want the ith bit in the second register to become  $|(Hx)_i\rangle$ , this can be realized by applying a CNOT gate with control on the jth qubit of the first register (in  $|x\rangle$ ) and acting on the ith qubit of the second register if  $H_{ij}=1$  (and doing nothing if  $H_{ij}=0$ ); each gate (or identity) realizing one term in the above sum.

#### Exercise 10.27

Show that the codes defined by

$$|x + C_2\rangle \equiv \frac{1}{\sqrt{|C_2|}} \sum_{y \in C_2} (-1)^{u \cdot y} |x + y + v\rangle$$

as parameterized by u and v are equivalent to  $CSS(C_1,C_2)$  in the sense that they have the same error-correcting properties. These codes, which we'll refer to as  $CSS_{u,v}(C_1,C_2)$  will be useful later in our study of quantum key distribution, in Section 12.6.5.

#### Exercise 10.28

Verify that the transpose of the matrix in (10.77) is the generator of the [7,4,3] Hamming code.

#### Solution

Concepts Involved: Linear Algebra, Parity Check Matrices, Generator Matrices

The matrix in (10.77) is:

$$H[C_2] = egin{bmatrix} 1 & 0 & 0 & 0 & 0 & 1 & 1 \ 0 & 1 & 0 & 0 & 1 & 0 & 1 \ 0 & 0 & 1 & 0 & 1 & 1 & 0 \ 0 & 0 & 0 & 1 & 1 & 1 & 1 \end{bmatrix}$$

Taking the transpose, we have:

$$H[C_2]^T = \begin{bmatrix} 1 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \\ 0 & 1 & 1 & 1 \\ 1 & 0 & 1 & 1 \\ 1 & 1 & 0 & 1 \end{bmatrix}$$

Looking at the parity check matrix for the [7,4,3] Hamming code, we have:

$$H[C_1] = \begin{bmatrix} 0 & 0 & 0 & 1 & 1 & 1 & 1 \\ 0 & 1 & 1 & 0 & 0 & 1 & 1 \\ 1 & 0 & 1 & 0 & 1 & 0 & 1 \end{bmatrix}$$

In order to confirm  $H[C_2]^T$  as the generator matrix of the code, we require that the columns are linearly independent and lie in the kernel of  $H[C_1]$ . The linearly independence is immediately clear by inspection (as columns i=1,2,3,4 uniquely have a non-zero ith entry). Let us then just verify that they lie in the kernel of H, which can equivalently be done by checking that  $H[C_1]H[C_2]^T=0$ :

Exercise 10.29

Show that an arbitrary linear combination of any two elements of  $V_S$  is also in  $V_S$ . Therefore,  $V_S$  is a subspace of the n qubit state space. Show that  $V_S$  is the intersection of the subspaces fixed by each operator in S (that is, the eigenvalue one eigenspaces of elements of S).

### Exercise 10.30

Show that  $-I \notin S$  implies  $\pm iI \notin S$ .

## Solution

Concepts Involved: Group Theory, Pauli Group

Subgroups are groups and therefore closed under multiplication. If  $iI \in S$ , then  $(iI)^2 = (i)^2I = -I \in S$ . The claim is thus shown via the contrapositive.

Suppose that S is a subgroup of  $G_n$  generated by elements  $g_1, \ldots, g_l$ . Show that all the elements of S commute if and only if  $g_i$  and  $g_j$  commute for each pair i, j.

### Solution

Concepts Involved: Group Theory, Pauli Group

 $\Longrightarrow$ : Suppose all elements of S commute. Since each generator is itself an element of S, any pair of generators will commute.

 $\sqsubseteq$  : Suppose any pair of generators  $g_i, g_j$  of S commute. Any two elements  $s_1, s_2 \in S$  can be written as a product of generators:

$$s_1 = g_1^{n_1} g_2^{n_2} \dots g_k^{n_k}, \quad s_2 = g_1^{m_1} g_2^{m_2} \dots g_k^{m_k}$$

Using the pairwise commutation of the generators, we can move the  $g_i$ s together to find that  $s_1s_2=s_2s_1=g_1^{n_1+m_1}g_2^{n_2+m_2}\dots g_k^{n_k+m_k}$  and so all elements of S commute.  $\square$ 

**Remark:** The above argument holds for any group, not just subgroups of  $G_n$ .

### Exercise 10.32

Verify that the generators in Figure 10.6 stabilize the codewords for the Steane code, as described in Section 10.4.2.

#### Exercise 10.33

Show that g and g' commute if and only if  $r(g)\Lambda r(g')^T=0$ . (In the check matrix representation, arithmetic is done modulo two.)

### Exercise 10.34

Let  $S = \langle g_1, \dots, g_l \rangle$ . Show that -I is not an element of S if and only if  $g_j^2 = I$  for all j, and  $g_j \neq -I$  for all j.

### Exercise 10.35

Let S be a subgroup of  $G_n$  such that -I is not an element of S. Show that  $g^2 = I$  for all  $g \in S$ , and thus  $g^{\dagger} = g$ .

### Exercise 10.36

Explicitly verify that  $UX_1U^{\dagger}=X_1X_2$ ,  $UX_2U^{\dagger}=X_2$ ,  $UZ_1U^{\dagger}=Z_1$ , and  $UZ_2U^{\dagger}=Z_1Z_2$ . These and other useful conjugation relations for the Hadamard, phase, and Pauli gates are summarized in Figure 10.7.

Concepts Involved: Linear Algebra, Unitary Operators, Controlled Operations

First note our ability to write the controlled-NOT operation in block diagonal form:

$$U = U^{\dagger} = \begin{bmatrix} I & 0 \\ 0 & X \end{bmatrix}$$

By explicit (block) matrix multiplication we then have:

$$UX_1U^\dagger = \begin{bmatrix} I & 0 \\ 0 & X \end{bmatrix} \begin{bmatrix} 0 & I \\ I & 0 \end{bmatrix} \begin{bmatrix} I & 0 \\ 0 & X \end{bmatrix} = \begin{bmatrix} 0 & I \\ X & 0 \end{bmatrix} \begin{bmatrix} I & 0 \\ 0 & X \end{bmatrix} = \begin{bmatrix} 0 & X \\ X & 0 \end{bmatrix} = X_1X_2$$

$$UX_2U^\dagger = \begin{bmatrix} I & 0 \\ 0 & X \end{bmatrix} \begin{bmatrix} X & 0 \\ 0 & X \end{bmatrix} \begin{bmatrix} I & 0 \\ 0 & X \end{bmatrix} = \begin{bmatrix} X & 0 \\ 0 & X^2 \end{bmatrix} \begin{bmatrix} I & 0 \\ 0 & X \end{bmatrix} = \begin{bmatrix} X & 0 \\ 0 & I \end{bmatrix} \begin{bmatrix} I & 0 \\ 0 & X \end{bmatrix} = \begin{bmatrix} X & 0 \\ 0 & X \end{bmatrix} = X_2$$

$$UZ_1U^\dagger = \begin{bmatrix} I & 0 \\ 0 & X \end{bmatrix} \begin{bmatrix} I & 0 \\ 0 & -I \end{bmatrix} \begin{bmatrix} I & 0 \\ 0 & X \end{bmatrix} = \begin{bmatrix} I & 0 \\ 0 & -X \end{bmatrix} \begin{bmatrix} I & 0 \\ 0 & X \end{bmatrix} = \begin{bmatrix} I & 0 \\ 0 & -X^2 \end{bmatrix} = \begin{bmatrix} I & 0 \\ 0 & -I \end{bmatrix} = Z_1$$

$$UZ_2U^\dagger = \begin{bmatrix} I & 0 \\ 0 & X \end{bmatrix} \begin{bmatrix} Z & 0 \\ 0 & Z \end{bmatrix} \begin{bmatrix} I & 0 \\ 0 & X \end{bmatrix} = \begin{bmatrix} Z & 0 \\ 0 & XZ \end{bmatrix} \begin{bmatrix} I & 0 \\ 0 & X \end{bmatrix} = \begin{bmatrix} Z & 0 \\ 0 & XZX \end{bmatrix} = \begin{bmatrix} Z & 0 \\ 0 & -Z \end{bmatrix} = Z_1Z_2$$

### Exercise 10.37

What is  $UY_1U^{\dagger}$ ?

## Solution

Concepts Involved: Linear Algebra, Unitary Operators, Controlled Operations

Note from the previous Exercise that  $UX_1U^\dagger=X_1X_2$  and  $UZ_1U^\dagger=Z_1$ . Hence writing Y=iZX and using that  $U^\dagger U=I$ :

$$UY_1U^{\dagger} = iUZXU^{\dagger} = iUZU^{\dagger}UXU^{\dagger} = iZ_1X_1X_2 = iY_1X_2.$$

#### 

#### Exercise 10.38

Suppose U and V are unitary operators on two qubits which transform  $Z_1, Z_2, X_1$  and  $X_2$  by conjugation in the same way. Show this implies that U = V.

Concepts Involved: Linear Algebra, Unitary Operators, Pauli Group

We observe that  $UAU^\dagger=VAV^\dagger$  for  $A\in\{Z_1,Z_2,X_1,X_2\}$  implies that  $UAU^\dagger=VAV^\dagger$  for any two qubit Pauli as  $UAU^\dagger UBU^\dagger=VAV^\dagger VBV^\dagger\implies UABU^\dagger=VABV^\dagger$  for any two operators A,B, and  $\{Z_1,Z_2,X_1,X_2\}$  generates all two qubit Paulis via multiplication (which is seen from  $Z^2=I$  and ZX=iY). Then, observing that the two-qubit Pauli operators form a basis for operators acting on the Hilbert space  $\mathbb{C}^4$  (formally - they are 16 linearly independent vectors in  $\mathbb{C}^{4\times4}$  over field  $\mathbb{C}$ , and since  $\dim\mathbb{C}^{4\times4}=16$  they form a basis) we can write any operator  $O\in\mathbb{C}^{4\times4}$  as a linear combination of the two qubit Paulis,  $O=\sum_i c_i A_i$ . We then have:

$$UOU^{\dagger} = U\left(\sum_{i} c_{i} A_{i}\right) U^{\dagger} = \sum_{i} c_{i} U A_{i} U^{\dagger} = \sum_{i} c_{i} V A_{i} V^{\dagger} = V\left(\sum_{i} c_{i} A_{i}\right) V^{\dagger} = VOV^{\dagger}$$

So  $UOU^\dagger = VOV^\dagger$  for all  $O \in \mathbb{C}^{4 \times 4}$ , which is only possible if U = V.

### Exercise 10.39

Verify (10.91).

#### Solution

Concepts Involved: Linear Algebra, Unitary Operators

We verify by matrix multiplication:

$$SXS^{\dagger} = \begin{bmatrix} 1 & 0 \\ 0 & i \end{bmatrix} \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 0 & -i \end{bmatrix} = \begin{bmatrix} 0 & 1 \\ i & 0 \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 0 & -i \end{bmatrix} = \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} = Y$$

$$SZS^\dagger = \begin{bmatrix} 1 & 0 \\ 0 & i \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 0 & -i \end{bmatrix} = \begin{bmatrix} 1 & 0 \\ 0 & -i \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 0 & -i \end{bmatrix} = \begin{bmatrix} 1 & 0 \\ 0 & (-i)^2 \end{bmatrix} = \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} = Z$$

Provide an inductive proof of Theorem 10.6 as follows.

- (1) Prove that the Hadamard and phase gates can be used to perform any normalizer operation on a single qubit.
- (2) Suppose U is an n+1 qubit gate in  $N(G_{n+1})$  such that  $UZ_1U^\dagger=X_1\otimes g$  and  $UX_1U^\dagger=Z_1\otimes g'$  for some  $g,g'\in G_n$ . Define U' on n qubits by  $U'|\psi\rangle\equiv\sqrt{2}\langle 0|U(|0\rangle\otimes|\psi\rangle)$ . Use the inductive hypothesis to show that the construction for U in Figure 10.9 may be implemented using  $O(n^2)$  Hadamard, phase, and controlled-NOT gates.
- (3) Show that any gate  $U \in N(G_{n+1})$  may be implemented using  $O(n^2)$  Hadamard, pahse, and controlled-NOT gates.

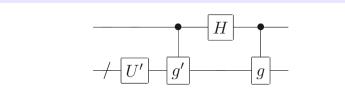


Figure 10.9. Construction used to prove that the Hadamard, phase and controlled-NOT gates generate the normalizer  $N(G_n)$ .

#### Exercise 10.41

Verify Equations (10.92) through (10.95).

#### Exercise 10.42

Use the stabilizer formalism to verify that the circuit of Figure 1.13 on page 27 teleports qubits, as claimed. Note that the stabilizer formalism restricts the class of states being teleported, so in some sense this is not a complete description of teleportation, nevertheless it does allow an understanding of the dynamics of teleportation.

#### Exercise 10.43

Show that  $S \subseteq N(S)$  for any subgroup S of  $G_n$ .

# Exercise 10.44

Show that N(S) = Z(S) for any subgroup S of  $G_n$  not containing -I.

# **Exercise 10.45: Correcting located errors**

Suppose C(S) is an [n,k,d] stabilizer code. Suppose k qubits are encoded in n qubits using this code, which is then subjected to noise. Fortunately, however, we are told that only d-1 of the qubits are affected by the noise, and moreover, we are told precisely which d-1 qubits have been affected. Show that it is possible to correct the effects of such *located* errors.

Show that the stabilizer for the three qubit pahse flip code is generated by  $X_1X_2$  and  $X_2X_3$ .

#### Exercise 10.47

Verify that the generators of Figure 10.11 generate the two codewords of Equation (10.13).

Name	Operator
$g_1$	ZZIIIIIII
$g_2$	IZZIIIIII
$g_3$	IIIZZIIII
$g_4$	IIIIZZIII
$g_5$	IIIIIIZZI
$g_6$	IIIIIIIZZ
$g_7$	XXXXXX I I I
$g_8$	IIIXXXXXX
$ar{Z}$	XXXXXXXXX
$\bar{X}$	ZZZZZZZZZ

Figure 10.11. The eight generators for the Shor nine qubit code, and the logical Z and logical X operations. (Yes, they really are the reverse of what one might naively expect!)

### Exercise 10.48

Show that the operations  $\bar{Z}=X_1X_2X_3X_4X_5X_6X_7X_8X_9$  and  $\bar{X}=Z_1Z_2Z_3Z_4Z_5Z_6Z_7Z_8Z_9$  act as logical Z and X operations on a Shor-code encoded qubit. That is, show that this  $\bar{Z}$  is independent of and commutes with the generators of the Shor code, and that  $\bar{X}$  is independent of and commutes with the generators of the Shor code, and anti-commutes with  $\bar{Z}$ .

### Exercise 10.49

Use Theorem 10.8 to veirfy that the five qubit code can protect against an arbitrary single qubit error.

## Exercise 10.50

Show that the five qubit code saturates the quantum Hamming bound, that is, it satisfies the inequality of (10.51) with equality.

### Exercise 10.51

Verify that the check matrix defined in (10.106) corresponds to the stabilizer of the CSS code  $CSS(C_1,C_2)$ , and use Theorem 10.8 to show that arbitrary errors on up to t qubits may be corrected by this code.

# Exercise 10.52

Verify by direct operation on the codewords that the operators of (10.107) act appropriately, as logical Z and X.

Prove that the encoded Z operators are independent of one another.

#### Exercise 10.54

Prove that with the check matrix for the encoded X operators defined as above, the encoded X operators are idnependent of one another and of the gnerators, commute with the gnerators of the stabilizer, with each other, and  $\bar{X}_j$  commutes with all the  $\bar{Z}_k$  except with  $\bar{Z}_j$ , with which it anti-commutes.

### Exercise 10.55

Find the  $\bar{X}$  operator for the standard form of the Steane code.

### Exercise 10.56

Show that replacing an encoded X or Z operator by g times that operator, where g is an element of the stabilizer, does not change the action of the operator on the code.

### Exercise 10.57

Give the check matrices for the five and nine qubit codes in standard form.

### Exercise 10.58

Verify that the circuits in Figures 10.13–10.15 work as described, and check the claimed circuit equivalences.

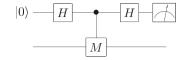


Figure 10.13. Quantum circuit for measuring a single qubit operator M with eigenvalues  $\pm 1$ . The top qubit is the ancilla used for the measurement, and the bottom qubit is being measured.



Figure 10.14. Quantum circuit for measuring the X operator. Two equivalent circuits are given; the one on the left is the usual construction (as in Figure 10.13), and the one on the right is a useful equivalent circuit.



Figure 10.15. Quantum circuit for measuring the Z operator. Two equivalent circuits are given; the one on the left is the usual construction (as in Figure 10.13), and the one on the right is a useful simplification.

Show that by using the identities of Figures 10.14 and 10.15, the syndrome circuit of Figure 10.16 can be replaced with the circuit of Figure 10.17.

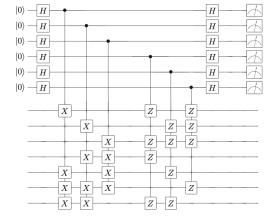


Figure 10.16. Quantum circuit for measuring the generators of the Steane code, to give the error syndrome. The top six qubits are the ancilla used for the measurement, and the bottom seven are the code qubits.

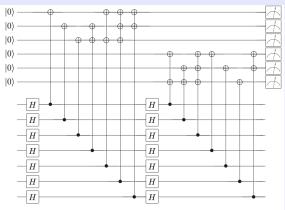


Figure 10.17. Quantum circuit equivalent to the one in Figure 10.16.

### Exercise 10.60

Construct a syndrome measuring circuit analogous to that in Figure 10.16, but for the nine and five qubit codes.

### Exercise 10.61

Describe explicit recovery operations  $E_j^{\dagger}$  corresponding to the different possible error syndromes that may be measured using the circuit in Figure 10.16.

### Exercise 10.62

Show by explicit construction of generators for the stabilizer that concatenating an  $[n_1, 1]$  stabilizer code with an  $[n_2, 1]$  stabilizer code gives an  $[n_1n_2, 1]$  stabilizer code.

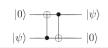
# Exercise 10.63

Suppose U is any unitary operation mapping the Steane code into itself, and such that  $U\bar{Z}U^\dagger=\bar{X}$  and  $U\bar{X}U^\dagger=\bar{Z}$ . Prove that up to a global phase the action of U on the encoded states  $|0_L\rangle$  and  $|1_L\rangle$  is  $|0_L\rangle\mapsto (|0_L\rangle+|1_L\rangle)/\sqrt{2}$  and  $|1_L\rangle\mapsto (|0_L\rangle-|1_L\rangle)/\sqrt{2}$ .

#### **Exercise 10.64: Back propogation of errors**

It is clear that an X error on the control qubit of a CNOT gate propagates to the target qubit. In addition, it turns out that a Z error on the target propagates back to the control! Show this using the stabilizer formalism, and also directly using quantum circuit identities. You may find Exercise 4.20 on page 179 useful.

An unknown qubit in the state  $|\psi\rangle$  can be swapped with a second qubit which is prepared in the state  $|0\rangle$  using only two controlled-NOT gates, with the circuit

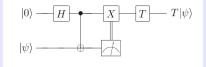


Show that the two circuits below, which use only a single CNOT gate, with measurement and a classically controlled single qubit operation, also accomplish the same task:



#### Exercise 10.66

One way to implement a  $\pi/8$  gate is to first swap the qubit state  $|\psi\rangle$  you wish to transform with some unknown state  $|0\rangle$ , then to apply a  $\pi/8$  gate to the resulting qubit. Here is a quantum circuit which does that:



Doing this does not seem particularly useful, but actually it leads to something which is! Show that by using the relations  $TX = \exp(-i\pi/4)SX$  and TU = UT (U is the controlled-NOT gate, and T acts on the control qubit) we may obtain the circuit of Figure 10.25.

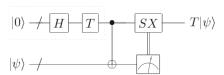
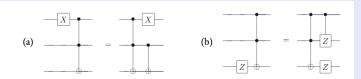


Figure 10.25. Quantum circuit which fault-tolerantly implements a  $\pi/8$  gate. The dashed box represents a (non-fault-tolerant) preparation procedure for the ancilla state  $(|0\rangle+\exp(i\pi/4)|1\rangle)/\sqrt{2}$ ; how to do this preparation fault-tolerantly is explained in the text. The slash on the wire denotes a bundle of seven qubits, and the double-line wire represents the classical bit resulting from the measurement. Note that the final SX operation is controlled by the measurement result.

## Exercise 10.67

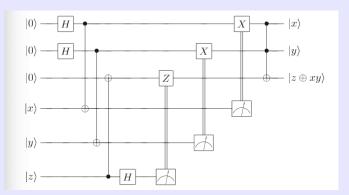
Show that the following circuit identities hold:



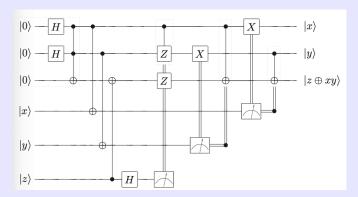
### Exercise 10.68: Fault-tolerant Toffoli gate construction

A procedure similar to the above sequence of exercises for the  $\pi/8$  gate gives a fault-tolerant Toffoli gate.

(1) First, swap the three qubit state  $|xyz\rangle$  you wish to transform with some known state  $|000\rangle$ , then apply a Toffoli gate to the resulting qubits. Show that the following circuit accomplishes this task:



(2) Using the commutation rules from Exercise 10.67, show that moving the final Toffoli gate all the way back to the left side gives the circuit



(3) Assuming the ancilla preparation shown in the leftmost dotted box can be done fault-tolerantly, show that this circuit can be used to give a fault-tolerant implementation of the Toffoli gate using the Steane code.

### Exercise 10.69

Show that a single failure anywhere in the ancilla preparation and verification can lead to at most one X or Y error in the ancilla output.

### Exercise 10.70

Show that Z errors in the ancilla do not propagate to affect the encoded data, but result in an incorrect measurement result being observed.

Verify that when  $M=e^{-i\pi/4}SX$  the procedure we have described gives a fault-tolerant method for measuring M.

### Exercise 10.72: Fault-tolerant Toffoli ancilla state construction

Show how to fault-tolerantly prepare the state created by the circuit in the dotted box of exercise 10.68, that is,

$$\frac{|000\rangle + |010\rangle + |100\rangle + |111\rangle}{2}$$

You may find it helpful to first give the stabilizer generators for this state.

#### Exercise 10.73: Fault-tolerant encoded state construction

Show that the Steane code encoded  $|0\rangle$  state can be constructed fault-tolerantly in the following manner.

- (1) Begin with the circuit of Figure 10.16, and replace the measurement of each generator, as shown in Figure 10.30, with each ancilla qubit becoming a cat state  $|00\dots0\rangle + |11\dots1\rangle$ , and the operations rearranged to have their controls on different qubits, so that errors do not propagate within the code block.
- (2) Add a stage to fault-tolerantly measure Z.
- (3) Calculate the error probability of this circuit, and of the circuit when the generator measurements are repeated three times and majority voting is done.
- (4) Enumerate the operations which should be performed conditioned on the measurement results and show that they can be done fault-tolerantly.

### Exercise 10.74

Construct a quantum circuit to fault-tolerantly generate the encoded  $|0\rangle$  state for the five qubit code (Section 10.5.6).

### Problem 10.1

Channels  $\mathcal{E}_1$  and  $\mathcal{E}_2$  are said to be *equivalent* if there exist unitary channels  $\mathcal{U}$  and  $\mathcal{V}$  such that  $\mathcal{E}_2 = \mathcal{U} \circ \mathcal{E}_1 \circ \mathcal{V}$ .

- (1) Show that the relation of channel equivalence is an equivalence relation.
- (2) Show how to turn an error-correcting code for  $\mathcal{E}_1$  into an error-correcting code for  $\mathcal{E}_2$ . Assume that the error-correction procedure for  $\mathcal{E}_1$  is performed as a projective measurement followed by a conditional unitary operation, and explain how the error-correction procedure for  $\mathcal{E}_2$  can be performed in the same fashion.

**Concepts Involved:** Unitary operators, Quantum Channels, Quantum Measurement. Recall that a unitary channel is a quantum channel such that  $\mathcal{U}(\rho) = U\rho U^{\dagger}$  for a unitary operator U.

(1) Reflexivity. Take  $\mathcal{U} = \mathcal{V} = \mathcal{I}$  (the identity channel, which is unitary). Then  $\mathcal{E} = \mathcal{I} \circ \mathcal{E} \circ \mathcal{I}$  so  $\mathcal{E} \sim \mathcal{E}$  as desired.

Symmetry. Suppose  $\mathcal{E}_1 \sim \mathcal{E}_2$ . Then there exists  $\mathcal{U}, \mathcal{V}$  such that  $\mathcal{E}_2 = \mathcal{U} \circ \mathcal{E}_1 \circ \mathcal{V}$ . It then follows that  $\mathcal{E}_1 = \mathcal{U}^\dagger \circ \mathcal{E}_2 \circ \mathcal{V}^\dagger$ , where  $\dagger$  denotes the dual channel, i.e. if  $\mathcal{U}(\rho) = U \rho U^\dagger$  then  $\mathcal{U}^\dagger(\rho) = U^\dagger \rho U$ . The dual channel of a unitary channel is also a unitary channel as  $U^\dagger$  is unitary for unitary U. Hence  $\mathcal{E}_2 \sim \mathcal{E}_1$  as desired.

Transitivity. Suppose  $\mathcal{E}_1 \sim \mathcal{E}_2$  and  $\mathcal{E}_2 \sim \mathcal{E}_3$ . Then there exist  $\mathcal{U}, \mathcal{V}, \mathcal{U}', \mathcal{V}'$  such that  $\mathcal{E}_2 = \mathcal{U} \circ \mathcal{E}_1 \circ \mathcal{V}$  and  $\mathcal{E}_3 = \mathcal{U}' \circ \mathcal{E}_2 \circ \mathcal{V}'$ . It then follows that  $\mathcal{E}_3 = \mathcal{U}' \circ \mathcal{U} \circ \mathcal{E}_1 \circ \mathcal{V} \circ \mathcal{V}'$  and since the composition of two unitary channels is another unitary channel (as the composition of two unitary operators is again unitary), it follows that  $\mathcal{E}_1 \sim \mathcal{E}_3$  as desired.

We have therefore shown that channel equivalence is an equivalence relation.

(2) Suppose that  $\mathcal{E}_2 = \mathcal{U} \circ \mathcal{E}_1 \circ \mathcal{V}$ . If the error correcting code for  $\mathcal{E}_1$  consists of a projective measurement  $P_m$  followed by a conditional unitary  $C_m$  (where the subscript denotes a conditioning on the measurement outcome m), this means that for a given state  $\rho$  which lies in the support of the code that  $C_m P_m \mathcal{E}_1(\rho) P_m C_m^\dagger \propto \rho$ .

If we instead have  $\mathcal{E}_2$  as our noise channel, let us rotate the measurement basis and carry out the measurement  $P'_m = U P_m U^\dagger$  (note that a unitary conjugation of a projector remains a projector, as  $P'^2_m = (U P_m U^\dagger)^2 = U P_m U^\dagger U P_m U^\dagger = U P_m U^\dagger = P'_m$ ). Let us also modify the conditional unitary to be  $C'_m = V^\dagger C_m U^\dagger$ . We then have:

$$\begin{split} C'_{m}P'_{m}\mathcal{E}_{2}(\rho)P'_{m}C'^{\dagger}_{m} &= (V^{\dagger}C_{m}U^{\dagger})(UP_{m}U^{\dagger})U\mathcal{E}_{1}(V\rho V^{\dagger})U^{\dagger}(UP_{m}U^{\dagger})(V^{\dagger}C_{m}U^{\dagger})^{\dagger} \\ &= V^{\dagger}C_{m}U^{\dagger}UP_{m}U^{\dagger}U\mathcal{E}_{1}(V\rho V^{\dagger})U^{\dagger}UP_{m}U^{\dagger}UC^{\dagger}_{m}V \\ &= V^{\dagger}C_{m}P_{m}\mathcal{E}_{1}(V\rho V^{\dagger})P_{m}C^{\dagger}_{m}V \\ &\propto V^{\dagger}V\rho V^{\dagger}V \\ &= \rho \end{split}$$

which shows that the error-correction procedure for  $\mathcal{E}_2$  can be performed in the same fashion (with the unitary-modified measurements/unitaries).

### Problem 10.2: Gilbert-Varshamov bound

Prove the Gilbert-Varshamov bound for CSS codes, namely, that an [n,k] CSS code correcting t errors exists for some k such that

$$\frac{k}{n} \ge 1 - 2H\left(\frac{2t}{n}\right).$$

As a challenge, you may like to try proving the Gilbert-Varshamov bound for a general stabilizer code,

namely, that there exists an [n, k] stabilizer code correcting errors on t qubits, with

$$\frac{k}{n} \ge 1 - \frac{2\log(3)t}{n} - H\left(\frac{2t}{n}\right)$$

### Problem 10.3: Encoding stabilizer codes

Suppose we assume that the generators for the code are in standard form, and that the encoded Z and X operators have been constructed in standard form. Find a circuit taking the  $n \times 2n$  check matrix corresponding to a listing of all the generators for the code together with the encoded Z operations from

$$G = \begin{bmatrix} 0 & 0 & 0 & | & I & 0 & 0 \\ 0 & 0 & 0 & | & 0 & I & 0 \\ 0 & 0 & 0 & | & 0 & 0 & 0 \end{bmatrix}$$

to the standard form

$$\begin{bmatrix} I & A_1 & A_2 & B & 0 & C_2 \\ 0 & 0 & 0 & D & I & E \\ 0 & 0 & 0 & A_2^T & 0 & I \end{bmatrix}$$

### Problem 10.4: Encoding by teleportation

Suppsoe you are given a qubit  $|\psi\rangle$  to encode in a stabilizer code, but you are not told anything about how  $|\psi\rangle$  was constructed: it is an unknown state; Construct a circuit to perform the encoding in the following manner:

(1) Explain how to fault-tolerantly construct the partially encoded state

$$\frac{|0\rangle|0_L\rangle+|1\rangle|1_L\rangle}{\sqrt{2}},$$

by writing this as a stabilizer state, so it can be prepared by measuring stabilizer generators.

- (2) Show how to fault-tolerantly perform a Bell basis measurement with  $|\psi\rangle$  and the unencoded qubit from this state.
- (3) Give the Pauli operations which you need to fix up the remaining encoded qubit after this measurement, so that it becomes  $|\psi\rangle$ , as in the usual quantum teleportation scheme.

Compute the probability of error of this circuit. Also show how to modify the circuit to perform fault-tolerant decoding.

#### Problem 10.5

Suppose C(S) is an [n,1] stabilizer code capable of correcting errors on a single qubit. Explain how a fault-tolerant implementation of the controlled-NOT gate may be implemented between two logical qubits encoded using this code, using only fault-tolerant stabilizer state preparation, fault-tolerant measurement of elements of the stabilizer, and normalizer gates applied transversally.

# A1 Notes on basic probability theory

### Exercise A1.1

Prove Bayes' rule.

### Solution

Concepts Involved: Probability, Conditional Probability.

Recall that conditional probabilities were defined as:

$$P(Y = y | X = x) = \frac{P(X = x, Y = y)}{P(X = x)}$$

and also recall that Bayes' rule is given by:

$$p(x|y) = p(y|x)\frac{p(x)}{p(y)}$$

By the definition of conditional probability:

$$p(y|x)\frac{p(x)}{p(y)} = \frac{p(X = x, Y = y)}{p(x)} \frac{p(x)}{p(y)} = \frac{P(X = x, Y = y)}{p(y)} = p(x|y)$$

### Exercise A1.2

Prove the law of total probability.

#### Solution

Concepts Involved: Probability, Conditional Probability.

Recall that the law of total probability is given by:

$$p(y) = \sum_{x} p(y|x)p(x)$$

Using the identity  $p(Y=y) = \sum_x p(X=x,Y=y)$  and Bayes' rule, we have

$$p(y) = \sum_{x} p(x, y) = \sum_{x} p(y|x)p(x)$$

# Exercise A1.3

Prove that there exists a value of  $x \ge \mathbf{E}(X)$  such that p(x) > 0.

Concepts Involved: Probability, Expectation.

Recall that the expectation of a random variable X is defined by:

$$\mathbf{E}(X) = \sum_{x} p(x)x$$

Let  $\tilde{x} = \max\{x : x \text{ is a possible value of } X\}$ . This maximum exists as we assume X can only take on a finite set of values. We therefore have that:

$$\mathbf{E}(X) = \sum_{x} p(x)x \le \sum_{x} p(x)\tilde{x} = \tilde{x} \sum_{x} p(x) = \tilde{x}$$

Where in the last equality we use that the sum over all probabilities must be 1.

### Exercise A1.4

Prove that  $\mathbf{E}(X)$  is linear in X

### Solution

Concepts Involved: Probability, Expectation.

Let  $a,b\in\mathbb{R}$  and X,Y be random variables. We then have that:

$$\mathbf{E}(aX + bY) = \sum_{x} \sum_{y} p(x, y)(ax + by)$$

$$= \sum_{x} \sum_{y} p(x, y)ax + \sum_{x} \sum_{y} p(x, y)by$$

$$= a \sum_{x} \left(\sum_{y} p(x, y)\right) x + b \sum_{y} \left(\sum_{x} p(x, y)\right) y$$

$$= a \sum_{x} p(x)x + b \sum_{y} p(y)y$$

$$= a\mathbf{E}(X) + b\mathbf{E}(y)$$

which shows that expectation is linear.

### Exercise A1.5

Prove that for independent random variables X and Y,  $\mathbf{E}(XY) = \mathbf{E}(X)\mathbf{E}(Y)$ .

### Solution

Concepts Involved: Probability, Expectation, Independent Random Variables.

Recall two random variables X, Y are independent if

$$p(X=x,Y=y) = p(X=x)p(Y=y)$$

We have that:

$$\mathbf{E}(XY) = \sum_{x} \sum_{y} xyp(x,y) = \sum_{x} \sum_{y} xyp(x)p(y) = \left(\sum_{x} p(x)x\right) \left(\sum_{y} p(y)y\right) = \mathbf{E}(x)\mathbf{E}(y)$$

Exercise A1.6

(\*) Prove Chebyshev's inequality.

#### Solution

Concepts Involved: Probability, Expectation, Variance.

Recall the definition of the variance and standard deviaiton of a random variable X:

$$Var(X) = \mathbf{E}[(X - \mathbf{E}(X))^2], \quad \Delta(X) = \sqrt{Var(X)}$$

Also, recall that Chebyshev's inequality reads:

$$p(|X - \mathbf{E}(X)| \ge \lambda \Delta(X)) \le \frac{1}{\lambda^2}$$

where  $\lambda > 0$ .

We first establish Markov's inequality for the expectation value  $\mathbf{E}(X)$ . Let a>0, and then we have that:

$$\mathbf{E}(X) = \sum_{x} xp(x) = \sum_{x>a} xp(x) + \sum_{xa} ap(x) + 0 = ap(X \ge a)$$

Therefore, we obtain that:

$$p(X \ge a) \le \frac{\mathbf{E}(X)}{a}$$

for any random variable X and a>0. Next, substitute X with  $(X-\mathbf{E}(X))^2$  and let  $a=\lambda^2\operatorname{Var}(X)$  for  $\lambda>0$ . Markov's inequality then states that:

$$p((X - \mathbf{E}(X))^2 \ge \lambda^2 \operatorname{Var}(X)) \le \frac{\mathbf{E}(X - \mathbf{E}(X))^2}{\lambda^2 \operatorname{Var}(X)}$$

Since  $\mathbf{E}(X - \mathbf{E}(X))^2 = \mathrm{Var}(X)$ , we have that:

$$p((X - \mathbf{E}(X))^2 \ge \lambda^2 \operatorname{Var}(X)) \le \frac{1}{\lambda^2}$$

If  $\lambda > 0$ , then  $p((X - \mathbf{E}(X))^2 \ge \lambda^2 \operatorname{Var}(X)) = p(|X - \mathbf{E}(X)| \ge \lambda \Delta(X))$  by taking square roots, so we obtain:

$$p(\left|X - \mathbf{E}(X)\right| \ge \lambda \Delta(X)) \le \frac{1}{\lambda^2}$$

as desired.  $\Box$ 

# A2 Group theory

### Exercise A2.1

Prove that for any element of g of a finite group, there always exists a positive integer r such that  $g^r = e$ . That is, every element of such a group has an order.

#### Solution

Concepts Involved: Group Axioms, Order.

Suppose G is a finite group, and  $g \in G$ . Then, there exists some  $r_1, r_2 \in \mathbb{N}$  such that  $r_1 \neq r_2$  and  $g^{r_1} = g^{r_2}$ . If this was not the case, then  $g^n$  would be unique for each  $n \in \mathbb{N}$ , contradicting the finiteness of G. WLOG take  $r_1 < r_2$ , and let  $r = r_2 - r_1 \in \mathbb{N}$ . Using associativity, we then have that:

$$g^{r_1} = g^{r_2} = g^{r_1+r} = g^{r_1}g^r$$

from which we conclude that  $g^r = e$ .

#### Exercise A2.2

Prove Lagrange's Theorem.

#### Solution

Concepts Involved: Group Axioms, Subgroups, Order, Equivalence Relations.

Let H be a subgroup of a group G and define the relation  $\sim$  by  $a \sim b$  iff a = bh for some  $h \in H$ .  $\sim$  is reflexive as a = ae (so  $a \sim a$ ) where  $e \in H$  is the identity element.  $\sim$  is symmetric as if  $a \sim b$ , then a = bh for some  $h \in H$  so  $b = ah^{-1}$  (so  $b \sim a$ ) where  $h^{-1} \in H$  as H is closed under inverses. Finally  $\sim$  is transitive as if  $a \sim b$  and  $b \sim c$ , there exist  $h_1, h_2 \in H$  such that  $a = bh_1$  and  $b = ch_2$  so  $a = ch_2h_1$ . As  $h_2h_1 \in H$  (H is closed under multiplication) it follows that  $a \sim c$ . Having shown  $\sim$  to have these three properties, we conclude it is an equivalence relation. Then, the equivalence classes of  $\sim$  partition G, where the equivalence class of  $g \in G$  is  $[g] = \{gh|h \in H\}$ .

Now, let  $g \in G$  and define the map  $\varphi_g : H \to [g]$  as  $\varphi_g(h) = gh$ .  $\varphi_g$  is injective as if  $\varphi_g(h_1) = \varphi_g(h_2)$  then  $gh_1 = gh_2$  and multiplying by  $g^{-1}$  on both sides  $h_1 = h_2$ .  $\varphi_g$  is surjective as if  $k \in [g]$ , then there exists some  $h \in H$  such that k = gh by the definition of  $\sim$ . Hence  $\varphi_g$  is bijective.

As per our prior observation, the equivalence classes of  $\sim$  partition G, so  $G = \bigcup_{i=1}^n [g_i]$  and  $|G| = |\bigcup_{i=1}^n [g_i]| = \sum_{i=1}^n |[g_i]|$ . Further, there is a bijection  $\varphi_{g_i}$  from each equivalence class to H, so  $|[g_i]| = |H|$  for all i. Thus  $|G| = \sum_{i=1}^n |H| = nH$  and hence |H| divides |G|, as desired.  $\square$ 

#### Exercise A2.3

Show that the order of an element  $g \in G$  divides |G|.

Concepts Involved: Group Axioms, Subgroups, Order, Lagrange's Theorem.

Let  $g\in G$  with order r. Then, define  $H=\left\{g^n|n\in\mathbb{N}\right\}$ . We claim that H is a subgroup of G. First,  $g^n\in G$  for any n as G is closed under multiplication, so  $H\subset G$ . Next, if  $g^{n_1},g^{n_2}\in H$  then  $g^{n_1}\cdot g^{n_2}=g^{n_1+n_2}\in H$ . Associativity is inherited from the associativity of multiplication in G. Since  $g^r=e\in H$ , H contains the identity. Finally, for  $g^k\in H$  we have  $g^{r-k}\in H$  such that  $g^kg^{r-k}=g^{r-k}g^k=g^r=e$  so H is closed under inverses. Hence the claim is proven.

Next, we observe that |H|=r as H contains the r elements  $e,g,g^2,\ldots g^{r-1}$ . Hence by Lagrange's Theorem r divides |G|.

### Exercise A2.4

Show that if  $y \in G_x$  then  $G_y = G_x$ 

#### Solution

Concepts Involved: Group Axioms, Conjugacy Classes

Suppose  $y \in G_x$ . Then there exists some  $g \in G$  such that  $g^{-1}xg = y$ . Multiplying both sides on the left by g and on the right by  $g^{-1}$  we find that  $x = gyg^{-1}$ . We now show the two inclusions.

Suppose that  $k \in G_x$ . Then there exists some  $g' \in G$  such that  $k = g'^{-1}xg'$ . Then using  $x = gyg^{-1}$  we find  $k = g'^{-1}gyg^{-1}g'$ . Now,  $g^{-1}g' \in G$  (by closure) and it has inverse  $g'^{-1}g$ , and hence  $k = g'^{-1}gyg^{-1}g' \in G_y$ . So,  $G_x \subseteq G_y$ .

Suppose that  $l \in G_y$ . Then there exists some  $g'' \in G$  such that  $l = g''^{-1}yg''$ . Then with  $g^{-1}xg = y$  we find  $l = g''^{-1}g^{-1}xgg''$ . Much like before,  $gg'' \in G$  (by closure) with inverse  $g''^{-1}g^{-1}$  so  $l \in G_x$ . So,  $G_y \subseteq G_x$ .

We conclude that  $G_y = G_x$ .

### Exercise A2.5

Show that if x is an element of an Abelian group G, then  $G_x = \{x\}$ .

#### Solution

Concepts Involved: Abelian Groups, Conjugacy Classes.

Evidently  $x=e^{-1}xe\in G_x$  so  $\{x\}\subseteq G_x$ . Next, if  $k\in G_x$  then  $k=g^{-1}xg$  for some  $g\in G$ , but since G is abelian,  $g^{-1}x=xg^{-1}$  so  $k=xg^{-1}g=xe=x$  so  $k\in \{x\}$  and hence  $G_x\subseteq \{x\}$ . We conclude that  $G_x=\{x\}$ .

# Exercise A2.6

Show that any group of prime order is cyclic.

Concepts Involved: Order, Cyclic Groups.

Suppose |G|=p where p is prime. Since G is finite, every element of G has an order by Exercise A2.1. Since the order of any element  $g\in G$  divides |G|=p by Exercise A2.3, and since p is prime, the order of g is either 1 or p. Since |G|>1, there exists at least one  $g\in G$  with order p, and this g is a generator of G (with  $g^1=g,g^2,g^3,\ldots,g^p=e$  distinct and comprising all the elements of G. In fact this is true of any non-identity g). Hence G is cyclic.  $\Box$ 

#### Exercise A2.7

Show that every subgroup of a cyclic group is cyclic.

### Solution

Concepts Involved: Group Axioms, Subgroups, Cyclic Groups, Euclid's Division Algorithm.

First we prove a necessary Lemma, namely that any nonempty subset of the natural numbers contains a least element. We show this by proving the contrapositive. Suppose that  $A\subseteq\mathbb{N}$  has no least element. Then  $1\notin A$  as then 1 would be the least element. Suppose then that  $1,\ldots k-1\notin A$ ; then  $k\notin A$  as then k would be the least element. By strong induction, there exists no  $k\in\mathbb{N}$  such that  $k\in A$ , i.e.  $k\in A$  is empty. This concludes the proof of the lemma.

Let  $G=\langle a \rangle$  be a cyclic group and H a subgroup of G. If  $H=\{e\}$ , it is trivially cyclic and we are done. If  $H \neq \{e\}$ , then there exists some  $a^n \in H$  with  $n \neq 0$ . Since H is closed under inverses,  $(a^n)^{-1}=a^{-n}\in H$  as well which ensures that H contains some positive power of a. Then consider the set  $A=\{k\in \mathbb{N}|a^k\in H\}$ . Any nonempty subset of the naturals has a minimum element; therefore let  $d=\min A$ . It is immediate that  $\langle a^d\rangle$  is a subgroup of H as  $a^d\in H$  and H is a group. To show the reverse containment, suppose that  $g\in H$ . Since H is a subgroup of the cyclic G, it follows that  $g=a^p$  for some  $p\in \mathbb{Z}$ . We can then write p=qd+r for  $0\leq r< d$  by Euclid's Division algorithm (see Appendix 4). We then have that  $a^r=a^{p-qd}=a^p(a^d)^{-q}\in H$  by closure. Now, since d is the least positive integer for which  $a^d\in H$  and  $0\leq r< d$ , it must follow that r=0. Therefore, p=qd and hence  $a^{qd}=(a^d)^q\in \langle a^d\rangle$ . So, H is a subgroup of  $\langle a^d\rangle$ . We conclude that  $H=\langle a^d\rangle$  and hence H is cyclic.

### Exercise A2.8

Show that if  $g \in G$  has finite order r, then  $g^m = g^n$  if and only if  $m = n \pmod{r}$ .

### Solution

Concepts Involved: Order, Modular Arithmetic, Euclid's Division Algorithm

Suppose  $g \in G$  has finite order r.

First suppose that  $m=n\pmod r$ . Then m-n=kr for some  $k\in\mathbb{N}$ . Therefore  $g^{m-n}=g^{kr}$ . But  $g^{kr}=(g^r)^k=e^k=e$ , so  $g^{m-n}=g^mg^{-n}=e$ , and multiplying both sides by  $g^n$  we find  $g^m=g^n$ . Suppose  $g^m=g^n$ . Then multiplying both sides by  $g^{-n}$  we find  $g^{m-n}=e$ . By Euclid's Division algorithm there exist integers q,p such that m-n=qr+p with  $0\le p< r$ . We then have that

 $g^{m-n}=g^{qr+p}=g^{qr}g^p=e. \text{ Furthermore, } g^{qr}=(g^r)^q=e^q=e \text{ so } g^p=e. \text{ But since } g \text{ has order } r \text{ and } 0\leq p< r, \text{ it follows that } p=0. \text{ Hence } m-n=qr \text{ and so } m\equiv n \pmod{r}.$ 

### Exercise A2.9

Cosets define an equivalence relation between elements. Show that  $g_1, g_2 \in G$  are in the same coset of H in G if and only if there exists some  $h \in H$  such that  $g_2 = g_1 h$ .

#### Solution

Concepts Involved: Equivalence Relations, Cosets

In Exercise A2.2 we showed that the relation  $\sim$  on a group G defined by  $g_1 \sim g_2$  iff  $g_1 = g_2 h$  for some  $h \in H$  was an equivalence relation. The equivalence classes of this equivalence relation were  $\left\{gh|h \in H\right\}$ , i.e. precisely the left cosets of H in G. So,  $g_1,g_2$  are in the same coset of H in G if and only if  $g_1 = g_2 h$  for some  $h \in H$ , which is exactly what we wished to show.  $\square$ 

### Exercise A2.10

How many cosets of H are there in G?

#### Solution

Concepts Involved: Equivalence Relations, Cosets

We observe that the map  $\varphi_g: H \to [g]$  defined in the solution of Exercise A2.2 is a map from H to a right coset of H in G defined by g. Since we showed that this map was bijective, this shows that |H| = |Hg| for any  $g \in G$ . Furthermore, since the cosets define an equivalence relation between elements of G, the cosets of H in G partition G. So, we conclude that there are |G|/|H| cosets of H in G, each of cardinality |H|.

### Exercise A2.11: Characters

Prove the properties of characters given above.

### Solution

Concepts Involved: Matrix Groups, Character (Trace)

Recall that the character of a matrix group  $G \subset M_n$  is a function on the group defined by  $\chi(g) = \operatorname{tr}(g)$  where  $\operatorname{tr}$  is the trace function. It has the properties that (1)  $\chi(I) = n$ , (2)  $\left|\chi(g)\right| \leq n$ , (3)  $\left|\chi(g)\right| = n$  implies  $g = e^{i\theta}I$ , (4)  $\chi$  is constant on any given conjugacy class of G, (5)  $\chi(g^{-1}) = \chi^*(g)$  and (6)  $\chi(g)$  is an algebraic number for all g.

The six properties are proven below.

- (1)  $\chi(I) = \text{tr}(I) = \sum_{k=1}^{n} 1 = n$ .
- (2) Let  $g \in G$ . Since G is finite, by Exercise A2.1 it follows that g has order r such that  $g^r = I$ . So, g may be diagonalized with roots of unity  $e^{2\pi i j/r}$ ,  $j \in \{0, 1, \dots, r-1\}$  on the diagonal. We then

find using the triangle inequality that:

$$|\chi(g)| = |\operatorname{tr}(g)| = \left| \sum_{k=1}^{n} e^{2\pi i j_k/r} \right| \le \sum_{k=1}^{n} \left| e^{2\pi i j_k/r} \right| = \sum_{i} 1 = n$$

which proves the claim.

- (3) The (complex) triangle inequality  $|z_1+z_2| \leq |z_1|+|z_2|$  is saturated when  $z_1=kz_2$  for some  $k\geq 0$ . This can only occur in the above equation when every  $\lambda_i$  in the sum is identical (as distinct roots of unity are not related by a non-negative constant). If the  $\lambda_i$ s are identical, then g is diagonal with diagonal entries of unit modulus, so  $g=e^{i\theta}I$  as claimed.
- (4) Let  $G_x=\left\{g^{-1}xg|g\in G\right\}$  be the conjugacy class of x in G. We then have for any  $h\in G_x$  that  $\chi(h)=\chi(g^{-1}xg)=\operatorname{tr}\left(g^{-1}xg\right)=\operatorname{tr}\left(xgg^{-1}\right)=\operatorname{tr}(xI)=\operatorname{tr}(x)$ , using the cyclicity of the trace. We conclude that  $\chi$  is constant on the conjugacy class.
- (5) By the same argument as (2),  $g \in G$  can be diagonalized with roots of unity  $e^{2\pi i j/r}$  on the diagonal:

$$g = \begin{bmatrix} e^{2\pi i j_1/r} & & & \\ & e^{2\pi i j_2/r} & & & \\ & & \ddots & & \\ & & & e^{2\pi i j_n/r} \end{bmatrix}$$

It then follows that  $g^{-1}$  is:

$$g^{-1} = \begin{bmatrix} e^{-2\pi i j_1/r} & & & \\ & e^{-2\pi i j_2/r} & & \\ & & \ddots & \\ & & & e^{-2\pi i j_n/r} \end{bmatrix}$$

So we have that:

$$\chi(g^{-1}) = \operatorname{tr}\left(g^{-1}\right) = \sum_{j=1}^{n} e^{-2\pi i j_k/r} = \sum_{j=1}^{n} (e^{2\pi i j_k/r})^* = \left(\sum_{j=1}^{n} e^{2\pi i j_k/r}\right)^* = (\operatorname{tr}(g))^* = \chi^*(g).$$

which proves the claim.

(6)  $\chi(g)$  is the sum of r-th roots of unity, which are algebraic; hence  $\chi(g)$  is algebraic as the sum of algebraic numbers.

### Exercise A2.12: Unitary matrix groups

(\*) A unitary matrix group is comprised solely of unitary matrices (those who which satisfy  $U^{\dagger}U=I$ ). Show that every matrix group is equivalent to a unitary matrix group. If a representation of a group consists entirely of unitary matrices, we may refer to it as being a *unitary representation*.

Concepts Involved: Matrix Groups, Character (Trace), Equivalence, Unitary Operators.

Recall that two groups are equivalent if they are isomorphic (i.e. there is a bijection between the groups that respects the group multiplication) and the isomorphic element have the same character.

Let  $G = \{A_1, \dots A_n\}$  be a finite matrix group. Then define

$$A = \sum_{i=1}^{n} A_i^{\dagger} A_i.$$

By Ex. 2.25 each term of the above sum is positive, and by Ex. 2.24 each term is Hermitian. The sum of Hermitian operators is Hermitian, so A is Hermitian. By Ex. 2.21, A is diagonalizable. Let U be the unitary matrix that diagonalizes A. We then have that D is a diagonal matrix, with:

$$D = UAU^{\dagger}$$
.

Let  $D^{1/2}$  be the matrix obtained by taking the square root of the diagonal entries of D. Then define  $T=D^{1/2}U$ . We then claim that  $G_U=\{V_1,\dots V_n\}$  is a unitary matrix group equivalent to G, where:

$$V_i = TA_iT^{-1}$$
.

We have three points to verify; (i) That the  $V_i$ s are unitary, (ii) That  $\varphi: G \to G_u$  defined by  $\varphi(A_i) = TA_iT^{-1} = V_i$  is an isomorphism, and (iii) that the characters of  $A_i$  and  $V_i$  are equivalent.

(i) For any  $V_i$ , we have:

$$\begin{split} V_i^\dagger V_i &= (TA_i T^{-1})^\dagger (TA_i T^{-1}) \\ &= (D^{1/2} U A_i U^\dagger D^{-1/2})^\dagger (D^{1/2} U A_i U^\dagger D^{-1/2}) \\ &= (D^{-1/2} U A_i^\dagger U^\dagger D^{1/2}) (D^{1/2} U A_i U^\dagger D^{-1/2}) \\ &= (D^{-1/2} U A_i^\dagger U^\dagger) D (U A_i U^\dagger D^{-1/2}) \\ &= (D^{-1/2} U A_i^\dagger U^\dagger) (U A U^\dagger) (U A_i U^\dagger D^{-1/2}) \\ &= (D^{-1/2} U A_i^\dagger) A (A_i^\dagger U^\dagger D^{-1/2}) \\ &= (D^{-1/2} U A_i^\dagger) \left(\sum_{j=1}^n A_j^\dagger A_j\right) (A_i U^\dagger D^{-1/2}) \\ &= D^{-1/2} U \left(\sum_{j=1}^n (A_j A_i)^\dagger (A_j A_i)\right) U^\dagger D^{-1/2} \\ &= D^{-1/2} U \left(\sum_{k=1}^n A_k^\dagger A_k\right) U^\dagger D^{-1/2} \\ &= D^{-1/2} U A U^\dagger D^{-1/2} \\ &= D^{-1/2} D D^{-1/2} \\ &= I \end{split}$$

Where in the sixth equality we use the unitarity of U, and in the ninth equality we use that  $A_jA_i=A_k$  iterates over all the group elements. To see that this is the case, it suffices to show that the map  $\psi_i:M_n\to M_n$  defined by  $\psi_i(A_j)=A_jA_i$  is a bijection. To see that it is injective, suppose that  $\psi_i(A_{j_1})=\psi_i(A_{j_2})$ . Then it follows that  $A_{j_1}A_i=A_{j_2}A_i$ , and multiplying on the left by  $A_i^{-1}$  (which exists) we find that  $A_{j_1}=A_{j_2}$ . To see that it is surjective, suppose that  $A_{j'}\in M_n$ . Then, there exists  $A_{j'}A_i^{-1}\in M_n$  such that  $\psi_i(A_{j'}A_i^{-1})=A_{j'}A_i^{-1}A_i=A_{j'}$ . We conclude that  $\psi_i$  is bijective.

(ii) Firstly,  $\varphi$  is a homomorphism as for any  $A_i, A_j$  we have:

$$\varphi(A_i)\varphi(A_j) = V_iV_j = TA_iT^{-1}TA_jT^{-1} = TA_iA_jT^{-1} = \varphi(A_iA_j).$$

Next,  $\varphi$  is surjective by construction. Finally, it is injective; suppose that  $V_i=V_j$ . Then we have that:

$$TA_iT^{-1} = TA_jT^{-1}$$

And multiplying both sides on the left by  $T^{-1}$  on the left and T on the right we find that  $A_i=A_j$ . Hence we conclude that  $\varphi$  is a bijective homomorphism and hence an isomorphism.

(iii) This is immediate from the cyclicity of the trace:

$$\chi(V_i) = \operatorname{tr}\left(TA_iT^{-1}\right) = \operatorname{tr}\left(T^{-1}TA_i\right) = \operatorname{tr}(A_i) = \chi(A_i).$$

The claim is therefore proven.

#### Exercise A2.13

Show that every irreducible Abelian matrix group is one dimensional.

#### Exercise A2.14

Show that if  $\rho$  is an irreducible representation of G, then  $|G|/d_{\rho}$  is an integer.

### Exercise A2.15

Using the Fundamental Theorem, prove that characters are orthogonal, that is:

$$\sum_{i=1}^r r_i(\chi_i^p)^*\chi_i^q = |G|\delta_{pq} \text{ and } \sum_{p=1}^r (\chi_i^p)^*\chi_j^q = \frac{|G|}{r_i}\delta_{ij}$$

where p,q, and  $\delta_{pq}$  have the same meaning as in the theorem and  $\chi^p_i$  is the value the character of the pth irreducible representation takes on the ith conjugacy class of G and  $r_i$  is the size of the ith conjugacy class of G and  $r_i$  is the size of the ith conjugacy class.

#### Exercise A2.16

 $S_3$  is the group of permutations of three elements. Suppose we order these as mapping 123 to: 123;231;312;213;132, and 321, respectively. Show that there exist two one-dimensional irreducible representations of  $S_3$ , one of which is trivial, and the other of which is 1,1,1,-1,-1, corresponding in order to the six permutations given earlier. Also show that there exists a two dimensional irreducible representation, with the matrices

$$\begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix}, \quad \frac{1}{2} \begin{bmatrix} -1 & -\sqrt{3} \\ \sqrt{3} & -1 \end{bmatrix}, \quad \frac{1}{2} \begin{bmatrix} -1 & \sqrt{3} \\ -\sqrt{3} & 1 \end{bmatrix},$$
$$\begin{bmatrix} -1 & 0 \\ 0 & 1 \end{bmatrix}, \quad \frac{1}{2} \begin{bmatrix} 1 & \sqrt{3} \\ \sqrt{3} & -1 \end{bmatrix}, \quad \frac{1}{2} \begin{bmatrix} 1 & -\sqrt{3} \\ -\sqrt{3} & 1 \end{bmatrix}$$

Verify that the representations are orthogonal.

### Exercise A2.17

Prove that the regular representation is faithful.

### Exercise A2.18

Show that the character of the regular representation is zero except on the representation of the identity element, for which  $\chi(I) = |G|$ .

#### Exercise A2.19

Use Theorem A2.5 to show that the regular representation contains  $d_{\rho^p}$  instances of each irreducible representation  $\rho^p$ . Thus, if R denotes the regular representation, and  $\hat{G}$  denotes the set of all inequivalent irreducible representations, then:

$$\chi_i^R = \sum_{\rho \in G} d_\rho \chi_i^\rho$$

### Exercise A2.20

The character of the regular representation is zero except for the conjugacy class i containing e, the identity element in G. Show, therefore, that

$$\sum_{\rho \in G} d_{\rho} \chi^{\rho}(g) = N \delta_{ge}.$$

# Exercise A2.21

Show that  $\sum_{\rho \in \hat{G}} d_{\rho}^2 = |G|.$ 

### Exercise A2.22

Substitute (A2.10) into (A2.9) and prove that  $\hat{f}(\rho)$  is obtained.

### Exercise A2.23

Let us represent an Abelian group G by  $g\in[0,N-1]$ , with addition as the group operation, and define  $\rho_h(g)=\exp\left[-2\pi igh/N\right]$  at the h representation of g. This representation is one-dimensional, so  $d_\rho=1$ . Show that the Fourier transform relations for G are

$$\hat{f}(h) = \frac{1}{\sqrt{N}} \sum_{g=0}^{N-1} f(g) e^{-2\pi i g h/N} \text{ and } f(h) = \frac{1}{\sqrt{N}} \sum_{g=0}^{N-1} \hat{f}(g) e^{2\pi i g h/N}$$

# Exercise A2.24

Using the results of Exercise A2.16, construct the Fourier transform over  $S_3$  and express it as a 6x6 unitary matrix.