

Select Solutions for “Quantum Computation and Quantum Information:
10th Anniversary Edition" by Nielsen and Chuang

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Repository

As of November, 2021, the original source \LaTeX code, located at <https://github.com/goropikari/SolutionForQuantumComputationAndQuantumInformation> has not been updated since April 2020. The extended source \LaTeX code is located at <https://github.com/tlesaul2/SolutionQCQINielsenChuang> . It may be updated more actively.

For readers

This is an unofficial solution manual for "Quantum Computation and Quantum Information: 10th Anniversary Edition" (ISBN-13: 978-1107002173) by Michael A. Nielsen and Isaac L. Chuang.

From the original author:

I have studied quantum information theory as a hobby. And I'm not a researcher. So there is no guarantee that these solutions are correct. Especially because I'm not good at mathematics, proofs are often wrong. Don't trust me. Verify yourself!

If you find some mistake or have some comments, please feel free to open an issue or a PR.

goropikari

From the second author:

I'm a mathematician relatively new to quantum information theory as of the adoption of this repo, so hope to supplement the original author's work by checking and formalizing the mathematics, overly at times, while I use the task to learn the field. The original author's sentiments about self-verification are echoed.

tlesaul2

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The exercise in this chapter is only interesting for it's mathematics, so it was moved to the end to avoid dissuading non-mathematicians from continuing to chapters more interesting for their quantum information theory.

Errata list

- p.101. eq (2.150) $\rho = \sum_m p(m)\rho_m$ should be $\rho' = \sum_m p(m)\rho_m$.
- p.408. eq (9.49) $\sum_i p_i D(\rho_i, \sigma_i) + D(p_i, q_i)$ should be $\sum_i p_i D(\rho_i, \sigma_i) + 2D(p_i, q_i)$.

$$\begin{aligned}
 \text{eqn (9.48)} &= \sum_i p_i \text{tr}(P(\rho_i - \sigma_i)) + \sum_i (p_i - q_i) \text{tr}(P\sigma_i) \\
 &\leq \sum_i p_i \text{tr}(P(\rho_i - \sigma_i)) + \sum_i |p_i - q_i| \text{tr}(P\sigma_i) \quad (\because p_i - q_i \leq |p_i - q_i|) \\
 &\leq \sum_i p_i \text{tr}(P(\rho_i - \sigma_i)) + \sum_i |p_i - q_i| \quad (\because \text{tr}(P\sigma_i) \leq 1) \\
 &= \sum_i p_i \text{tr}(P(\rho_i - \sigma_i)) + 2 \frac{\sum_i |p_i - q_i|}{2} \\
 &= \sum_i p_i \text{tr}(P(\rho_i - \sigma_i)) + 2D(p_i, q_i)
 \end{aligned}$$

- p.409. Exercise 9.12. If $\rho = \sigma$, then $D(\rho, \sigma) = 0$. Furthermore trace distance is non-negative. Therefore $0 \leq D(\mathcal{E}(\rho), \mathcal{E}(\sigma)) \leq 0 \Rightarrow D(\mathcal{E}(\rho), \mathcal{E}(\sigma)) = 0$. So I think the map \mathcal{E} is not strictly contractive. If $p \neq 1$ and $\rho \neq \sigma$, then $D(\mathcal{E}(\rho), \mathcal{E}(\sigma)) < D(\rho, \sigma)$ is satisfied.
- p.411. Exercise 9.16. eqn(9.73) $\text{tr}(A^\dagger B) = \langle m|A \otimes B|m\rangle$ should be $\text{tr}(A^{\textcolor{red}{T}} B) = \langle m|A \otimes B|m\rangle$.

Simple counter example is the case that $A = \begin{bmatrix} i & 0 \\ 0 & 0 \end{bmatrix}$. $B = \begin{bmatrix} 1 & 0 \\ 0 & 0 \end{bmatrix}$, In this case,

$$\begin{aligned}
 A^\dagger B &= \begin{bmatrix} -i & 0 \\ 0 & 0 \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 0 & 0 \end{bmatrix} = \begin{bmatrix} -i & 0 \\ 0 & 0 \end{bmatrix}, \\
 \text{tr}(A^\dagger B) &= -i, \\
 A \otimes B &= \begin{bmatrix} i & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 \end{bmatrix} \\
 \langle m|A \otimes B|m\rangle &= (\langle 00| + \langle 11|)(A \otimes B)(|00\rangle + |11\rangle) = i.
 \end{aligned}$$

Thus $\text{tr}(A^\dagger B) \neq \langle m|A \otimes B|m\rangle$.

By using following relation, we can prove.

$$\begin{aligned}
 (I \otimes A) |m\rangle &= (A^T \otimes I) |m\rangle \\
 \text{tr}(A) &= \langle m|I \otimes A|m\rangle
 \end{aligned}$$

$$\begin{aligned}
\mathrm{tr}(A^T B) &= \mathrm{tr}(B A^T) = \langle m | I \otimes B A^T | m \rangle \\
&= \langle m | (I \otimes B)(I \otimes A^T) | m \rangle \\
&= \langle m | (I \otimes B)(A \otimes I) | m \rangle \\
&= \langle m | A \otimes B | m \rangle .
\end{aligned}$$

- p.515. eqn (11.67) $S(\rho' || \rho)$ should be $S(\rho || \rho')$.

Chapter 2

Introduction to quantum mechanics

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

Soln: It is enough to express $(0, 0)$ as a linear combination of the specified vectors.

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

2.2) 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 .

Soln: With specified operations, it is enough to solve for the entries of a 2x2 matrix which converts the input vectors expressed as linear combinations of one basis, say $(|a_1\rangle, |a_2\rangle)$, into vectors expressed as linear combinations of another basis, say $(|b_1\rangle, |b_2\rangle)$.

$$A = \begin{matrix} & \begin{matrix} |b_1\rangle & |b_2\rangle \end{matrix} \\ \begin{matrix} |a_1\rangle \\ |a_2\rangle \end{matrix} & \begin{bmatrix} A_{11} & A_{12} \\ A_{21} & A_{22} \end{bmatrix} \end{matrix}$$

With $(|a_1\rangle, |a_2\rangle) = (|0\rangle, |1\rangle)$ and $(|b_1\rangle, |b_2\rangle) = (|0\rangle, |1\rangle)$, we have

$$A|0\rangle := |1\rangle = 0|0\rangle + 1|1\rangle = A_{11}|b_1\rangle + A_{21}|b_2\rangle = A_{11}|0\rangle + A_{21}|1\rangle \Rightarrow A_{11} = 0, A_{21} = 1$$

$$A|1\rangle := |0\rangle = 1|0\rangle + 0|1\rangle = A_{12}|b_1\rangle + A_{22}|b_2\rangle = A_{12}|0\rangle + A_{22}|1\rangle \Rightarrow A_{12} = 1, A_{22} = 0$$

$$\therefore A = \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix}$$

If the output basis was $(|b_1\rangle, |b_2\rangle) = (|1\rangle, |0\rangle)$ instead, then $A = I$. More formally:

$$A|0\rangle := |1\rangle = 1|1\rangle + 0|0\rangle = A_{11}|b_1\rangle + A_{21}|b_2\rangle = A_{11}|1\rangle + A_{21}|0\rangle \Rightarrow A_{11} = 1, A_{21} = 0$$

$$A|1\rangle := |0\rangle = 0|1\rangle + 1|0\rangle = A_{12}|b_1\rangle + A_{22}|b_2\rangle = A_{12}|1\rangle + A_{22}|0\rangle \Rightarrow A_{12} = 0, A_{22} = 1$$

$$\therefore A = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix}$$

With a more interesting orthonormal output basis $(|b_1\rangle, |b_2\rangle) = (|+\rangle, |-\rangle)$:

$$\begin{aligned} A|0\rangle &:= |1\rangle = \frac{\sqrt{2}}{2}|+\rangle - \frac{\sqrt{2}}{2}|-\rangle = A_{11}|b_1\rangle + A_{21}|b_2\rangle = A_{11}|+\rangle + A_{21}|-\rangle \Rightarrow A_{11} = \frac{\sqrt{2}}{2}, A_{21} = -\frac{\sqrt{2}}{2} \\ A|1\rangle &:= |0\rangle = \frac{\sqrt{2}}{2}|+\rangle + \frac{\sqrt{2}}{2}|-\rangle = A_{12}|b_1\rangle + A_{22}|b_2\rangle = A_{12}|+\rangle + A_{22}|-\rangle \Rightarrow A_{12} = \frac{\sqrt{2}}{2}, A_{22} = \frac{\sqrt{2}}{2} \\ \therefore A &= \frac{\sqrt{2}}{2} \begin{bmatrix} 1 & -1 \\ 1 & 1 \end{bmatrix} \end{aligned}$$

Note: This is similar, but not equal to **H**. Had A been the identity transformation when expressed with the same input and output bases, then the result would have been exactly **H**.

2.3) 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 $|v_i\rangle, |w_j\rangle$, and $|x_k\rangle$ 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.

Soln: Fix i . We'll show that $(B \circ A)_{ki} = (B \cdot A)_{ki}$.

$$\begin{aligned} (B \circ A)|v_i\rangle &= \sum_k (B \circ A)_{ki} |x_k\rangle = B \left(\sum_j A_{ji} |w_j\rangle \right) && \text{(Eqn 2.12, composition)} \\ &= \sum_j A_{ji} B|w_j\rangle && \text{(linearity)} \\ &= \sum_{j,k} A_{ji} B_{kj} |x_k\rangle && \text{(Eqn 2.12)} \\ &= \sum_k \left(\sum_j B_{kj} A_{ji} \right) |x_k\rangle && \text{(finiteness, commutativity)} \\ &= \sum_k ((B \cdot A)_{ki}) |x_k\rangle && \text{(definition)} \end{aligned}$$

$$\therefore (B \circ A)_{ki} = (B \cdot A)_{ki}$$

2.4) 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 representation is taken with respect to the same input and output bases. This matrix is known as the *identity matrix*

Soln: Let I be the matrix in question.

$$\begin{aligned} I|v_j\rangle &:= |v_j\rangle = \sum_i I_{ij} |v_i\rangle, \forall j. \\ \Rightarrow I_{ij} &= \delta_{ij} := \begin{cases} 1 & i = j \\ 0 & o/w \end{cases} \end{aligned}$$

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

Soln: Defined inner product on \mathbb{C}^n is

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

Equation (2.13.1), linearity in second argument:

$$\begin{aligned}
 \left((y_1, \dots, y_n), \sum_i \lambda_i (z_{i1}, \dots, z_{in}) \right) &= \left((y_1, \dots, y_n), \left(\sum_i \lambda_i z_{i1}, \dots, \sum_i \lambda_i z_{in} \right) \right) && \text{(definition)} \\
 &= \sum_j y_j^* \left(\sum_i \lambda_i z_{ij} \right) && \text{(definition)} \\
 &= \sum_j \left(\sum_i y_j^* \lambda_i z_{ij} \right) && \text{(linearity of multiplication)} \\
 &= \sum_j \left(\sum_i \lambda_i y_j^* z_{ij} \right) && \text{(associativity/commutativity)} \\
 &= \sum_i \left(\sum_j \lambda_i y_j^* z_{ij} \right) && \text{(finiteness)} \\
 &= \sum_i \lambda_i \left(\sum_j y_j^* z_{ij} \right) && \text{(linearity)} \\
 &= \sum_i \lambda_i ((y_1, \dots, y_n), (z_{i1}, \dots, z_{in})) && \text{(definition)}
 \end{aligned}$$

Equation (2.13.2), conjugate symmetry:

$$\begin{aligned}
 ((y_1, \dots, y_n), (z_1, \dots, z_n))^* &= \left(\sum_i y_i^* z_i \right)^* && \text{(definition)} \\
 &= \left(\sum_i y_i z_i^* \right) && \text{(conjugate symmetry in } \mathbb{C}^1) \\
 &= \left(\sum_i z_i^* y_i \right) && \text{(commutativity in } \mathbb{C}^1) \\
 &= ((z_1, \dots, z_n), (y_1, \dots, y_n)) && \text{(definition)}
 \end{aligned}$$

Equation (2.13.3), positive definiteness:

$$\begin{aligned}
 ((y_1, \dots, y_n), (y_1, \dots, y_n)) &= \sum_i y_i^* y_i && \text{(definition)} \\
 &= \sum_i |y_i|^2 && \text{(definition)} \\
 &\geq 0 && \text{(positive definiteness of } |\cdot|^2 \text{ over } \mathbb{C}^1)
 \end{aligned}$$

Now:

$$\begin{aligned}
 ((y_1, \dots, y_n), (y_1, \dots, y_n)) &= \sum_i |y_i|^2 \stackrel{?}{=} 0 && \text{(hypothesis)} \\
 \iff |y_i|^2 &= 0 \quad \forall i && \text{(positivity of } |\cdot|^2) \\
 \iff y_i &= 0 \quad \forall i && \text{(positive definiteness of } |\cdot|^2 \text{ over } \mathbb{C}^1) \\
 \iff (y_1, \dots, y_n) &= \mathbf{0} && \text{(definition)}
 \end{aligned}$$

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

$$\left(\sum_i \lambda_i |w_i\rangle, |v\rangle \right) = \sum_i \lambda_i^* (|w_i\rangle, |v\rangle).$$

Soln:

$$\begin{aligned}
 \left(\sum_i \lambda_i |w_i\rangle, |v\rangle \right) &= (|v\rangle, \sum_i \lambda_i |w_i\rangle)^* && \text{(conjugate symmetry)} \\
 &= \left(\sum_i \lambda_i (|v\rangle, |w_i\rangle) \right)^* && \text{(linearity in the 2nd arg.)} \\
 &= \sum_i \lambda_i^* (|v\rangle, |w_i\rangle)^* && \text{(distributivity of complex conjugate)} \\
 &= \sum_i \lambda_i^* (|w_i\rangle, |v\rangle) && \text{(conjugate symmetry)}
 \end{aligned}$$

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

Soln:

$$\begin{aligned}
 (|w\rangle, |v\rangle) &= \langle w|v\rangle && \text{(notation)} \\
 &= \begin{bmatrix} 1^* & 1^* \end{bmatrix} \begin{bmatrix} 1 \\ -1 \end{bmatrix} && \text{(definition)} \\
 &= 1^* \cdot 1 + 1^* \cdot (-1) && \text{(matrix multiplication)} \\
 &= 1 \cdot 1 - 1 \cdot 1 = 0 && \text{(arithmetic)} \\
 \frac{|w\rangle}{\| |w\rangle \|} &= \frac{|w\rangle}{\sqrt{\langle w|w\rangle}} = \frac{1}{\sqrt{2}} \begin{bmatrix} 1 \\ 1 \end{bmatrix} = |+\rangle \\
 \frac{|v\rangle}{\| |v\rangle \|} &= \frac{|v\rangle}{\sqrt{\langle v|v\rangle}} = \frac{1}{\sqrt{2}} \begin{bmatrix} 1 \\ -1 \end{bmatrix} = |-\rangle
 \end{aligned}$$

2.8) Prove that the Gram-Schmidt procedure produces an orthonormal basis.

Soln: We prove inductively. For $d = 1$, the only requirement is that the procedure normalize $|w_d\rangle$, which it does by definition for all d . For $d = 2$, suppose $|v_1\rangle, \dots, |v_{d-1}\rangle$ is a orthonormal basis for the subspace spanned by $|w_1\rangle, \dots, |w_{d-1}\rangle$. Being a basis, the subspace spanned by $|v_1\rangle, \dots, |v_{d-1}\rangle$ is the same. Linear independence of $|w_1\rangle, \dots, |w_d\rangle$ implies that $|w_d\rangle$ is not in this subspace, so $|v_1\rangle, \dots, |v_{d-1}\rangle, |w_d\rangle$ is easily seen to be linearly independent as well. It remains to be shown that $|v_d\rangle$ is linearly independent of $|v_1\rangle, \dots, |v_{d-1}\rangle$, and is orthogonal to all such vectors. For independence, note that any dependence relation between $|v_1\rangle, \dots, |v_d\rangle$ immediately induces one between $|v_1\rangle, \dots, |v_{d-1}\rangle, |w_d\rangle$, violating their independence. For orthogonality, let $1 \leq j \leq d-1$. We show $\langle v_j | v_d \rangle = 0$, completing the proof.

$$\begin{aligned}
 \langle v_j | v_d \rangle &= \langle v_j | \left(\frac{|w_d\rangle - \sum_{i=1}^{d-1} \langle v_i | w_d \rangle |v_i\rangle}{\left\| |w_d\rangle - \sum_{i=1}^{d-1} \langle v_i | w_d \rangle |v_i\rangle \right\|} \right) && \text{(definition)} \\
 &= \frac{\langle v_j | w_d \rangle - \sum_{i=1}^{d-1} \langle v_i | w_d \rangle \langle v_j | v_i \rangle}{\left\| |w_d\rangle - \sum_{i=1}^{d-1} \langle v_i | w_d \rangle |v_i\rangle \right\|} && \text{(linearity in the 2nd argument)} \\
 &= \frac{\langle v_j | w_d \rangle - \sum_{i=1}^{d-1} \langle v_i | w_d \rangle \delta_{ij}}{\left\| |w_d\rangle - \sum_{i=1}^{d-1} \langle v_i | w_d \rangle |v_i\rangle \right\|} && \text{(orthonormality of } |v_1\rangle, \dots, |v_{d-1}\rangle) \\
 &= \frac{\langle v_j | w_d \rangle - \langle v_j | w_d \rangle}{\left\| |w_d\rangle - \sum_{i=1}^{d-1} \langle v_i | w_d \rangle |v_i\rangle \right\|} && \text{(definition of } \delta_{ij}) \\
 &= 0. && \text{(arithmetic)}
 \end{aligned}$$

2.9) (Pauli operators and the outer product) The Pauli matrices 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.

$$\begin{aligned}\sigma_0 &= I = |0\rangle\langle 0| + |1\rangle\langle 1| \\ \sigma_x &= \sigma_1 = X = |1\rangle\langle 0| + |0\rangle\langle 1| \\ \sigma_y &= \sigma_2 = Y = i|1\rangle\langle 0| - i|0\rangle\langle 1| \\ \sigma_z &= \sigma_3 = Z = |0\rangle\langle 0| - |1\rangle\langle 1|\end{aligned}$$

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

Soln:

$$\begin{aligned}|v_j\rangle\langle v_k| &= I_V |v_j\rangle\langle v_k| I_V && \text{(multiply by identity)} \\ &= \left(\sum_p |v_p\rangle\langle v_p|\right) |v_j\rangle\langle v_k| \left(\sum_q |v_q\rangle\langle v_q|\right) && \text{(completeness)} \\ &= \sum_{p,q} |v_p\rangle\langle v_p|v_j\rangle\langle v_k|v_q\rangle\langle v_q| && \text{(linearity and outer product definition)} \\ &= \sum_{p,q} \delta_{pj}\delta_{kq} |v_p\rangle\langle v_q| && \text{(orthonormality)}\end{aligned}$$

Thus

$$(|v_j\rangle\langle v_k|)_{pq} = \delta_{pj}\delta_{kq} = \begin{cases} 1 & p=j, k=q \\ 0 & \text{o/w} \end{cases}.$$

That is, $|v_j\rangle\langle v_k|$ is a square matrix with a 1 in row j , column k , and 0s everywhere else.

(Cauchy-Schwartz inequality) A brief expansion from a mathematician: in equation (2.26), other $|i\rangle$ -basis vectors appear, but since $\langle i|v\rangle = \langle v|i\rangle^*$, $a \cdot a^* = \|a\|^2 \geq 0$ for all $a \in \mathbb{C}$, and $\langle \cdot | \cdot \rangle$ is positive definite, all terms but the first constructed in terms of $|w\rangle$ are non-negative and can be removed, leaving the inequality.

2.11) Eigendecomposition of the Pauli matrices: Find the eigenvectors, eigenvalues, and diagonal representations of the Pauli matrices X, Y , and Z .

Soln:

$$X = \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix}, \det(X - \lambda I) = \det\left(\begin{bmatrix} -\lambda & 1 \\ 1 & -\lambda \end{bmatrix}\right) = \lambda^2 - 1 = 0 \Rightarrow \lambda = \pm 1$$

If $\lambda = 1$,

$$\begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} \begin{bmatrix} c_1 \\ c_2 \end{bmatrix} = \begin{bmatrix} c_2 \\ c_1 \end{bmatrix} = 1 \begin{bmatrix} c_1 \\ c_2 \end{bmatrix} \Rightarrow c_2 = c_1$$

The eigenspace corresponding to $\lambda = 1$ is the set of vectors $\begin{bmatrix} c \\ c \end{bmatrix}$. The vector $\frac{1}{\sqrt{2}} \begin{bmatrix} 1 \\ 1 \end{bmatrix} = |+\rangle$ is such a unit (normalized) vector. If $\lambda = -1$,

$$\begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} \begin{bmatrix} c_1 \\ c_2 \end{bmatrix} = \begin{bmatrix} c_2 \\ c_1 \end{bmatrix} = -1 \begin{bmatrix} c_1 \\ c_2 \end{bmatrix} \Rightarrow c_2 = -c_1$$

The eigenspace corresponding to $\lambda = -1$ is the set of vectors $\begin{bmatrix} c \\ -c \end{bmatrix}$. The vector $\frac{1}{\sqrt{2}} \begin{bmatrix} 1 \\ -1 \end{bmatrix} = |- \rangle$ is such a unit (normalized) vector. So, a diagonal representation of X (when expressed in terms of the computational basis) is $(|+\rangle \langle +|) - (|- \rangle \langle -|) = \begin{bmatrix} \frac{1}{2} & \frac{1}{2} \\ \frac{1}{2} & \frac{1}{2} \end{bmatrix} - \begin{bmatrix} \frac{1}{2} & -\frac{1}{2} \\ -\frac{1}{2} & \frac{1}{2} \end{bmatrix} (= X)$.

$$Y = \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix}, \det(Y - \lambda I) = \det \left(\begin{bmatrix} -\lambda & -i \\ i & -\lambda \end{bmatrix} \right) = \lambda^2 - (i)(-i) = \lambda^2 - 1 = 0 \Rightarrow \lambda = \pm 1$$

If $\lambda = 1$,

$$\begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} \begin{bmatrix} c_1 \\ c_2 \end{bmatrix} = \begin{bmatrix} -i \cdot c_2 \\ i \cdot c_1 \end{bmatrix} = 1 \begin{bmatrix} c_1 \\ c_2 \end{bmatrix} \Rightarrow c_2 = i \cdot c_1$$

The eigenspace corresponding to $\lambda = 1$ is the set of vectors $\begin{bmatrix} c \\ i \cdot c \end{bmatrix}$. The vector $\frac{1}{\sqrt{2}} \begin{bmatrix} 1 \\ i \end{bmatrix} \equiv |\psi_{y+} \rangle$ is such a unit (normalized) vector. If $\lambda = -1$,

$$\begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} \begin{bmatrix} c_1 \\ c_2 \end{bmatrix} = \begin{bmatrix} -i \cdot c_2 \\ i \cdot c_1 \end{bmatrix} = -1 \begin{bmatrix} c_1 \\ c_2 \end{bmatrix} \Rightarrow c_2 = -i \cdot c_1$$

The eigenspace corresponding to $\lambda = -1$ is the set of vectors $\begin{bmatrix} c \\ -i \cdot c \end{bmatrix}$. The vector $\frac{1}{\sqrt{2}} \begin{bmatrix} 1 \\ -i \end{bmatrix} \equiv |\psi_{y-} \rangle$ is such a unit (normalized) vector. So, a diagonal representation of Y (when expressed in terms of the computational basis) is $(|\psi_{y+} \rangle \langle \psi_{y+}|) - (|\psi_{y-} \rangle \langle \psi_{y-}|) = \begin{bmatrix} \frac{1}{2} & -\frac{i}{2} \\ \frac{i}{2} & \frac{1}{2} \end{bmatrix} - \begin{bmatrix} \frac{1}{2} & \frac{i}{2} \\ -\frac{i}{2} & \frac{1}{2} \end{bmatrix} (= Y)$.

$$Z = \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix}, \det(Z - \lambda I) = \det \left(\begin{bmatrix} 1 - \lambda & 0 \\ 0 & -1 - \lambda \end{bmatrix} \right) = (\lambda + 1)(\lambda - 1) = 0 \Rightarrow \lambda = \pm 1$$

If $\lambda = 1$,

$$\begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} \begin{bmatrix} c_1 \\ c_2 \end{bmatrix} = \begin{bmatrix} c_1 \\ -c_2 \end{bmatrix} = 1 \begin{bmatrix} c_1 \\ c_2 \end{bmatrix} \Rightarrow c_2 = -c_2 \Rightarrow c_2 = 0$$

The eigenspace corresponding to $\lambda = 1$ is the set of vectors $\begin{bmatrix} c \\ 0 \end{bmatrix}$. The vector $\begin{bmatrix} 1 \\ 0 \end{bmatrix} = |0 \rangle$ is such a unit (normalized) vector. If $\lambda = -1$,

$$\begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} \begin{bmatrix} c_1 \\ c_2 \end{bmatrix} = \begin{bmatrix} c_1 \\ -c_2 \end{bmatrix} = -1 \begin{bmatrix} c_1 \\ c_2 \end{bmatrix} \Rightarrow c_1 = -c_1$$

The eigenspace corresponding to $\lambda = -1$ is the set of vectors $\begin{bmatrix} 0 \\ c \end{bmatrix}$. The vector $\begin{bmatrix} 0 \\ 1 \end{bmatrix} = |1 \rangle$ is such a unit (normalized) vector. So, the computation basis *is* the eigenbasis for Z , and a diagonal representation of Z is $(|0 \rangle \langle 0|) - (|1 \rangle \langle 1|) = \begin{bmatrix} 1 & 0 \\ 0 & 0 \end{bmatrix} - \begin{bmatrix} 0 & 0 \\ 0 & 1 \end{bmatrix} (= Z)$.

2.12) Prove that the matrix $\begin{bmatrix} 1 & 0 \\ 1 & 1 \end{bmatrix} (\equiv A)$ is not diagonalizable.

Soln:

$$\det(A - \lambda I) = \det \left(\begin{bmatrix} 1 & 0 \\ 1 & 1 \end{bmatrix} - \lambda I \right) = \det \left(\begin{bmatrix} 1 - \lambda & 0 \\ 1 & 1 - \lambda \end{bmatrix} \right) = (1 - \lambda)^2 = 0 \Rightarrow \lambda = 1 \text{ (with multiplicity 2)}$$

All eigenvectors $\begin{bmatrix} c_1 \\ c_2 \end{bmatrix}$ satisfy:

$$\begin{bmatrix} 1 & 0 \\ 1 & 1 \end{bmatrix} \begin{bmatrix} c_1 \\ c_2 \end{bmatrix} = \begin{bmatrix} c_1 \\ c_1 + c_2 \end{bmatrix} = 1 \begin{bmatrix} c_1 \\ c_2 \end{bmatrix} \Rightarrow c_1 = 0$$

So, the eigenspace corresponding to eigenvalue 1 of A is 1-dimensional, with a single unit (normalized) vector $\begin{bmatrix} 0 \\ 1 \end{bmatrix} = |1\rangle$. The only possible diagonal representation of A would then be $A = |1\rangle\langle 1|$, but this equality does not hold. We conclude that A has no diagonal representation and is not diagonalizable.

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

Soln: We show that $|v\rangle\langle w|$ has the defining property of $(|w\rangle\langle v|)^\dagger$, *i.e.* if $|\psi\rangle, |\phi\rangle$ are arbitrary vectors in V , then $(|\psi\rangle, (|w\rangle\langle v|)|\phi\rangle) = ((|v\rangle\langle w|)|\psi\rangle, |\phi\rangle)$. We do so by expanding $(|\psi\rangle, (|w\rangle\langle v|)|\phi\rangle)^*$ in two different ways.

$$\begin{aligned} (|\psi\rangle, (|w\rangle\langle v|)|\phi\rangle)^* &= ((|w\rangle\langle v|)^\dagger|\psi\rangle, |\phi\rangle)^* && \text{(definition of } \dagger) \\ &= (|\phi\rangle, (|w\rangle\langle v|)^\dagger|\psi\rangle) && \text{(conjugate symmetry)} \end{aligned}$$

On the other hand,

$$\begin{aligned} (|\psi\rangle, (|w\rangle\langle v|)|\phi\rangle)^* &= (\langle\psi|w\rangle, \langle v|\phi\rangle)^* && \text{(associativity of } \langle\cdot|\cdot\rangle, \langle\cdot|\cdot\rangle, \text{ and } |\cdot\rangle\langle\cdot|) \\ &= (\langle\phi|v\rangle, \langle w|\psi\rangle)^* && \text{(conjugate symmetry)} \\ &= (\langle\phi|, (|v\rangle\langle w|)|\psi\rangle). && \text{(notation)} \end{aligned}$$

Thus

$$(|\phi\rangle, (|w\rangle\langle v|)^\dagger|\psi\rangle) = (\langle\phi|, (|v\rangle\langle w|)|\psi\rangle) \text{ for arbitrary vectors } |\psi\rangle, |\phi\rangle$$

We conclude that $(|w\rangle\langle v|)^\dagger$ and $|v\rangle\langle w|$ are the same operator, so $(|w\rangle\langle v|)^\dagger = |v\rangle\langle w|$.

2.14) Anti-linearity of the adjoint: Show that the adjoint operation is anti-linear,

$$\left(\sum_i a_i A_i\right)^\dagger = \sum_i a_i^* A_i^\dagger$$

Soln: It is tempting to assume that $(\sum_i a_i A_i)^\dagger = \sum_i (a_i A_i)^\dagger$, *i.e.* that the † transformation is additive, but we don't yet know this. It will follow from the fact that $A^\dagger \equiv (A^*)^T$ given after problem 2.15, and that both $*$ and T are linear. This itself is not hard to prove by observing that $(A^*)^T$ has the defining property of A^\dagger , making use of the matrix formulation of the inner product. Without the assumption though, we must be careful to carry around the full sums until additivity (and in-fact full linearity) is known.

$$\begin{aligned}
\left(\left(\sum_i a_i A_i \right)^\dagger |\phi\rangle, |\psi\rangle \right) &= \left(|\phi\rangle, \left(\sum_i a_i A_i \right) |\psi\rangle \right) && \text{(definition of } \dagger \text{)} \\
&= \left(|\phi\rangle, \sum_i a_i A_i |\psi\rangle \right) && \text{(distributivity of matrix multiplication)} \\
&= \sum_i a_i (|\phi\rangle, A_i |\psi\rangle) && \text{(linearity in the second argument)} \\
&= \sum_i a_i (A_i^\dagger |\phi\rangle, |\psi\rangle) && \text{(definition of } \dagger \text{)} \\
&= \sum_i (a_i^* A_i^\dagger |\phi\rangle, |\psi\rangle) && \text{(conjugate-linearity in the first argument)} \\
&= \left(\left(\sum_i a_i^* A_i^\dagger \right) |\phi\rangle, |\psi\rangle \right) && \text{(distributivity of matrix multiplication)} \\
\text{therefore } \left(\sum_i a_i A_i \right)^\dagger &= \sum_i a_i^* A_i^\dagger && \square
\end{aligned}$$

2.15) Show that $(A^\dagger)^\dagger = A$.

Soln: We show that A has the defining property of the adjoint of A^\dagger .

$$\begin{aligned}
\left((A^\dagger)^\dagger |\psi\rangle, |\phi\rangle \right) &= (|\psi\rangle, A^\dagger |\phi\rangle) && \text{(definition of } (A^\dagger)^\dagger \text{)} \\
&= (A^\dagger |\phi\rangle, |\psi\rangle)^* && \text{(conjugate symmetry)} \\
&= (|\phi\rangle, A |\psi\rangle)^* && \text{(definition of } A^\dagger \text{)} \\
&= (A |\psi\rangle, |\phi\rangle) && \text{(conjugate symmetry)} \\
\text{therefore } (A^\dagger)^\dagger &= A && \square
\end{aligned}$$

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

$$\begin{aligned}
P &= \sum_i |i\rangle \langle i|. && \text{(definition)} \\
P^2 &= \left(\sum_i |i\rangle \langle i| \right) \left(\sum_j |j\rangle \langle j| \right) && \text{(square definition)} \\
&= \sum_{i,j} |i\rangle \langle i|j\rangle \langle j| && \text{(distributivity)} \\
&= \sum_{i,j} |i\rangle \langle j| \delta_{ij} && \text{(evaluate } \langle i|j\rangle \text{)} \\
&= \sum_i |i\rangle \langle i| && \text{(evaluate sum over } j \text{)} \\
&= P && \text{(definition)}
\end{aligned}$$

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

Proof. (\Rightarrow) Suppose A is Hermitian. Then $A = A^\dagger$. Let λ be an eigenvalue of A with unit-eigenvector $|\lambda\rangle$.

We have:

$$\begin{aligned} A|\lambda\rangle &= \lambda|\lambda\rangle && \text{(definition)} \\ \langle\lambda|A|\lambda\rangle &= \lambda\langle\lambda|\lambda\rangle && \text{(multiply by } \langle\lambda|) \\ &= \lambda. && (\lambda \text{ is a unit-vector}) \end{aligned}$$

Now:

$$\begin{aligned} \lambda^* &= \langle\lambda|A|\lambda\rangle^* && \text{(conjugate)} \\ &= (|\lambda\rangle, A|\lambda\rangle)^* && \text{(change notation)} \\ &= (A|\lambda\rangle, |\lambda\rangle) && \text{(conjugate symmetry)} \\ &= (A^\dagger|\lambda\rangle, |\lambda\rangle) && \text{(hypothesis)} \\ &= (|\lambda\rangle, A|\lambda\rangle) && \text{(definition of } \dagger) \\ &= \lambda && \text{(from above)} \end{aligned}$$

So the eigenvalue λ is real, since only real numbers are equal to their conjugates.

(\Leftarrow) To prove the converse we make use of the spectral decomposition theorem. It's proof does *not* use the fact that a normal matrix is Hermitian if and only if it's eigenvalues are real, so using it here does not make this proof circular. Suppose the eigenvalues of A are real. From the spectral decomposition theorem there exists a set of eigenvalues λ_i and a corresponding orthonormal basis $|\lambda_i\rangle$ such that

$$A = \sum_i \lambda_i |\lambda_i\rangle\langle\lambda_i| \quad \text{(spectral decomposition)}$$

From this we have:

$$\begin{aligned} A^\dagger &= \left(\sum_i \lambda_i |\lambda_i\rangle\langle\lambda_i| \right)^\dagger && \text{(apply adjoint)} \\ &= \sum_i \lambda_i^* (|\lambda_i\rangle\langle\lambda_i|)^\dagger && \text{(anti-linearity)} \\ &= \sum_i \lambda_i |\lambda_i\rangle\langle\lambda_i| && (\lambda_i \text{ real, projectors are Hermitian)} \\ &= A && \text{(from spectral decomposition)} \end{aligned}$$

Thus A is Hermitian. □

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 θ .

Soln: Suppose λ is an eigenvalue with corresponding unit-eigenvector $|v\rangle$

$$\begin{aligned} 1 &= \langle v|v\rangle && (|v\rangle \text{ is a unit vector}) \\ &= \langle v|I|v\rangle && \text{(multiply by identity)} \\ &= \langle v|U^\dagger U|v\rangle && (U \text{ is unitary}) \\ &= (\langle v|U^\dagger)(U|v\rangle) && \text{(associativity of matrix multiplication)} \\ &= (U|v\rangle)^\dagger (U|v\rangle) && \text{(arithmetic properties of } \dagger) \\ &= (\lambda|v\rangle)^\dagger (\lambda|v\rangle) && (|v\rangle \text{ is an eigenvector)} \\ &= \lambda^* \lambda \langle v|v\rangle && \text{(re-apply } \dagger \text{ and simplify)} \\ &= \|\lambda\|^2 && \text{(definition of } \|\cdot\|, |v\rangle \text{ is a unit-vector)} \end{aligned}$$

Now $\|\lambda\| = 1$, and all complex numbers with modulus 1 are located on the unit-circle in \mathbb{C} and can be expressed as $e^{i\theta}$ for some real $\theta \in [0, 2\pi)$

2.19) Show that the Pauli matrices are Hermitian and unitary

Soln: It is easy to see that the Pauli matrices are Hermitian (self-adjoint) given the conjugate-transpose formula. We still must show that their squares are the identity:

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

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

$$Z^\dagger Z = Z^2 = \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} = \begin{bmatrix} 1 & 0 \\ 0 & (-1)^2 \end{bmatrix} = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix} = I$$

2.20) Suppose A' and A'' are matrix representations of an operator 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'' .

Soln:

$$U \equiv \sum_i |w_i\rangle \langle v_i|, \quad U^\dagger = \sum_j |v_j\rangle \langle w_j| \quad (\text{construct a unitary operator and its adjoint})$$

$$\begin{aligned} A'_{ij} &= \langle v_i | A | v_j \rangle && (\text{given}) \\ &= \langle v_i | U U^\dagger A U U^\dagger | v_j \rangle && (U \text{ is unitary; } U U^\dagger = I) \\ &= \sum_{p,q,r,s} \langle v_i | w_p \rangle \langle v_p | v_q \rangle \langle w_q | A | w_r \rangle \langle v_r | v_s \rangle \langle w_s | v_j \rangle && (\text{expand } U, U^\dagger, \text{ apply linearity}) \\ &= \sum_{p,q,r,s} \langle v_i | w_p \rangle \delta_{pq} A''_{qr} \delta_{rs} \langle w_s | v_j \rangle && (|v_i\rangle \text{ is orthonormal, apply given for } A'') \\ &= \sum_{p,r} \langle v_i | w_p \rangle \langle w_r | v_j \rangle A''_{pr} && (\text{collect non-zero terms and re-index}) \end{aligned}$$

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.

Theorem 2.1 (Spectral decomposition) A Hermitian operator M on a vector space V is diagonal with respect to some orthonormal basis for V .

Proof. We induct on the dimension of V , as in the boxed proof. Let λ be an eigenvalue of M , P be the projector onto the λ eigenspace, and Q the projector onto the orthogonal complement.

$$\begin{aligned} M &= I M I && (\text{trivial}) \\ &= (P + Q) M (P + Q) && (\text{definition of } Q) \\ &= P M P + Q M P + P M Q + Q M Q && (\text{expand}) \end{aligned}$$

Now $PMP = \lambda P$ and $QMP = 0$ as before. To show that $PMQ = 0$ is as easy as substituting M^\dagger :

$$\begin{aligned}
 PMQ &= PM^\dagger Q && (M \text{ is Hermitian}) \\
 &= P(M^{*T}Q) && (\dagger = {}^{*T}) \\
 &= (QM^*P)^T && (\text{properties of } T) \\
 &= ((QMP)^*)^T && (\text{properties of } *) \\
 &= 0 && (QMP = 0)
 \end{aligned}$$

Thus $M = PMP + QMQ$. Next, we prove QMQ is normal.

$$\begin{aligned}
 QMQ(QMQ)^\dagger &= QMQQ^\dagger M^\dagger Q^\dagger && (\text{properties of } \dagger, \text{ and symmetry}) \\
 &= QMQQM^\dagger Q && (\text{projectors are Hermitian}) \\
 &= QM^\dagger QMQ && (M = M^\dagger) \\
 &= Q^\dagger M^\dagger Q^\dagger QMQ && (\text{projectors are Hermitian}) \\
 &= (QMQ)^\dagger QMQ && (\text{properties of } \dagger, \text{ and symmetry})
 \end{aligned}$$

Therefore QMQ is normal. By induction, QMQ is diagonal. The rest follows Box 2.2 identically. \square

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

Soln: Suppose A is a Hermitian operator and $|v_1\rangle, |v_2\rangle$ are eigenvectors of A with eigenvalues λ_1, λ_2 , with $\lambda_1 \neq \lambda_2$. Then

$$\langle v_1|A|v_2\rangle = \lambda_2 \langle v_1|v_2\rangle. \quad (\text{definition of } v_1, \text{ linearity of } \langle \cdot | \cdot \rangle)$$

On the other hand,

$$\begin{aligned}
 \langle v_1|A|v_2\rangle &= \langle v_1|A^\dagger|v_2\rangle && (A \text{ is Hermitian}) \\
 &= \langle v_2|A|v_1\rangle^* && (\text{properties of } \dagger, \text{ Hermitian} \Rightarrow \text{self-transpose}) \\
 &= \lambda_1 \langle v_2|v_1\rangle^* && (\text{definition of } v_1, \text{ linearity of } \langle \cdot | \cdot \rangle) \\
 &= \lambda_1 \langle v_1|v_2\rangle && (\text{properties of } *)
 \end{aligned}$$

Thus

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

Since $\lambda_1 - \lambda_2 \neq 0$, we must have $\langle v_1|v_2\rangle = 0$, so v_1 and v_2 are orthogonal.

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

Soln: Suppose P is projector and $|v\rangle$ is an eigenvector of P with eigenvalue λ . By exercise 2.16, $P^2 = P$. We have $P|v\rangle = \lambda|v\rangle$ by hypothesis. Alternatively,

$$\begin{aligned}
 P|v\rangle &= P^2|v\rangle && (\text{exercise 2.16}) \\
 &= \lambda P|v\rangle && (\text{hypothesis, linearity}) \\
 &= \lambda^2|v\rangle && (\text{hypothesis})
 \end{aligned}$$

Therefore

$$\begin{aligned}
 \lambda &= \lambda^2 \\
 \lambda^2 - \lambda &= 0 \\
 \lambda(\lambda - 1) &= 0 \\
 \lambda &= 0 \text{ or } 1.
 \end{aligned}$$

2.24) (Hermiticity of positive operators) Show that a positive operator is necessarily Hermitian.

Soln: Let A be a positive operator, that is, suppose $\langle v|A|v \rangle$ is real and ≥ 0 for all $|v \rangle$. Define $B = \frac{A+A^\dagger}{2}$ and $C = \frac{A-A^\dagger}{2i}$. Simple complex arithmetic will show that $A = B+iC$. B is clearly Hermitian by commutativity of operator addition. C is also Hermitian by linearity of the adjoint, noting that $\left(\frac{1}{2i}\right)^* = -\frac{1}{2i}$. There are two ways to proceed: one heuristic, and one mathematically rigorous. We'll start with a heuristic outline of the proof, then provide some mathematically rigorous detail after the fact.

Let v be a vector and note that it can be proven (below) that $\langle v|B|v \rangle$ and $\langle v|C|v \rangle$ are both real numbers. Now $\langle v|A|v \rangle = \langle v|B+iC|v \rangle = \langle v|B|v \rangle + i\langle v|C|v \rangle$ by the construction of B and C , and the linearity of $\langle \cdot | \cdot \rangle$. By hypothesis, $\langle v|A|v \rangle$ is a non-negative real number, so $\langle v|C|v \rangle = 0$, since both $\langle v|B|v \rangle$ and $\langle v|C|v \rangle$ are real. This will be enough to show that $C = 0$ which yields $A = A^\dagger$ by the definition of C , that is, A is Hermitian.

Now, to complete the proof, we need to rigorously show that both $\langle v|B|v \rangle$ and $\langle v|C|v \rangle$ are real numbers, and that if $\langle v|C|v \rangle = 0$ for all $|v \rangle$, then $C = 0$. Let W be Hermitian, thus normal, and note that by exercise 2.17, W has real eigenvalues, say ω_i . By the spectral decomposition theorem there is an orthonormal basis, say $|w_i \rangle$, such that $W = \sum_i \omega_i |w_i \rangle \langle w_i|$. Let $|v \rangle$ be an arbitrary vector, expressed in the orthonormal $|w_i \rangle$ -basis as $\sum_i \alpha_i |w_i \rangle$.

$$\begin{aligned}
 \langle v|W|v \rangle &= \left\langle \sum_j \alpha_j |w_j \rangle \left| W \right| \sum_i \alpha_i |w_i \rangle \right\rangle && \text{(by construction)} \\
 &= \sum_i \sum_j \alpha_i \alpha_j^* \langle w_j | W | w_i \rangle && \text{((conjugate) linearity of } \langle \cdot | \cdot \rangle \text{)} \\
 &= \sum_i \sum_j \alpha_i \alpha_j^* \left\langle w_j \left| \sum_k \omega_k |w_k \rangle \langle w_k| \right| w_i \right\rangle && \text{(spectral decomposition)} \\
 &= \sum_i \sum_j \sum_k \alpha_i \alpha_j^* \omega_k \langle w_j | w_k \rangle \langle w_k | w_i \rangle && \text{(linearity of } \langle \cdot | \cdot \rangle \text{)} \\
 &= \sum_i \sum_j \sum_k \alpha_i \alpha_j^* \omega_k \delta_{jk} \delta_{ki} && \text{(orthonormality of the } |w_i \rangle \text{ basis)} \\
 &= \sum_k \alpha_k \alpha_k^* \omega_k && \text{(collecting non-zero terms)} \\
 &= \sum_k \|\alpha_k\|^2 \omega_k && \text{(definition of } \|\cdot\| \text{)}
 \end{aligned}$$

The ω_k are real numbers by exercise 2.17, and the $\|\alpha_k\|^2$ are real by the definition of $\|\cdot\|$, so $\langle v|W|v \rangle$ is a sum of real number, and hence also real itself. Applying this to B and C above completes the first missing part. To finally complete the proof we'll require Theorem 2.0.1 below, more generally applicable to linear operators on complex vector spaces, without the assumption of Hermiticity. The proof follows an MIT 8.05 Quantum Physics II lecture note by Prof. Barton Zwiebach (https://ocw.mit.edu/courses/physics/8-05-quantum-physics-ii-fall-2013/lecture-notes/MIT8_05F13_Chap_03.pdf)

Proposition. 2.0.1. *Let T be a linear operator on a complex vector space V . If $\langle u|T|v \rangle = 0$ for all $|u \rangle, |v \rangle \in V$, then $T = 0$.*

Proof. Let $|u \rangle = T|v \rangle$. Then $\langle T|v \rangle |T|v \rangle = \langle T|v \rangle |T|v \rangle = \|T|v \rangle\|^2 = 0$, which implies $T|v \rangle = 0$ for all v by property 3 of the inner product (page 65). T is identically 0, so is the zero operator, i.e. $T = 0$. \square

Theorem. 2.0.1. *Let T be a linear operator on a complex vector space V . If $\langle v|T|v \rangle = 0$ for all $|v \rangle \in V$, then $T = 0$.*

Proof. Note that the weakened hypothesis doesn't directly apply if $|u \rangle \neq |v \rangle$. We show that the "off-diagonal", distinct vector hypothesis of Proposition 2.0.1 can be derived from the weakened "diagonal"

hypothesis' of this theorem, that is, if $\langle v|T|v\rangle = 0$ for all $|v\rangle$, then $\langle u|T|v\rangle = 0$ for all $|u\rangle, |v\rangle$. Then apply proposition 2.0.1

Suppose $|u\rangle, |v\rangle \in V$. Then note that by “foiling” the $\langle \cdot | \cdot \rangle$'s, we can show a “polarization” identity, expressing $\langle u|T|v\rangle$ as follows

$$\begin{aligned} & \frac{1}{4} \left(\langle u+v|T|u+v\rangle - \langle u-v|T|u-v\rangle + \frac{1}{i} \langle u+iv|T|u+iv\rangle - \frac{1}{i} \langle u-iv|T|u-iv\rangle \right) = \\ & \frac{1}{4} \left((\langle u|T|u\rangle + \langle u|T|v\rangle + \langle v|T|u\rangle + \langle v|T|v\rangle) - (\langle u|T|u\rangle - \langle u|T|v\rangle - \langle v|T|u\rangle + \langle v|T|v\rangle) + \dots \right. \\ & \left. \frac{1}{i} (\langle u|T|u\rangle + i \langle u|T|v\rangle - i \langle v|T|u\rangle + \langle v|T|v\rangle) - \frac{1}{i} (\langle u|T|u\rangle - i \langle u|T|v\rangle + i \langle v|T|u\rangle + \langle v|T|v\rangle) \right) = \\ & \frac{1}{4} (0 \langle u|T|u\rangle + 4 \langle u|T|v\rangle + 0 \langle v|T|u\rangle + 0 \langle v|T|v\rangle) = \\ & \langle u|T|v\rangle \end{aligned}$$

Applying the diagonal hypothesis to $|u+v\rangle, |u-v\rangle, |u+iv\rangle$, and $|u-iv\rangle$ in the first expression above gives that $\langle u|T|v\rangle = 0$ for all $|u\rangle, |v\rangle$, hence by Proposition 2.0.1, $T = 0$. \square

Applying Theorem 2.0.1 to C from above finally completes the proof of the Hermiticity of positive operators.

2.25) Show that for any operator A , $A^\dagger A$ is positive.

Soln: Its enough to show that $\langle v|A^\dagger A|v\rangle \geq 0$ for all v , but note that $\langle v|A^\dagger A|v\rangle = \|Av\|^2$, which is non-negative, so $A^\dagger A$ is a positive operator.

2.26) Let $|\psi\rangle = (|0\rangle + |1\rangle)/\sqrt{2}$ ($= |+\rangle$). 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.

Soln:

$$\begin{aligned} |\psi\rangle^{\otimes 2} &= \frac{1}{\sqrt{2}}(|0\rangle + |1\rangle) \otimes \frac{1}{\sqrt{2}}(|0\rangle + |1\rangle) \\ &= \frac{1}{2}(|00\rangle + |01\rangle + |10\rangle + |11\rangle) \\ &= \frac{1}{2} \begin{bmatrix} 1 \\ 1 \\ 1 \\ 1 \end{bmatrix} \end{aligned}$$

$$\begin{aligned} |\psi\rangle^{\otimes 3} &= \frac{1}{\sqrt{2}}(|0\rangle + |1\rangle) \otimes \frac{1}{\sqrt{2}}(|0\rangle + |1\rangle) \otimes \frac{1}{\sqrt{2}}(|0\rangle + |1\rangle) \\ &= \frac{1}{2\sqrt{2}}(|000\rangle + |001\rangle + |010\rangle + |011\rangle + |100\rangle + |101\rangle + |110\rangle + |111\rangle) \\ &= \frac{1}{2\sqrt{2}} \begin{bmatrix} 1 \\ 1 \\ 1 \\ 1 \\ 1 \\ 1 \\ 1 \\ 1 \end{bmatrix} \end{aligned}$$

2.27) Calculate the matrix representations 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?

Soln:

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

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

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

In general, the tensor product is not commutative.

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

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

Soln: Let A be $n_1 \times m_1$ and B be $n_2 \times m_2$, so that $A \otimes B$ is $n \times m$, where $n = n_1 n_2$ and $m = m_1 m_2$. The

entries in $A \otimes B$ are products of a single entry in A and a single entry in B . Specifically, if $i = i_1 n_1 + i_2$ and $j = j_1 m_1 + j_2$, with $0 \leq i_2 < n_1$ and $0 \leq j_2 < m_1$, then $(A \otimes B)_{ij} = A_{i_1 j_1} B_{i_2 j_2}$.

$$\begin{aligned} (A \otimes B)^* &= [A_{i_1 j_1} B_{i_2 j_2}]^* && \text{(from above)} \\ &= [A_{i_1 j_1}^* B_{i_2 j_2}^*] && \text{(piecewise conjugation)} \\ &= A^* \otimes B^* && \text{(consistent indexing)} \end{aligned}$$

To see that $(A \otimes B)^T = A^T \otimes B^T$, note that $A^T \otimes B^T$ is $m \times n$, and $(A^T \otimes B^T)_{kl}$ is the product of a single entry in A^T and a single entry in B^T . Specifically, if $k = k_1 m_1 + k_2$ and $\ell = \ell_1 n_1 + \ell_2$, with $0 \leq k_2 < m_2$ and $0 \leq \ell_2 < n_2$, then $(A^T \otimes B^T)_{k\ell} = (A^T)_{k_1 \ell_1} (B^T)_{k_2 \ell_2} = A_{\ell_1 k_1} B_{\ell_2 k_2}$. Now, the hypotheses on k match the hypotheses on j above, and similarly for ℓ and i . This implies $(A^T \otimes B^T)_{k\ell} = (A \otimes B)_{\ell k} = (A \otimes B)_{k\ell}^T$. All entries in $A^T \otimes B^T$ and $(A \otimes B)^T$ are equal, so $(A \otimes B)^T = A^T \otimes B^T$.

Distributivity of † follows by applying distributivity of * and T in turn:

$$\begin{aligned} (A \otimes B)^\dagger &= ((A \otimes B)^*)^T && \text{(definition of } \dagger) \\ &= (A^* \otimes B^*)^T && \text{(distribute } *) \\ &= (A^*)^T \otimes (B^*)^T && \text{(distribute } ^T) \\ &= A^\dagger \otimes B^\dagger. && \text{(definition of } \dagger) \end{aligned}$$

2.29) Show that the tensor product of two unitary operators is unitary

Soln: Suppose U_1 and U_2 are unitary operators. To avoid implicit assumptions on multiplication of tensor products, let $|v\rangle$ and $|w\rangle$ be vectors in the spaces on which U_1 and U_2 operate. Then:

$$\begin{aligned} (U_1 \otimes U_2)(U_1 \otimes U_2)^\dagger(|v\rangle \otimes |w\rangle) &= (U_1 \otimes U_2)(U_1^\dagger \otimes U_2^\dagger)(|v\rangle \otimes |w\rangle) && \text{(distributivity of } \dagger) \\ &= (U_1 \otimes U_2)(U_1^\dagger |v\rangle \otimes U_2^\dagger |w\rangle) && \text{(definition of tensor product of operators)} \\ &= U_1 U_1^\dagger |v\rangle \otimes U_2 U_2^\dagger |w\rangle && \text{(definition of tensor product of operators)} \\ &= I |v\rangle \otimes I |w\rangle && (U_1 \text{ and } U_2 \text{ are unitary)} \\ &= (I \otimes I)(|v\rangle \otimes |w\rangle) && \text{(definition of tensor product of operators)} \\ &= I(|v\rangle \otimes |w\rangle) && (I \otimes I = I \text{ by construction)} \end{aligned}$$

So, $(U_1 \otimes U_2)(U_1 \otimes U_2)^\dagger = I$. Similarly, $(U_1 \otimes U_2)^\dagger(U_1 \otimes U_2) = I \otimes I = I$, so $U_1 \otimes U_2$ is unitary.

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

Soln: Suppose A and B are Hermitian operators. Then by distributivity of † and Hermiticity:

$$(A \otimes B)^\dagger = A^\dagger \otimes B^\dagger = A \otimes B.$$

Thus $A \otimes B$ is Hermitian.

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

Soln: Suppose A and B are positive operators. Then

$$\begin{aligned} (|\psi\rangle \otimes |\phi\rangle, (A \otimes B)(|\psi\rangle \otimes |\phi\rangle)) &= (|\psi\rangle \otimes |\phi\rangle, A|\psi\rangle \otimes B|\phi\rangle) && \text{(definition of } A \otimes B) \\ &= (|\psi\rangle \otimes |\phi\rangle)^\dagger (A|\psi\rangle \otimes B|\phi\rangle) && \text{(definition of inner-product)} \\ &= (\langle\psi| \otimes \langle\phi|)(A|\psi\rangle \otimes B|\phi\rangle) && \text{(distributivity of } \dagger) \\ &= \langle\psi|A|\psi\rangle \langle\phi|B|\phi\rangle. \end{aligned}$$

Since A and B are positive operators, $\langle \psi | A | \psi \rangle \geq 0$ and $\langle \phi | B | \phi \rangle \geq 0$ for all $|\psi\rangle, |\phi\rangle$, so $\langle \psi | A | \psi \rangle \langle \phi | B | \phi \rangle \geq 0$, from which we conclude that $A \otimes B$ is positive.

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

Soln: Suppose P_1 and P_2 are projectors. It is tempting to think that by applying exercise 2.16, which yields

$$\begin{aligned} (P_1 \otimes P_2)^2 &= P_1^2 \otimes P_2^2 && \text{(tensor product is multiplicative)} \\ &= P_1 \otimes P_2, && \text{(exercise 2.16)} \end{aligned}$$

exercise 2.16 would then imply that $P_1 \otimes P_2$ is also projector. However, this implication is the converse of exercise 2.16, which we have not proven. Instead, we need to prove that if $P_1 = \sum_{i=0}^k |v_i\rangle \langle v_i|$ and $P_2 = \sum_{j=0}^t |w_j\rangle \langle w_j|$, where $|v_i\rangle_{i=0}^k$ is a subset of an orthonormal basis $|v_i\rangle_{i=0}^\kappa$, and $|w_j\rangle_{j=0}^t$ is a subset of an orthonormal basis $|w_j\rangle_{j=0}^\tau$, then $P_1 \otimes P_2 = \sum_{q=0}^s |r_q\rangle \langle r_q|$, where $|r_q\rangle_{q=0}^s$ is a subset of an orthonormal basis $|r_q\rangle_{q=0}^\sigma$. First, the fact that $P_1 \otimes P_2 = \sum_{\substack{0 \leq i \leq k \\ 0 \leq j \leq t}} (|v_i\rangle \langle v_i|) \otimes (|w_j\rangle \langle w_j|)$ follows easily from distributivity of

operator tensor products, having illustrated how easily that follows from distributivity of tensor products of vectors in exercise 2.29. We need to show that $\left\{ (|v_i\rangle \langle v_i|) \otimes (|w_j\rangle \langle w_j|) \right\}_{\substack{0 \leq i \leq k \\ 0 \leq j \leq t}}$ is a subset of an orthonormal basis. It is automatically a subset of the set of vector tensor products resulting from loosening the restrictions on i and j to $0 \leq i \leq \kappa$ and $0 \leq j \leq \tau$, which we may assume is a basis, as stated on page 72. We need only show that the inner product of tensor products is multiplicative so that orthonormality is preserved. Let v_1, v_2, w_1 and w_2 be vectors.

$$\begin{aligned} \langle v_1 \otimes w_1 | v_2 \otimes w_2 \rangle &= \langle v_1 \otimes w_1 |^\dagger | v_2 \otimes w_2 \rangle && \text{(definition of } \langle \cdot | \cdot \rangle) \\ &= (|v_1\rangle^\dagger \otimes |w_1\rangle^\dagger) (|v_2\rangle \otimes |w_2\rangle) && \text{(distributivity of } \dagger \text{ over } \otimes) \\ &= (|v_1\rangle^\dagger |v_2\rangle) \otimes (|w_1\rangle^\dagger |w_2\rangle) && \text{(mixed-product property of Kronecker product)} \\ &= \langle v_1 | v_2 \rangle \otimes \langle w_1 | w_2 \rangle && \text{(definition of } \langle \cdot | \cdot \rangle) \\ &= \langle v_1 | v_2 \rangle \langle w_1 | w_2 \rangle && (\langle \cdot | \cdot \rangle \text{ is a scalar}) \end{aligned}$$

So, in the basis $\left\{ (|v_i\rangle \langle v_i|) \otimes (|w_j\rangle \langle w_j|) \right\}_{\substack{0 \leq i \leq \kappa \\ 0 \leq j \leq \tau}}$, the inner product of two vectors $|v_{i_1}\rangle \otimes |w_{j_1}\rangle$ and $|v_{i_2}\rangle \otimes |w_{j_2}\rangle$ is $\langle v_{i_1} \otimes w_{j_1} | v_{i_2} \otimes w_{j_2} \rangle = \langle v_{i_1} | v_{i_2} \rangle \langle w_{j_1} | w_{j_2} \rangle = \delta_{i_1 i_2} \delta_{j_1 j_2}$, from which it follows that this basis is orthonormal, completing the proof.

2.33) The Hadamard operator on one qubit may be written as

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

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 $H^{\otimes 2}$.

Soln: It is important to note what is meant by $x \cdot y$ in this formula. Here \cdot does **not** mean integer multiplication. It can be taken to mean popparity of the binary AND of x and y . Note that this property is multiplicative across dimensions, when 1 is used for even popparity, and -1 for odd. We proceed by

induction on n . Assume the preceding formula for $n - 1$. We must prove the formula for n .

$$\begin{aligned}
 H^{\otimes n} &= H \otimes H^{\otimes n-1} && \text{(notation)} \\
 &= \left(\frac{1}{\sqrt{2}} [(|0\rangle + |1\rangle) \langle 0| + (|0\rangle - |1\rangle) \langle 1|] \right) \otimes \left(\frac{1}{\sqrt{2^{n-1}}} \sum_{x,y} (-1)^{x \cdot y} |x\rangle \langle y| \right) && \text{(hypothesis)} \\
 &= \frac{1}{\sqrt{2^n}} \sum_{x,y} (-1)^{x \cdot y} (|0\rangle \langle 0| + |1\rangle \langle 0| + |0\rangle \langle 1| - |1\rangle \langle 1|) \otimes |x\rangle \langle y| && \text{(rearrange)} \\
 &= \frac{1}{\sqrt{2^n}} \sum_{x',y'} (-1)^{x' \cdot y'} |x'\rangle \langle y'| && \text{(multiplicativity of } \cdot \text{ across dimensions)}
 \end{aligned}$$

Now, explicitly

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

and

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

2.34) Find the square root and logarithm of the matrix

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

Soln: Suppose $A = \begin{bmatrix} 4 & 3 \\ 3 & 4 \end{bmatrix}$. We will need the “spectral” decomposition A . First, to find the eigenvalues of A :

$$\begin{aligned}
 0 &= \det(A - \lambda I) = (4 - \lambda)^2 - 3^2 \\
 &= \lambda^2 - 8\lambda + 7 \\
 &= (\lambda - 1)(\lambda - 7)
 \end{aligned}$$

So, the eigenvalues of A are $\lambda = 1$, and $\lambda = 7$. The corresponding eigenvectors can easily be seen to be $\begin{bmatrix} 1 \\ -1 \end{bmatrix}$, corresponding to $\lambda = 1$, and $\begin{bmatrix} 1 \\ 1 \end{bmatrix}$, corresponding to $\lambda = 7$. To construct an orthonormal basis from these eigenvectors, we can scale both by $\frac{1}{\sqrt{2}}$, and denote these scaled vectors by $|\lambda = 1\rangle = \frac{1}{\sqrt{2}} \begin{bmatrix} 1 \\ -1 \end{bmatrix}$ and $|\lambda = 7\rangle = \frac{1}{2} \begin{bmatrix} 1 \\ 1 \end{bmatrix}$. Now, seeing as A is real and self-transpose/adjoint, it is a normal matrix/operator, and as such, A can be “spectrally decomposed”/diagonalized as:

$$\begin{aligned}
 A &= |\lambda = 1\rangle \langle \lambda = 1| + 7 |\lambda = 7\rangle \langle \lambda = 7| \\
 &= \frac{1}{2} \begin{bmatrix} 1 & -1 \\ -1 & 1 \end{bmatrix} + 7 \left(\frac{1}{2} \begin{bmatrix} 1 & 1 \\ 1 & 1 \end{bmatrix} \right).
 \end{aligned}$$

The square root of A is “defined” as:

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

and $\log(A)$ is defined as

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

Note that one would hope that operator functions would respect various properties of the functions, such as inverses. To formally prove such a thing, let us consider the diagonal representation $A = \sum_a a |\lambda = a\rangle\langle\lambda = a|$ and the induced definition of $f(A) = \sum_a f(a) |\lambda = a\rangle\langle\lambda = a|$. Note that the $f(a)$ are eigenvalues of $f(A)$, with corresponding eigenvectors $|\lambda = a\rangle$, since the fact that the $|\lambda = a\rangle$ are orthonormal gives $(\sum_a f(a) |\lambda = a\rangle\langle\lambda = a|) |\lambda = a'\rangle = f(a') |\lambda = a'\rangle$. So $\sum_a f(a) |\lambda = a\rangle\langle\lambda = a|$ is a diagonal representation of $f(A)$. Now $f^{-1}(f(A)) = f^{-1}(\sum_a f(a) |\lambda = a\rangle\langle\lambda = a|) = \sum_a f^{-1}(f(a)) |\lambda = a\rangle\langle\lambda = a| = \sum_a a |\lambda = a\rangle\langle\lambda = a| = A$.

2.35) (Exponentiation of Pauli Matrices) let \vec{v} be any real three-dimensional unit vector and θ a real number. Prove that

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

where $\vec{v} \cdot \vec{\sigma} \equiv \sum_{i=1}^3 v_i \sigma_i$. This exercise is generalized in problem 2.1 on page 117.

Soln: To find the eigenvalues of $\vec{v} \cdot \vec{\sigma}$, we first express in matrix form:

$$\vec{v} \cdot \vec{\sigma} = \sum_{i=1}^3 v_i \sigma_i (= v_1 \sigma_x + v_2 \sigma_y + v_3 \sigma_z = v_1 X + v_2 Y + v_3 Z) \quad (\text{definition})$$

$$= v_1 \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} + v_2 \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} + v_3 \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} \quad (\text{substitute})$$

$$= \begin{bmatrix} v_3 & v_1 - iv_2 \\ v_1 + iv_2 & -v_3 \end{bmatrix} \quad (\text{collect terms})$$

$$0 = \det(\vec{v} \cdot \vec{\sigma} - \lambda I) = (v_3 - \lambda)(-v_3 - \lambda) - (v_1 - iv_2)(v_1 + iv_2) \quad (\text{characteristic equation})$$

$$= \lambda^2 - (v_1^2 + v_2^2 + v_3^2) \quad (\text{expand and collect } v\text{'s})$$

$$= \lambda^2 - 1 \quad (v \text{ is a unit vector})$$

Solving yields that the eigenvalues of $\vec{v} \cdot \vec{\sigma}$ are $\lambda = \pm 1$. Let $|\lambda_{\pm 1}\rangle$ be eigenvectors with eigenvalues ± 1 . Since \vec{v} is a real valued vector, it is easily seen that $\vec{v} \cdot \vec{\sigma}$ is Hermitian, and so is diagonalizable, and we may take the $|\lambda_{\pm 1}\rangle$ to be orthonormal by Exercise 2.22. In particular, this gives

$$|\lambda_1\rangle\langle\lambda_1| + |\lambda_{-1}\rangle\langle\lambda_{-1}| = I \quad (\text{completeness})$$

$$|\lambda_1\rangle\langle\lambda_1| - |\lambda_{-1}\rangle\langle\lambda_{-1}| = \vec{v} \cdot \vec{\sigma} \quad (\text{diagonalization})$$

Now

$$\exp(i\theta \vec{v} \cdot \vec{\sigma}) = e^{i\theta} |\lambda_1\rangle\langle\lambda_1| + e^{-i\theta} |\lambda_{-1}\rangle\langle\lambda_{-1}| \quad (\text{definition})$$

$$= (\cos \theta + i \sin \theta) |\lambda_1\rangle\langle\lambda_1| + (\cos \theta - i \sin \theta) |\lambda_{-1}\rangle\langle\lambda_{-1}| \quad (\text{Euler's formula})$$

$$= \cos \theta (|\lambda_1\rangle\langle\lambda_1| + |\lambda_{-1}\rangle\langle\lambda_{-1}|) + i \sin \theta (|\lambda_1\rangle\langle\lambda_1| - |\lambda_{-1}\rangle\langle\lambda_{-1}|) \quad (\text{group cos and sin terms})$$

$$= \cos(\theta)I + i \sin(\theta) \vec{v} \cdot \vec{\sigma}. \quad (\text{from above})$$

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

Soln:

$$\begin{aligned}\text{tr}(\sigma_1) &= \text{tr} \begin{pmatrix} 0 & 1 \\ 1 & 0 \end{pmatrix} = 0 \\ \text{tr}(\sigma_2) &= \text{tr} \begin{pmatrix} 0 & -i \\ i & 0 \end{pmatrix} = 0 \\ \text{tr}(\sigma_3) &= \text{tr} \begin{pmatrix} 1 & 0 \\ 0 & -1 \end{pmatrix} = 1 - 1 = 0\end{aligned}$$

2.37) (Cyclic property of the trace) If A and B are two linear operators show that

$$\text{tr}(AB) = \text{tr}(BA).$$

$$\begin{aligned}\text{tr}(AB) &= \sum_i \langle i|AB|i\rangle && \text{(using matrix representation of } AB\text{)} \\ &= \sum_i \langle i|AIB|i\rangle && \text{(insert } I\text{)} \\ &= \sum_{i,j} \langle i|A|j\rangle \langle j|B|i\rangle && \text{(completeness: } I = \sum_j |j\rangle\langle j|\text{)} \\ &= \sum_{i,j} \langle j|B|i\rangle \langle i|A|j\rangle && \text{(commutativity in } \mathbb{C}\text{)} \\ &= \sum_j \langle j|BA|j\rangle && \text{(completeness: } I = \sum_i |i\rangle\langle i|\text{)} \\ &= \text{tr}(BA) && \text{(using matrix representation of } BA\text{)}\end{aligned}$$

2.38) (Linearity of the trace) If A and B are two linear operators, show that

$$\text{tr}(A + B) = \text{tr}(A) + \text{tr}(B)$$

and if z is an arbitrary complex number show that

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

Soln:

$$\begin{aligned}\text{tr}(A + B) &= \sum_i \langle i|A + B|i\rangle && \text{(using matrix representation of } A + B\text{)} \\ &= \sum_i (\langle i|A|i\rangle + \langle i|B|i\rangle) && \text{(linearity of } \langle \cdot | \cdot \rangle\text{)} \\ &= \sum_i \langle i|A|i\rangle + \sum_i \langle i|B|i\rangle && \text{(separate terms)} \\ &= \text{tr}(A) + \text{tr}(B). && \text{(using matrix representation of } A \text{ and } B\text{)}\end{aligned}$$

$$\begin{aligned}
\text{tr}(zA) &= \sum_i \langle i|zA|i\rangle && \text{(matrix representation)} \\
&= \sum_i z \langle i|A|i\rangle && \text{(linearity)} \\
&= z \sum_i \langle i|A|i\rangle && \text{(linearity of sum)} \\
&= z \text{tr}(A). && \text{(matrix representation)}
\end{aligned}$$

2.39) (The Hilbert-Schmidt inner product on operators) The set L_V of linear operators on a Hilbert space V is a vector space. 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

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

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 .

Soln: (1) We check the three properties of inner products on page 65:

(i) linearity in the second argument:

$$\begin{aligned}
\left(A, \sum_i \lambda_i B_i\right) &= \text{tr}\left(A^\dagger \left(\sum_i \lambda_i B_i\right)\right) && \text{(definition of } (\cdot, \cdot)) \\
&= \text{tr}\left(\sum_i \lambda_i A^\dagger B_i\right) && \text{(linearity in } L_V) \\
&= \sum_i \lambda_i \text{tr}(A^\dagger B_i) && \text{(linearity of tr from Exercise 2.38)} \\
&= \sum_i \lambda_i (A, B_i) \quad \square && \text{(definition of } (\cdot, \cdot))
\end{aligned}$$

(ii) conjugate symmetry:

$$\begin{aligned}
(A, B)^* &= \left(\text{tr}(A^\dagger B)\right)^* && \text{(definition)} \\
&= \left(\sum_{i,j} \langle i|A^\dagger|j\rangle \langle j|B|i\rangle\right)^* && \text{(matrix representation, insert } I, \text{ apply completeness)} \\
&= \sum_{i,j} \langle j|B|i\rangle^* \langle i|A^\dagger|j\rangle^* && \text{(distributivity and multiplicativity of } *, \text{ commutativity in } \mathbb{C}) \\
&= \sum_{i,j} |j\rangle^{\dagger*} B^* |i\rangle^* |i\rangle^{\dagger*} A^{\dagger*} |j\rangle^* && \text{(distribute conjugation and make } \dagger\text{'s explicit)} \\
&= \sum_{i,j} |j\rangle^T B^{\dagger T} |i\rangle^{\dagger T} |i\rangle^T A^T |j\rangle^{\dagger T} && (\dagger^* = {}^T, {}^* = \dagger^T) \\
&= \sum_{i,j} (|i\rangle^\dagger B^\dagger |j\rangle)^T (|j\rangle^\dagger A |i\rangle)^T && ({}^T \text{ is cyclic)} \\
&= \sum_{i,j} \langle i|B^\dagger|j\rangle \langle j|A|i\rangle && \text{(definition of } \langle \cdot | \cdot \rangle, \text{ transpose in } \mathbb{C}) \\
&= \sum_i \langle i|B^\dagger A|i\rangle && \text{(completeness)} \\
&= \text{tr}(B^\dagger A) = (B, A) && \text{(matrix representation, definition of } (\cdot, \cdot))
\end{aligned}$$

(iii) positivity:

$$\begin{aligned}
 (A, A) &= \text{tr}(A^\dagger A) && \text{(definition)} \\
 &= \sum_i \langle i | A^\dagger A | i \rangle && \text{(matrix representation)} \\
 &\geq 0 && (A^\dagger A \text{ is positive by Exercise 2.25})
 \end{aligned}$$

We are left only show that if $(A, A) = 0$, then $A = 0$. Note that $(A, A) = 0$ implies that $\langle i | A^\dagger A | i \rangle = 0$ for all i , where here, the $|i\rangle$ are an orthonormal *basis* with respect to which the matrix representation of A is constructed. Note that $A |i\rangle$ is the i -th column of A , as constructed on page 64. Also $\|A |i\rangle\|^2 = \langle i | A^\dagger A | i \rangle = 0$, so the i -th column of A is 0, for all i . This gives that $(A, A) = 0$ iff $A = 0$, completing the proof of positivity.

(2) To show that L_V has dimension d^2 , note that the elements of L_V are linear operators from V to V and thus have a matrix representation by a $d \times d$ matrix. The vector space of $d \times d$ matrices over \mathbb{C} is clearly at most d^2 -dimensional. An easily constructed basis is the set of matrices populated with all 0s except for a single 1, where the position of the 1's ranges over all d^2 positions. This set is clearly linearly independent and spans L_V , from which dimension d^2 follows.. More generally, if $|k\rangle_k = 0^{d-1}$ is any orthonormal basis for V , then $\{|k\rangle\langle\ell|\}_{k,\ell}$ is a basis. The first example is produced using the standard (computational) basis.

(3) The obvious choices of basis matrices discussed above are orthonormal, but not Hermitian. Note that for $d = 2$, the Pauli matrices are orthonormal with respect to the Hilbert-Schmidt inner product, and are Hermitian. Other bases can be constructed using a “discrete Weyl system”. See Example 1.6 on page 11 of the lecture notes titled “Mathematical Introduction to Quantum Information Processing”, from Professor Michael Wolf of the Technical University of Munich (https://www-m5.ma.tum.de/foswiki/pub/M5/Allgemeines/MA5057_2019S/QIPlecture.pdf), and a convenient diagram included in “Holevo Capacity of Discrete Weyl Channels”, by Rehman et. al (<https://www.nature.com/articles/s41598-018-35777-7>). Unfortunately, neither of those constructions are Hermitian. To construct an orthonormal basis of Hermitian matrices, we use a related construction. First, let $|k\rangle_{k=0}^{d-1}$ be an orthonormal basis for V . The matrices we construct, $U_{k,\ell}$, will be doubly indexed by elements of the basis. When $|k\rangle_{k=0}^{d-1}$ is the standard basis, each will have two non-zero entries, unless $k = \ell$, in which case there will be a single non-zero entry. Define

$$\alpha_{k,\ell} := \begin{cases} \frac{1}{\sqrt{2}} & k < \ell \\ \frac{1}{2} & k = \ell \\ \frac{i}{\sqrt{2}} & k > \ell \end{cases}$$

Now let $U_{k,\ell} = \alpha_{k,\ell} |k\rangle\langle\ell| + \alpha_{k,\ell}^* |\ell\rangle\langle k|$. When $k = \ell$, $|k\rangle\langle\ell|$ and $|\ell\rangle\langle k|$ coincide, placing $\alpha_{k,\ell} + \alpha_{k,\ell}^* = 1$ somewhere along the diagonal. More generally, for fixed k, ℓ

$$\begin{aligned}
 U_{k,\ell}^\dagger &= \left(\alpha_{k,\ell} |k\rangle\langle\ell| + \alpha_{k,\ell}^* |\ell\rangle\langle k| \right)^\dagger && \text{(definition)} \\
 &= \alpha_{k,\ell}^* \langle\ell|^\dagger |k\rangle^\dagger + \alpha_{k,\ell} \langle\ell|^\dagger |k\rangle^\dagger && \text{(properties of } \dagger \text{)} \\
 &= \alpha_{k,\ell}^* |\ell\rangle\langle k| + \alpha_{k,\ell} |k\rangle\langle\ell| && \text{(simplify } \dagger \text{)} \\
 &= U_{k,\ell},
 \end{aligned}$$

so each $U_{k,\ell}$ is Hermitian. To show they are orthonormal, consider the inner-product of two arbitrary matrices $U_{k,\ell}$ and $U_{n,m}$.

$$\begin{aligned}
 (U_{k,\ell}, U_{n,m}) &= \text{tr}(U_{k,\ell}^\dagger U_{n,m}) && \text{(definition of } (\cdot, \cdot) \text{)} \\
 &= \text{tr}(U_{k,\ell} U_{n,m}) && \text{(Hermiticity of } U_{k,\ell} \text{)} \\
 &= \text{tr} \left(\left(\alpha_{k,\ell} |k\rangle\langle\ell| + \alpha_{k,\ell}^* |\ell\rangle\langle k| \right) \left(\alpha_{n,m} |n\rangle\langle m| + \alpha_{n,m}^* |m\rangle\langle n| \right) \right) && \text{(definition of } U_{i,j} \text{)}
 \end{aligned}$$

$$\begin{aligned}
&= \text{tr} \left(\begin{array}{cc} \alpha_{k,\ell} \alpha_{n,m} |k\rangle \langle \ell| n\rangle \langle m| & + \alpha_{k,\ell} \alpha_{n,m}^* |k\rangle \langle \ell| m\rangle \langle n| \\ + \alpha_{k,\ell}^* \alpha_{n,m} |\ell\rangle \langle k| n\rangle \langle m| & + \alpha_{k,\ell}^* \alpha_{n,m}^* |\ell\rangle \langle k| m\rangle \langle n| \end{array} \right) \quad (\text{F.O.I.L.}) \\
&= \begin{array}{cc} \alpha_{k,\ell} \alpha_{n,m} \delta_{\ell,n} \text{tr}(|k\rangle \langle m|) & + \alpha_{k,\ell} \alpha_{n,m}^* \delta_{\ell,m} \text{tr}(|k\rangle \langle n|) \\ + \alpha_{k,\ell}^* \alpha_{n,m} \delta_{k,n} \text{tr}(|\ell\rangle \langle m|) & + \alpha_{k,\ell}^* \alpha_{n,m}^* \delta_{k,m} \text{tr}(|\ell\rangle \langle n|) \end{array} \quad (\text{linearity, orthonormality}) \\
&= \begin{array}{cc} \alpha_{k,\ell} \alpha_{n,m} \delta_{\ell,n} \text{tr}(\langle m| k\rangle) & + \alpha_{k,\ell} \alpha_{n,m}^* \delta_{\ell,m} \text{tr}(\langle n| k\rangle) \\ + \alpha_{k,\ell}^* \alpha_{n,m} \delta_{k,n} \text{tr}(\langle m| \ell\rangle) & + \alpha_{k,\ell}^* \alpha_{n,m}^* \delta_{k,m} \text{tr}(\langle n| \ell\rangle) \end{array} \quad (\text{cyclicity of tr}) \\
&= \begin{array}{cc} \alpha_{k,\ell} \alpha_{n,m} \delta_{\ell,n} \delta_{k,m} & + \alpha_{k,\ell} \alpha_{n,m}^* \delta_{\ell,m} \delta_{k,n} \\ + \alpha_{k,\ell}^* \alpha_{n,m} \delta_{k,n} \delta_{\ell,m} & + \alpha_{k,\ell}^* \alpha_{n,m}^* \delta_{k,m} \delta_{\ell,n} \end{array} \quad (\text{orthonormality}) \\
&= \begin{array}{cc} (\alpha_{k,\ell} \alpha_{n,m} + (\alpha_{k,\ell} \alpha_{n,m})^*) \delta_{k,m} \delta_{\ell,n} & \\ + (\alpha_{k,\ell} \alpha_{n,m}^* + (\alpha_{k,\ell} \alpha_{n,m}^*)^*) \delta_{k,n} \delta_{\ell,m} & \end{array} \quad (\text{group like } \delta\text{s, property of } *) \\
&= \begin{array}{cc} 2\Re(\alpha_{k,\ell} \alpha_{n,m}) \delta_{k,m} \delta_{\ell,n} & \\ + 2\Re(\alpha_{k,\ell} \alpha_{n,m}^*) \delta_{k,n} \delta_{\ell,m} & \end{array} \quad (\text{property of } *)
\end{aligned}$$

When $k \neq n$ or $\ell \neq m$ only the first term contributes, so $(U_{k,\ell}, U_{n,m}) = 2\Re(\alpha_{k,\ell} \alpha_{n,m}) \delta_{k,m} \delta_{\ell,n}$. This can only be nonzero if $m = k$ and $n = \ell$, in which case $(U_{k,\ell}, U_{n,m}) = (U_{k,\ell}, U_{\ell,k}) = 2\Re(\alpha_{k,\ell} \alpha_{\ell,k}) = 2\Re(\frac{i}{2}) = 0$, since $k \neq n = \ell$ implies one of the $\alpha_{k,\ell}$ and $\alpha_{\ell,k}$ is $\frac{1}{\sqrt{2}}$, and the other $\frac{i}{\sqrt{2}}$. Hence, the $U_{k,\ell}$ are orthogonal. If $k = n$ and $\ell = m$, then $(U_{k,\ell}, U_{k,\ell}) = 2\Re(\alpha_{k,\ell} \alpha_{k,\ell}) \delta_{k,\ell}^2 + 2\Re(\alpha_{k,\ell} \alpha_{k,\ell}^*)$. There are two cases. If $k = \ell$, then $\alpha_{k,\ell} = \alpha_{k,\ell}^* = \frac{1}{2}$, so $(U_{k,\ell}, U_{k,\ell}) = 2(\frac{1}{4}) + 2(\frac{1}{4}) = 1$. When $k \neq \ell$, $(U_{k,\ell}, U_{k,\ell}) = 2\Re(\alpha_{k,\ell} \alpha_{k,\ell}^*) = 2\Re(\|\alpha_{k,\ell}\|^2) = 2(\frac{1}{2}) = 1$. So in all cases $(U_{k,\ell}, U_{k,\ell}) = 1$ and thus the $U_{k,\ell}$ are a set of d^2 Hermitian matrices which are orthonormal. Orthonormality implies linear independence, in which case the subspace spanned by the $U_{k,\ell}$ has dimension d^2 and so must be L_V itself, making $U_{k,\ell}$ a basis.

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 ϵ_{jkl} , the antisymmetric tensor on three indices, for which $\epsilon_{jkl} = 0$ except for $\epsilon_{123} = \epsilon_{231} = \epsilon_{312} = 1$ and $\epsilon_{321} = \epsilon_{213} = \epsilon_{132} = -1$:

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

Soln:

$$\begin{aligned}
[\sigma_1, \sigma_2] &= [X, Y] = XY - YX \\
&= \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} - \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} \\
&= \begin{bmatrix} i & 0 \\ 0 & -i \end{bmatrix} - \begin{bmatrix} -i & 0 \\ 0 & i \end{bmatrix} \\
&= \begin{bmatrix} 2i & 0 \\ 0 & -2i \end{bmatrix} \\
&= 2iZ = 2i\epsilon_{12}\sigma_3
\end{aligned}$$

$$\begin{aligned}
[\sigma_2, \sigma_3] &= [Y, Z] = \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} - \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} \\
&= \begin{bmatrix} 0 & 2i \\ 2i & 0 \end{bmatrix} \\
&= 2iX = 2i\epsilon_{231}\sigma_1
\end{aligned}$$

$$\begin{aligned}
[\sigma_3, \sigma_1] &= [Z, X] = \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} - \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} \\
&= \begin{bmatrix} 0 & 2 \\ -2 & 0 \end{bmatrix} \\
&= 2i \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} \\
&= 2iY = 2i\epsilon_{312}\sigma_2
\end{aligned}$$

2.41) (Anti-commutation relations for the Pauli Matrices) Verify that the anticommutation relations

$$\sigma_i, \sigma_j = 0$$

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

$$\sigma_i^2 = I.$$

Soln:

$$\begin{aligned}
\{X, Y\} &= \{\sigma_1, \sigma_2\} = \sigma_1\sigma_2 + \sigma_2\sigma_1 \\
&= \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} + \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} \\
&= \begin{bmatrix} i & 0 \\ 0 & -i \end{bmatrix} + \begin{bmatrix} -i & 0 \\ 0 & i \end{bmatrix} \\
&= 0
\end{aligned}$$

$$\begin{aligned}
\{Y, Z\} &= \{\sigma_2, \sigma_3\} = \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} + \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} \\
&= 0
\end{aligned}$$

$$\begin{aligned}
\{Z, X\} &= \{\sigma_3, \sigma_1\} = \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} + \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} \\
&= 0
\end{aligned}$$

$$\begin{aligned}\sigma_0^2 &= I^2 = I \\ X^2 &= \sigma_1^2 = \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix}^2 = I \\ Y^2 &= \sigma_2^2 = \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix}^2 = I \\ Z^2 &= \sigma_3^2 = \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix}^2 = I\end{aligned}$$

2.42) Verify that

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

Soln:

$$\frac{[A, B] + \{A, B\}}{2} = \frac{AB - BA + AB + BA}{2} = AB \quad (\text{definition of } [\cdot, \cdot] \text{ and } \{\cdot, \cdot\})$$

2.43) Show that for $j, k = 1, 2, 3$,

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

From Exercises 2.41 (eq (2.75), (2.76)), $\{\sigma_j, \sigma_k\} = 2\delta_{jk} I$.

$$\begin{aligned}\sigma_j \sigma_k &= \frac{[\sigma_j, \sigma_k] + \{\sigma_j, \sigma_k\}}{2} && (\text{Exercise 2.42 (eq 2.77)}) \\ &= \frac{2i \sum_{\ell=1}^3 \epsilon_{jkl} \sigma_\ell + 2\delta_{jk} I}{2} && (\text{Exercise 2.40 (eq 2.74) and above}) \\ &= \delta_{jk} I + i \sum_{\ell=1}^3 \epsilon_{jkl} \sigma_\ell && (\text{cancel 2s})\end{aligned}$$

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

Soln: By assumption, $[A, B] = AB - BA = 0$ implies $AB = BA$, now $\{A, B\} = AB + BA = 2AB = 0$, so $AB = 0$. Since A is invertible, multiplying by A^{-1} from the left gives

$$\begin{aligned}A^{-1}AB &= 0 && (A \text{ is invertible, multiply both sides by } A^{-1}) \\ IB &= 0 && (A^{-1}A = I) \\ B &= 0.\end{aligned}$$

2.45) Show that $[A, B]^\dagger = [B^\dagger, A^\dagger]$.

Soln:

$$\begin{aligned}[A, B]^\dagger &= (AB - BA)^\dagger && (\text{definition of } [\cdot, \cdot]) \\ &= B^\dagger A^\dagger - A^\dagger B^\dagger && (\text{properties of } \dagger) \\ &= [B^\dagger, A^\dagger] && (\text{definition of } [\cdot, \cdot])\end{aligned}$$

2.46) Show that $[A, B] = -[B, A]$.

Soln:

$$\begin{aligned} [A, B] &= AB - BA && \text{(definition of } [\cdot, \cdot]) \\ &= -(BA - AB) && \text{(reverse signs)} \\ &= -[B, A] && \text{(definition of } [\cdot, \cdot]) \end{aligned}$$

2.47) Suppose A and B are Hermitian. Show the $i[A, B]$ is Hermitian.

Soln:

$$\begin{aligned} (i[A, B])^\dagger &= -i[A, B]^\dagger && \text{(distribute } \dagger) \\ &= -i[B^\dagger, A^\dagger] && \text{(Exercise 2.45)} \\ &= -i[B, A] && (A \text{ and } B \text{ are Hermitian)} \\ &= i[A, B] && \text{(Exercise 2.46)} \end{aligned}$$

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

Soln:

(Positive) Since P is positive, it is Hermitian by Exercise 2.24, thus normal, so the spectral decomposition theorem gives that it is diagonalizable. Then $P = \sum_i \lambda_i |i\rangle\langle i|$, where $\lambda_i \geq 0$ by positivity. Note that this implies that $\sqrt{\lambda_i^2} = \lambda_i$.

$$\begin{aligned} J &= \sqrt{P^\dagger P} && \text{(uniqueness of } J) \\ &= \sqrt{P^2} && \text{(Hermiticity)} \\ &= \sum_i \sqrt{\lambda_i^2} |i\rangle\langle i| && \text{(spectral decomposition, definition of } \sqrt{\cdot}) \\ &= \sum_i \lambda_i |i\rangle\langle i| && (\sqrt{\cdot} \text{ in } \mathbb{R}) \\ &= P. && \text{(spectral decomposition)} \end{aligned}$$

Therefore, for any positive operator P , the polar decomposition of P is $P = UP$. *If P were positive *definite**, it would be easy to show that P is invertible. If $a = \sum_j a_j |j\rangle$, then if $Pa = 0$, $Pa = P\left(\sum_j a_j |j\rangle\right) = \left(\sum_i \lambda_i |i\rangle\langle i|\right)\left(\sum_j a_j |j\rangle\right) = \sum_i \lambda_i a_i |i\rangle = 0$. Since $|i\rangle$ is a basis, we must have $a_i \lambda_i = 0$ for all i , but the λ_i cannot be 0 since, being a basis vector, $|i\rangle \neq 0$ implies $\langle i|P|i\rangle = \lambda_i > 0$ by positive definiteness. Now all $a_i = 0$. Having 0-dimensional null-space, P must be invertible, in which case $U = I$ by uniqueness of U , since I satisfies $P = IP$. Then $P = P$ is the polar decomposition of P .

(Unitary) Suppose unitary U is decomposed by $U = WJ$ where W is unitary and J is positive, $J = \sqrt{U^\dagger U}$.

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

Since unitary operators are invertible, $W = UJ^{-1} = UI^{-1} = UI = U$. Thus, the polar decomposition of U is $U = U$.

Alternatively, since J is unique, note that U satisfies $U = UJ$, where $J = I$, so $U = U$ is again the polar decomposition of U . This is essentially the same argument, as above, where we use the stated uniqueness of J instead of the unique formula for it.

(Hermitian) Suppose $H = UJ$. By Hermiticity

$$J = \sqrt{H^\dagger H} = \sqrt{H H} = \sqrt{H^2}.$$

Thus $H = U\sqrt{H^2}$.

In general, $H \neq \sqrt{H^2}$.

From spectral decomposition, $H = \sum_i \lambda_i |i\rangle\langle i|$, $\lambda_i \in \mathbb{R}$.

$$\sqrt{H^2} = \sqrt{\sum_i \lambda_i^2 |i\rangle\langle i|} = \sum_i \sqrt{\lambda_i^2} |i\rangle\langle i| = \sum_i |\lambda_i| |i\rangle\langle i| \neq H$$

unless H is positive.

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

Soln: Let A be a normal matrix, which by the spectral decomposition theorem we may write $A = \sum_i \lambda_i |i\rangle\langle i|$ for an orthonormal basis $|i\rangle$. As in the proof of the polar decomposition theorem, define $|e_i\rangle = A|i\rangle$ for those $|i\rangle$ for which $\lambda_i \neq 0$, and extend via Gram-Schmidt to find $|e_i\rangle$ for those $|i\rangle$ for which $\lambda_i = 0$.

$$J = \sqrt{A^\dagger A} = \sum_i |\lambda_i| |i\rangle\langle i|. \quad (\text{uniqueness})$$

$$U = \sum_i |e_i\rangle\langle i|. \quad (\text{as constructed in the proof})$$

$$A = UJ \quad (\text{polar decomposition})$$

$$= \left(\sum_i |e_i\rangle\langle i| \right) \left(\sum_j |\lambda_j| |j\rangle\langle j| \right) \quad (\text{defined above})$$

$$= \sum_i |\lambda_i| |e_i\rangle\langle i|. \quad (\text{orthonormality})$$

2.50) Find the left and right polar decomposition of the matrix

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

Soln: We have $A^\dagger A = \begin{bmatrix} 1 & 1 \\ 0 & 1 \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 1 & 1 \end{bmatrix} = \begin{bmatrix} 2 & 1 \\ 1 & 1 \end{bmatrix}$. To construct J , we must find a spectral decomposition of $A^\dagger A$. The characteristic equation of $A^\dagger A$ is $\det(A^\dagger A - \lambda I) = (2 - \lambda)(1 - \lambda) - 1 = \lambda^2 - 3\lambda + 1 = 0$. By the quadratic formula, the eigenvalues of $A^\dagger A$ are $\lambda_\pm = \frac{3 \pm \sqrt{5}}{2}$. The associated orthonormal eigenvectors are $|\lambda_\pm\rangle = \frac{5 \pm \sqrt{5}}{10} \begin{bmatrix} \frac{1 \pm \sqrt{5}}{2} \\ 1 \end{bmatrix}$. We have $|\lambda_\pm\rangle\langle\lambda_\pm| = \begin{bmatrix} 1 \pm \frac{2\sqrt{5}}{5} & \frac{1}{2} \pm \frac{3\sqrt{5}}{10} \\ \frac{1}{2} \pm \frac{3\sqrt{5}}{10} & \frac{1}{2} \pm \frac{\sqrt{5}}{10} \end{bmatrix}$. By the spectral decomposition theorem, $A^\dagger A = \lambda_+ |\lambda_+\rangle\langle\lambda_+| + \lambda_- |\lambda_-\rangle\langle\lambda_-|$, and

$$\begin{aligned} J = \sqrt{A^\dagger A} &= \sqrt{\lambda_+} |\lambda_+\rangle\langle\lambda_+| + \sqrt{\lambda_-} |\lambda_-\rangle\langle\lambda_-| \quad (\text{definition of } \sqrt{\cdot}) \\ &= \sum_{\pm} \sqrt{\frac{3 \pm \sqrt{5}}{2}} \begin{bmatrix} 1 \pm \frac{2\sqrt{5}}{5} & \frac{1}{2} \pm \frac{3\sqrt{5}}{10} \\ \frac{1}{2} \pm \frac{3\sqrt{5}}{10} & \frac{1}{2} \pm \frac{\sqrt{5}}{10} \end{bmatrix} \end{aligned}$$

$$J^{-1} = \frac{1}{\sqrt{\lambda_+}} |\lambda_+\rangle\langle\lambda_+| + \frac{1}{\sqrt{\lambda_-}} |\lambda_-\rangle\langle\lambda_-|. \quad (\text{Applying the }^{-1} \text{ function to } J)$$

$$U = AJ^{-1} \quad (A \text{ is invertible, } U \text{ is unique})$$

It is not worth simplifying J, J^{-1} , or U , nor is it worth finding the right polar decomposition.

2.51) Verify that the Hadamard gate H is unitary.

Soln:

$$H^\dagger H = \left(\frac{1}{\sqrt{2}} \begin{bmatrix} 1 & 1 \\ 1 & -1 \end{bmatrix} \right)^\dagger \frac{1}{\sqrt{2}} \begin{bmatrix} 1 & 1 \\ 1 & -1 \end{bmatrix} = \frac{1}{\sqrt{2}} \begin{bmatrix} 1 & 1 \\ 1 & -1 \end{bmatrix} \frac{1}{\sqrt{2}} \begin{bmatrix} 1 & 1 \\ 1 & -1 \end{bmatrix} = \frac{1}{2} \begin{bmatrix} 2 & 0 \\ 0 & 2 \end{bmatrix} = I.$$

2.52) Verify that $H^2 = I$.

Soln: It was shown above that $H^\dagger = H$, so $H^2 = H^\dagger H = I$.

2.53) What are the eigenvalues and eigenvectors of H ?

Soln: Using the characteristic equation

$$\begin{aligned} \det(H - \lambda I) &= \left(\frac{1}{\sqrt{2}} - \lambda \right) \left(-\frac{1}{\sqrt{2}} - \lambda \right) - \frac{1}{2} \\ &= \lambda^2 - \frac{1}{2} - \frac{1}{2} \\ &= \lambda^2 - 1 = 0, \end{aligned}$$

the eigenvalues are $\lambda_\pm = \pm 1$. The associated orthonormal eigenvectors can be calculated to be $|\lambda_\pm\rangle = \frac{1}{\sqrt{4 \mp 2\sqrt{2}}} \begin{bmatrix} 1 \\ -1 \pm \sqrt{2} \end{bmatrix}$.

2.54) Suppose A and B are commuting Hermitian operators. Prove that $\exp(A)\exp(B) = \exp(A+B)$.

Soln: Since $[A, B] = 0$, A and B are simultaneously diagonalizable. Let $|i\rangle$ be an orthonormal basis such that $A = \sum_i a_i |i\rangle\langle i|$, $B = \sum_i b_i |i\rangle\langle i|$. Note that $A+B$ is also simultaneously diagonalizable, since $A+B = \sum_i (a_i + b_i) |i\rangle\langle i|$.

$$\begin{aligned} \exp(A)\exp(B) &= \left(\sum_i \exp(a_i) |i\rangle\langle i| \right) \left(\sum_i \exp(b_i) |i\rangle\langle i| \right) && \text{(from above)} \\ &= \sum_{i,j} \exp(a_i + b_j) |i\rangle\langle i|j\rangle\langle j| && \text{(group sum)} \\ &= \sum_{i,j} \exp(a_i + b_j) |i\rangle\langle j| \delta_{i,j} && (|i\rangle \text{ is orthonormal}) \\ &= \sum_i \exp(a_i + b_i) |i\rangle\langle i| && \text{(group non-zero terms)} \\ &= \exp(A+B). && \text{(definition of } \exp(\cdot) \text{)} \end{aligned}$$

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

Soln: Some definitions: H is the Hamiltonian of some closed system. It is Hermitian, and hence has a spectral decomposition: $H = \sum E |E\rangle\langle E|$. Note that $H^\dagger = (\sum E |E\rangle\langle E|)^\dagger = \sum E^* (|E\rangle\langle E|)^\dagger = \sum E^* |E\rangle\langle E|$ by Exercise 2.13. U is defined as $U(t_1, t_2) \equiv \exp\left[\frac{-iH(t_2-t_1)}{\hbar}\right]$. To prove U is unitary, we show $UU^\dagger = I$.

Now

$$\begin{aligned}
 (U(t_1, t_2))^\dagger &= \left(\exp \left[\frac{-iH(t_2 - t_1)}{\hbar} \right] \right)^\dagger && \text{(definition of } U) \\
 &= \left(\sum_E \exp \left(\frac{-iE(t_2 - t_1)}{\hbar} \right) |E\rangle\langle E| \right)^\dagger && \text{(definition of } \exp(A)) \\
 &= \sum_E \exp \left(\frac{-iE(t_2 - t_1)}{\hbar} \right)^* |E\rangle\langle E|^\dagger && \text{(linearity of } \dagger) \\
 &= \sum_E \exp \left(\frac{iE(t_2 - t_1)}{\hbar} \right) |E\rangle\langle E| && \text{(complex conjugation, exercise 2.13)} \\
 &= \exp \left(\frac{iH(t_2 - t_1)}{\hbar} \right) && \text{(definition of } \exp(A))
 \end{aligned}$$

We have

$$\begin{aligned}
 U(t_2 - t_1)(U(t_2 - t_1))^\dagger &= \exp \left(-\frac{iH(t_2 - t_1)}{\hbar} \right) \exp \left(\frac{iH(t_2 - t_1)}{\hbar} \right) && \text{(from above)} \\
 &= \sum_{E, E'} \left(\exp \left(\frac{-iE(t_2 - t_1)}{\hbar} \right) |E\rangle\langle E| \right) \left(\exp \left(\frac{iE'(t_2 - t_1)}{\hbar} \right) |E'\rangle\langle E'| \right) && \\
 &&& (E' \text{ used to distinguish from } E) \\
 &= \sum_{E, E'} \exp \left(-\frac{i(E - E')(t_2 - t_1)}{\hbar} \right) |E\rangle\langle E'| \delta_{E, E'} && \text{(orthonormality)} \\
 &= \sum_E \exp(0) |E\rangle\langle E| && \text{(group nonzero terms)} \\
 &= \sum_E |E\rangle\langle E| && (e^0 = 1) \\
 &= I && \text{(completeness)}
 \end{aligned}$$

Similarly, $(U(t_2 - t_1))^\dagger U(t_2 - t_1) = I$. So, U is unitary.

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 .

Soln: Since U is unitary, it has a spectral decomposition with respect to some orthonormal basis, say $|\lambda\rangle$. For each eigenvalue λ , note that exercise 2.18 gives that $\|\lambda\| = 1$, so express $\lambda = e^{i\theta}$ for some real valued argument θ . Then $\log(U) = \sum \log(\lambda) |\lambda\rangle\langle\lambda| = \sum i\theta |\lambda\rangle\langle\lambda|$. Now $K \equiv -i \log(U) = \sum \theta |\lambda\rangle\langle\lambda|$, and $K^\dagger = (\sum \theta |\lambda\rangle\langle\lambda|)^\dagger = \sum \theta^* (|\lambda\rangle\langle\lambda|)^\dagger = \sum \theta |\lambda\rangle\langle\lambda| = K$, since θ is real and exercise 2.13 gives that $(|\lambda\rangle\langle\lambda|)^\dagger = |\lambda\rangle\langle\lambda|$. So, K is Hermitian, and by applying the fact that \exp and \log are inverse complex valued functions, at least for the λ , we can conclude that $\exp(iK) = \exp(i \cdot (-i \log(U))) = \exp(\log(U)) = U$.

2.57) (Cascaded measurements are single measurements) Suppose $\{L_\ell\}$ and $\{M_m\}$ are two sets of measurement operators. Show that a measurement defined by the measurement operators $\{L_\ell\}$ followed by a measurement defined by the measurement operators $\{M_m\}$ is physically equivalent to a single measurement defined by measurement operators $\{N_{\ell m}\}$ with the representation $N_{\ell m} \equiv M_m L_\ell$.

Soln: Let the state of a physical system be ψ . Note that by definition the state after applying measurement operators $\{L_\ell\}$ is $|\phi\rangle \equiv \frac{L_\ell |\psi\rangle}{\sqrt{\langle\psi|L_\ell^\dagger L_\ell|\psi\rangle}}$ for some ℓ . The state after applying measurement operators

$\{M_m\}$ is $\xi \equiv \frac{M_m|\phi\rangle}{\sqrt{\langle\phi|M_m^\dagger M_m|\phi\rangle}}$ for some m . Now

$$\begin{aligned}
 \xi &= \frac{M_m|\phi\rangle}{\sqrt{\langle\phi|M_m^\dagger M_m|\phi\rangle}} && \text{(from above)} \\
 &= \frac{M_m \frac{L_\ell|\psi\rangle}{\sqrt{\langle\psi|L_\ell^\dagger L_\ell|\psi\rangle}}}{\sqrt{\left\langle \frac{L_\ell|\psi\rangle}{\sqrt{\langle\psi|L_\ell^\dagger L_\ell|\psi\rangle}} \middle| M_m^\dagger M_m \middle| \frac{L_\ell|\psi\rangle}{\sqrt{\langle\psi|L_\ell^\dagger L_\ell|\psi\rangle}} \right\rangle}} && \text{(substitute for } \phi \text{)} \\
 &= \frac{\frac{M_m L_\ell|\psi\rangle}{\sqrt{\langle\psi|L_\ell^\dagger L_\ell|\psi\rangle}}}{\sqrt{\frac{\langle L_\ell|\psi| M_m^\dagger M_m | L_\ell|\psi\rangle}{\langle\psi|L_\ell^\dagger L_\ell|\psi\rangle}}} && \text{(group internal scalar values)} \\
 &= \frac{\frac{\sqrt{\langle\psi|L_\ell^\dagger L_\ell|\psi\rangle}}{1} \cdot \frac{M_m L_\ell|\psi\rangle}{\sqrt{\langle\psi|L_\ell^\dagger L_\ell|\psi\rangle}}}{\sqrt{\langle\psi|L_\ell^\dagger M_m^\dagger M_m L_\ell|\psi\rangle}} && \text{(move the scalar to numerator)} \\
 &= \frac{(M_m L_\ell)\psi}{\sqrt{\langle\psi|(M_m L_\ell)^\dagger (M_m L_\ell)|\psi\rangle}} && \text{(cancel scalars, group operators)} \\
 &= \frac{N_{\ell m}|\psi\rangle}{\sqrt{\langle\psi|N_{\ell m}^\dagger N_{\ell m}|\psi\rangle}} && \text{(definition of } N_{\ell m} \text{)}
 \end{aligned}$$

This gives that the state of the system after applying measurement operators $\{L_\ell\}$ followed by measurement operators $\{M_m\}$ can be expressed in terms of applying measurement operators $N_{\ell m} = M_m L_\ell$. Note, though, that the state of the system after applying $N_{\ell m}$ is exactly the state of the system after applying M_m , for any ℓ . In practice, the intermediate measurement result ℓ would be unknown, and to find the probability that the system is in state m after application of the measurement operators $\{N_{\ell m}\}$, one would have to sum over the possible intermediate measurement results.

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?

Soln: Express $M = \sum_\mu \mu P_\mu$, where the $|\mu\rangle$ are an orthonormal set of eigenvectors each with eigenvalue μ . We may assume $\| |\psi\rangle \| = 1$, so that $|\psi\rangle = |\mu\rangle$ for some $|\mu\rangle$ (and $m = \mu$ for some μ).

$$\begin{aligned}
 \mathbf{E}(M) &= \langle M \rangle = \langle \psi | M | \psi \rangle && \text{(definition of } \mathbf{E}(M) \text{)} \\
 &= \left\langle \psi \middle| \sum_\mu \mu P_\mu \middle| \psi \right\rangle && \text{(spectral decomposition of } M \text{)} \\
 &= \sum_\mu \mu \langle \psi | P_\mu | \psi \rangle && \text{(linearity)} \\
 &= \sum_\mu \mu \langle \psi | \mu \rangle \langle \mu | \psi \rangle && (P_\mu \text{ is a projector)} \\
 &= \sum_\mu \mu \delta_{\psi,\mu}^2 && \text{(orthonormality)} \\
 &= m && \text{(collect nonzero terms)}
 \end{aligned}$$

Similarly, $\langle M^2 \rangle = \langle \psi | M^2 | \psi \rangle = m^2$. Now $\Delta(M) = \sqrt{\langle M^2 \rangle - \langle M \rangle^2} = \sqrt{m^2 - (m)^2} = 0$.

2.59) Suppose we have (a) qubit in the state $|0\rangle$, and we measure the observable X . What is the average

value of X ? What is the standard deviation?

Soln: There are two ways to proceed. First, we can apply X , noting that $X|0\rangle = |1\rangle$ and $X|1\rangle = |0\rangle$.

$$\begin{aligned}\langle X \rangle &= \langle 0|X|0\rangle = \langle 0|1\rangle = 0 \\ \langle X^2 \rangle &= \langle 0|X^2|0\rangle = \langle 0|X|1\rangle = \langle 0|0\rangle = 1 \\ \Delta(X) &= \sqrt{\langle X^2 \rangle - \langle X \rangle^2} = \sqrt{1 - 0^2} = 1\end{aligned}$$

Alternatively, write $X = (|+\rangle\langle +|) - (|-\rangle\langle -|)$, and note that $X^2 = (|+\rangle\langle +|) + (|-\rangle\langle -|)$. Now

$$\begin{aligned}\langle X \rangle &= \langle 0|X|0\rangle = \langle 0|+\rangle\langle +|0\rangle - \langle 0|-\rangle\langle -|0\rangle = \frac{1}{2} - \frac{1}{2} = 0 \\ \langle X^2 \rangle &= \langle 0|X^2|0\rangle = \langle 0|+\rangle\langle +|0\rangle + \langle 0|-\rangle\langle -|0\rangle = \frac{1}{2} + \frac{1}{2} = 1 \\ \Delta(X) &= \sqrt{\langle X^2 \rangle - \langle X \rangle^2} = \sqrt{1 - 0^2} = 1\end{aligned}$$

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

Soln: We calculate eigenvalues by expressing $\vec{v} \cdot \vec{\sigma}$ explicitly in terms of v_1, v_2 , and v_3 .

$$\vec{v} \cdot \vec{\sigma} = \sum_{i=1}^3 v_i \sigma_i \quad (\text{definition})$$

$$= v_1 \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} + v_2 \begin{bmatrix} 0 & -i \\ i & 0 \end{bmatrix} + v_3 \begin{bmatrix} 1 & 0 \\ 0 & -1 \end{bmatrix} \quad (\sigma_1 = X, \sigma_2 = Y, \sigma_3 = Z)$$

$$= \begin{bmatrix} v_3 & v_1 - iv_2 \\ v_1 + iv_2 & -v_3 \end{bmatrix} \quad (\text{add corresponding entries})$$

$$\det(\vec{v} \cdot \vec{\sigma} - \lambda I) = (v_3 - \lambda)(-v_3 - \lambda) - (v_1 - iv_2)(v_1 + iv_2) \quad (\text{characteristic equation})$$

$$= \lambda^2 - (v_1^2 + v_2^2 + v_3^2) \quad (\text{simplify})$$

$$= \lambda^2 - 1 \quad (\vec{v} \text{ is a unit vector})$$

So the eigenvalues of $\vec{v} \cdot \vec{\sigma}$ are $\lambda = \pm 1$.

To calculate the projectors onto the corresponding eigenspaces, note that $|\lambda_1\rangle = \begin{bmatrix} 1 \\ \frac{1-v_3}{v_1-iv_2} \end{bmatrix}$ can easily be seen to be an eigenvector with eigenvalue $\lambda = 1$. To normalize, we'll use the fact that $1 - v_3^2 = v_1^2 + v_2^2$, since \vec{v} is a unit vector. Factoring both sides and dividing yields $\frac{1 \pm v_3}{v_1 \pm iv_2} = \frac{v_1 \mp iv_2}{1 \mp v_3}$, that is, such rational functions can be flipped by negating both binary operations.

$$\| |\lambda_1\rangle \|^2 = 1 + \left\| \frac{1 - v_3}{v_1 - iv_2} \right\|^2 \quad (\text{definition of } \|\cdot\|^2)$$

$$= 1 + \left(\frac{v_1 + iv_2}{1 + v_3} \right) \left(\frac{v_1 - iv_2}{1 + v_3} \right) \quad (\text{flip, } \|c\|^2 = c \cdot c^*)$$

$$= 1 + \frac{v_1^2 + v_2^2}{(1 + v_3)^2} \quad (\text{expand})$$

$$= \frac{1 + 2v_3 + v_1^2 + v_2^2 + v_3^2}{(1 + v_3)^2} \quad (\text{common denominator})$$

$$= \frac{2(1 + v_3)}{(1 + v_3)^2} \quad (\vec{v} \text{ is a unit vector})$$

$$= \frac{2}{1 + v_3} \quad (\text{cancel})$$

So $|\lambda_1\rangle = \sqrt{\frac{1+v_3}{2}} \begin{bmatrix} 1 \\ \frac{1-v_3}{v_1-iv_2} \end{bmatrix}$ is a normalized eigenvector with eigenvalue 1. Similarly, $|\lambda_{-1}\rangle = \sqrt{\frac{1-v_3}{2}} \begin{bmatrix} 1 \\ -\frac{1+v_3}{v_1-iv_2} \end{bmatrix}$ is a normalized eigenvector with eigenvalue -1 . For convenience, we write $|\lambda_{\pm 1}\rangle = \sqrt{\frac{1\pm v_3}{2}} \begin{bmatrix} 1 \\ \frac{v_3\mp 1}{v_1-iv_2} \end{bmatrix}$. Calculating the projectors:

$$\begin{aligned}
|\lambda_{\pm 1}\rangle\langle\lambda_{\pm 1}| &= \frac{1\pm v_3}{2} \begin{bmatrix} 1 \\ \frac{v_3\mp 1}{v_1-iv_2} \end{bmatrix} \begin{bmatrix} 1 & \left(\frac{v_3\mp 1}{v_1-iv_2}\right)^* \end{bmatrix} && \text{(definition)} \\
&= \frac{1\pm v_3}{2} \begin{bmatrix} 1 \\ \frac{v_1+iv_2}{v_3\pm 1} \end{bmatrix} \begin{bmatrix} 1 & \left(\frac{v_1+iv_2}{v_3\pm 1}\right)^* \end{bmatrix} && \text{(flip)} \\
&= \frac{1\pm v_3}{2} \begin{bmatrix} 1 \\ \frac{v_1+iv_2}{v_3\pm 1} \end{bmatrix} \begin{bmatrix} 1 & \frac{v_1-iv_2}{v_3\pm 1} \end{bmatrix} && (v_1, v_2 \text{ are real, conjugate}) \\
&= \frac{1\pm v_3}{2} \begin{bmatrix} 1 & \frac{v_1-iv_2}{1\pm v_3} \\ \frac{v_1+iv_2}{1\pm v_3} & \frac{v_1^2+v_2^2}{(1\pm v_3)^2} \end{bmatrix} && \text{(multiply)} \\
&= \frac{1}{2} \begin{bmatrix} 1\pm v_3 & v_1-iv_2 \\ v_1+iv_2 & \frac{1-v_3^2}{1\pm v_3} \end{bmatrix} && (v \text{ is a unit vector, cancel external numerator}) \\
&= \frac{1}{2} \begin{bmatrix} 1\pm v_3 & v_1-iv_2 \\ v_1+iv_2 & 1\mp v_3^2 \end{bmatrix} && \text{(cancel internal denominator)} \\
&= \frac{1}{2} \left(I \pm \begin{bmatrix} v_3 & v_1-iv_2 \\ v_1+iv_2 & -v_3 \end{bmatrix} \right) && \text{(separate)} \\
&= \frac{1}{2} (I \pm \vec{v} \cdot \vec{\sigma}) && \text{(definition)}
\end{aligned}$$

The first author points out that when $v_1 - iv_2 = 0$, or equivalently when $v_3 \pm 1 = 0$, various expressions above are indeterminant. The first author attempts to circumvent this by working more generally below, however, the second author is suspicious that the argument is circular. Some of the details are instructive however, so it is left below. To deal with the degenerate cases, note that, since \vec{v} is a real-valued vector, $v_1 - iv_2 = 0$ implies $v_1 = 0, v_2 = 0$, and $v_3 = \pm 1$, in which case $\vec{v} \cdot \vec{\sigma} = \pm Z$. The normalized eigenvectors are $|\lambda_1\rangle = \pm|0\rangle$ and $|\lambda_{-1}\rangle = \pm|1\rangle$, where here the \pm indicate the sign of v_3 instead of the sign of the eigenvalue. Now

$$\begin{aligned}
|\lambda_1\rangle\langle\lambda_1| &= (\pm 1)^2 |0\rangle\langle 0| = \begin{bmatrix} 1 & 0 \\ 0 & 0 \end{bmatrix} = \frac{1}{2}(I + Z) = \frac{1}{2}(I + \vec{v} \cdot \vec{\sigma}) \\
|\lambda_{-1}\rangle\langle\lambda_{-1}| &= (\pm 1)^2 |1\rangle\langle 1| = \begin{bmatrix} 0 & 0 \\ 0 & 1 \end{bmatrix} = \frac{1}{2}(I - Z) = \frac{1}{2}(I - \vec{v} \cdot \vec{\sigma}),
\end{aligned}$$

completing the proof. The first author's attempted resolution follows:

While I review my proof, I notice that my proof has a defect. The case $(v_1, v_2, v_3) = (0, 0, 1)$, second component of eigenstate, $\frac{1-v_3}{v_1-iv_2}$, diverges. So I implicitly assume $v_1 - iv_2 \neq 0$. Hence my proof is incomplete.

Since the exercise doesn't require explicit form of projector, we should prove the problem more abstractly. In order to prove, we use the following properties of $\vec{v} \cdot \vec{\sigma}$

- $\vec{v} \cdot \vec{\sigma}$ is Hermitian
- $(\vec{v} \cdot \vec{\sigma})^2 = I$ where \vec{v} is a real unit vector.

We can easily check above conditions.

$$\begin{aligned}
 (\vec{v} \cdot \vec{\sigma})^\dagger &= (v_1\sigma_1 + v_2\sigma_2 + v_3\sigma_3)^\dagger \\
 &= v_1\sigma_1^\dagger + v_2\sigma_2^\dagger + v_3\sigma_3^\dagger \\
 &= v_1\sigma_1 + v_2\sigma_2 + v_3\sigma_3 \quad (\because \text{Pauli matrices are Hermitian.}) \\
 &= \vec{v} \cdot \vec{\sigma}
 \end{aligned}$$

$$\begin{aligned}
 (\vec{v} \cdot \vec{\sigma})^2 &= \sum_{j,k=1}^3 (v_j\sigma_j)(v_k\sigma_k) \\
 &= \sum_{j,k=1}^3 v_jv_k\sigma_j\sigma_k \\
 &= \sum_{j,k=1}^3 v_jv_k \left(\delta_{jk}I + i \sum_{l=1}^3 \epsilon_{jkl}\sigma_l \right) \quad (\because \text{Exercise 2.43, eqn (2.78) page 78}) \\
 &= \sum_{j,k=1}^3 v_jv_k\delta_{jk}I + i \sum_{j,k,l=1}^3 \epsilon_{jkl}v_jv_k\sigma_l \\
 &= \sum_{j=1}^3 v_j^2 I \\
 &= I \quad \left(\because \sum_j v_j^2 = 1 \right)
 \end{aligned}$$

Proof. Suppose $|\lambda\rangle$ is an eigenstate of $\vec{v} \cdot \vec{\sigma}$ with eigenvalue λ . Then

$$\begin{aligned}
 \vec{v} \cdot \vec{\sigma} |\lambda\rangle &= \lambda |\lambda\rangle \\
 (\vec{v} \cdot \vec{\sigma})^2 |\lambda\rangle &= \lambda^2 |\lambda\rangle
 \end{aligned}$$

On the other hand $(\vec{v} \cdot \vec{\sigma})^2 = I$,

$$\begin{aligned}
 (\vec{v} \cdot \vec{\sigma})^2 |\lambda\rangle &= I |\lambda\rangle = |\lambda\rangle \\
 \therefore \lambda^2 |\lambda\rangle &= |\lambda\rangle.
 \end{aligned}$$

Thus $\lambda^2 = 1 \Rightarrow \lambda = \pm 1$. Therefore $\vec{v} \cdot \vec{\sigma}$ has eigenvalues ± 1 .

Let $|\lambda_1\rangle$ and $|\lambda_{-1}\rangle$ are eigenvectors with eigenvalues 1 and -1 , respectively. I will prove that $P_\pm = |\lambda_{\pm 1}\rangle\langle\lambda_{\pm 1}|$.

In order to prove above equation, all we have to do is prove following condition. (see Theorem 2.0.1)

$$\langle\psi|(P_\pm - |\lambda_{\pm 1}\rangle\langle\lambda_{\pm 1}|)|\psi\rangle = 0 \text{ for all } |\psi\rangle \in \mathbb{C}^2. \quad (2.1)$$

Since $\vec{v} \cdot \vec{\sigma}$ is Hermitian, $|\lambda_1\rangle$ and $|\lambda_{-1}\rangle$ are orthonormal vector (\because Exercise 2.22). Let $|\psi\rangle \in \mathbb{C}^2$ be an arbitrary state. $|\psi\rangle$ can be written as

$$|\psi\rangle = \alpha |\lambda_1\rangle + \beta |\lambda_{-1}\rangle \quad (|\alpha|^2 + |\beta|^2 = 1, \alpha, \beta \in \mathbb{C}).$$

$$\begin{aligned}
\langle \psi | (P_{\pm} - |\lambda_{\pm}\rangle\langle\lambda_{\pm}|) | \psi \rangle &= \langle \psi | P_{\pm} | \psi \rangle - \langle \psi | \lambda_{\pm} \rangle \langle \lambda_{\pm} | \psi \rangle. \\
\langle \psi | P_{\pm} | \psi \rangle &= \langle \psi | \frac{1}{2} (I \pm \vec{v} \cdot \vec{\sigma}) | \psi \rangle \quad \left(\begin{array}{l} \text{implicit assumption the above equals 0} \\ \text{and that } |\lambda_{\pm 1}\rangle\langle\lambda_{\pm 1}| = \frac{1}{2}(I \pm \vec{v} \cdot \vec{\sigma}) \end{array} \right) \\
&= \frac{1}{2} \pm \frac{1}{2} \langle \psi | \vec{v} \cdot \vec{\sigma} | \psi \rangle \\
&= \frac{1}{2} \pm \frac{1}{2} (|\alpha|^2 - |\beta|^2) \\
&= \frac{1}{2} \pm \frac{1}{2} (2|\alpha|^2 - 1) \quad (\because |\alpha|^2 + |\beta|^2 = 1) \\
\langle \psi | \lambda_1 \rangle \langle \lambda_1 | \psi \rangle &= |\alpha|^2 \\
\langle \psi | \lambda_{-1} \rangle \langle \lambda_{-1} | \psi \rangle &= |\beta|^2 = 1 - |\alpha|^2
\end{aligned}$$

Therefore $\langle \psi | (P_{\pm} - |\lambda_{\pm 1}\rangle\langle\lambda_{\pm 1}|) | \psi \rangle = 0$ for all $|\psi\rangle \in \mathbb{C}^2$. Thus $P_{\pm} = |\lambda_{\pm 1}\rangle\langle\lambda_{\pm 1}|$. \square

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

Soln:

$$\begin{aligned}
p(1) = \langle 0 | P_1 | 0 \rangle &= \left\langle 0 \left| \frac{1}{2} (I + \vec{v} \cdot \vec{\sigma}) \right| 0 \right\rangle && \text{(definition)} \\
&= \langle 0 | \lambda_1 \rangle \langle \lambda_1 | 0 \rangle && \text{(use eigenvector directly)} \\
&= \frac{1+v_3}{2} \begin{bmatrix} 1 & 0 \end{bmatrix} \begin{bmatrix} 1 \\ \frac{1-v_3}{v_1-iv_2} \end{bmatrix} \begin{bmatrix} 1 \\ 0 \end{bmatrix} && \text{(substitute)} \\
&= \frac{1+v_3}{2} \begin{bmatrix} 1 & 0 \end{bmatrix} \begin{bmatrix} 1 & \frac{v_1-iv_2}{1+v_3} \\ \frac{v_1+iv_2}{1+v_3} & \frac{1-v_3}{1+v_3} \end{bmatrix} \begin{bmatrix} 1 \\ 0 \end{bmatrix} && \text{(multiply)} \\
&= \frac{1}{2} (1+v_3) && \text{(extract 0,0 entry)}
\end{aligned}$$

The post-measurement state is

$$\begin{aligned}
\frac{|\lambda_1\rangle \langle \lambda_1 | 0 \rangle}{\sqrt{p(1)}} &= \left(\frac{\sqrt{\frac{1+v_3}{2}}}{\sqrt{\frac{1+v_3}{2}}} \begin{bmatrix} 1 & \left(\frac{1-v_3}{v_1-iv_2} \right)^* \end{bmatrix} \begin{bmatrix} 1 \\ 0 \end{bmatrix} \right) |\lambda_1\rangle && \left(\begin{array}{l} \text{numerator from normalization} \\ \text{denominator from above} \end{array} \right) \\
&= |\lambda_1\rangle && \text{(simplify, multiply)}
\end{aligned}$$

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

Soln: We can no longer avoid the converse of Exercise 2.16. Before we prove it, we quibble about semantics. The assumption in the exercise is that the POVM elements “coincide” with the measurement operators. This can be taken to mean one of two things. Let $\{M_m\}$ be the set of measurement operators. One interpretation is that $M_m = E_m = M_m^\dagger M_m$ for all m . We’ll call this direct coincidence. Another interpretation is that they coincide as sets, i.e., that $\{M_m\} = \{E_m\}$, without requiring $M_m = E_m$ for each m . A mathematician may argue that what is to be assumed is set-wise coincidence, but we’ll show that this is not the case. We’ll argue for the assumption of direct coincidence by contradiction. If it were the case that POVM measurements which satisfied the setwise coincidence assumption but not the direct coincidence assumption were projective measurements, then there would exist $M_m = E_\mu = \overline{M}_\mu^\dagger M_\mu$, where $M_m \neq M_\mu$. However, being a projective measurement, M_μ must be Hermitian, in which case $M_m = M_\mu^2 = M_\mu$, by Exercise 2.16 (not the converse). This is a contradiction, but it is important to note

what we can conclude from it. Our assumption was that setwise coincidence of measurement operators and POVM measurements was enough to prove the measurement projective. Having contradicted this, we conclude that the exercise must assume direct coincidence instead. We have not proven that setwise coincidence without direct coincidence is impossible.

Theorem. 2.0.2. *Let P be a normal linear operator. If $P^2 = P$, then P is a projector, that is $P = \sum_i |i\rangle\langle i|$ for some orthonormal basis $|i\rangle$.*

Proof. Having assumed normality (the statement is not true in general if we do not), we may apply the Structural Decomposition Theorem to write $P = \sum_i \lambda_i |i\rangle\langle i|$. Note, we may assume $\lambda_i \neq 0$. Being idempotent, $\sum_i \lambda_i |i\rangle\langle i| = P = P^2 = (\sum_i \lambda_i |i\rangle\langle i|)^2 = \sum_i \lambda_i^2 |i\rangle\langle i|$ by orthonormality. The $|i\rangle\langle i|$ are linearly independent (see Exercise 2.10), from which we conclude that all λ_i satisfy $\lambda_i = \lambda_i^2$, from which $\lambda_i(\lambda_i - 1) = 0$, or that $\lambda_i = 1$, since $\lambda_i \neq 0$. Now $P = \sum_i |i\rangle\langle i|$, as required.

Note: Some definitions of projectors define them in terms of relations between their kernels and images as opposed to the formulaic definition given in Equation 2.35. They are equivalent (at least for normal matrices). With the set theoretic definition, Exercise 2.16 and it's converse Theorem 2.0.2 follow by definition. \square

On to the exercise: Suppose M_m is a measurement operator such that $E_m = M_m^\dagger M_m = M_m$. Note that $M_m = M_m^\dagger M_m$ is positive by exercise 2.25, thus Hermitian, so $M_m = E_m = M_m^\dagger M_m = M_m^2$. Being Hermitian, M_m is normal, in which case Theorem 2.0.2 above applies, giving that M_m is a projector. Thus the measurement is a projective measurement.

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.

Soln: By the singular value decomposition (Corrolary 2.4, eq 2.80, p 79), there exists unitary U, V , and real-valued diagonal D such that $M_m = UDV$. Now

$$\begin{aligned}
 \sqrt{E_m} &= \sqrt{M_m^\dagger M_m} && \text{(definition)} \\
 &= \sqrt{V^\dagger D^\dagger U^\dagger U D V} && \text{(properties of }^\dagger) \\
 &= \sqrt{V^\dagger D^2 V} && (U \text{ unitary, } D \text{ real-valued diagonal)} \\
 &= V^\dagger D V && (V \text{ unitary} \rightarrow (V^\dagger D V)^2 = V^\dagger D^2 V) \\
 &= V^\dagger (U^\dagger U) D V && (U \text{ unitary, } U^\dagger U = I) \\
 UV \sqrt{E_m} &= M_m && (U, V \text{ unitary, solve for } M_m)
 \end{aligned}$$

Define $U_m \equiv UV$ so that $M_m = U_m \sqrt{E_m}$, completing the solution.

2.64) Suppose Bob is given a quantum state chosen from a set $|\psi_1\rangle, \dots, |\psi_m\rangle$ of linearly independent states. Construct a POVM $\{E_1, \dots, E_{m+1}\}$ such that if outcome E_i occurs, $1 \leq i \leq m$, then Bob knows with certainty that he was given the state $|\psi_i\rangle$.

Soln: Being linearly independent, S forms a basis for the subspace it spans. Consider the subspaces spanned by $S'_i = \{|\psi_j\rangle\}_{j \neq i}$. Let $|\psi'_i\rangle$ be a non-zero unit vector in the orthogonal complement. Note that if $i \neq j$, then $\langle \psi_i | \psi'_j \rangle = 0$, since ψ'_j is in the orthogonal complement of a subspace containing ψ_i . Also note that $\langle \psi_i | \psi'_i \rangle \neq 0$, since if it were, then ψ_i would be in the subspace spanned by S'_i , violating linear independence. Combining, we write $\langle \psi_i | \psi'_j \rangle = \delta_{i,j} \cdot p_i$, for some non-zero p_i . We may scale the ψ'_i so that $p_i = 1$. For, $1 \leq i \leq m$, define $E_i = |\psi'_i\rangle\langle \psi'_i|$, then, to cover the define $E_{m+1} = I - \sum_m E_i$. Now, for

$1 \leq i, j \leq m$,

$$\begin{aligned}
 p_i(j) &\equiv \langle \psi_i | E_j^\dagger E_j | \psi_i \rangle && \text{(definition)} \\
 &= \langle \psi_i | (|\psi'_j\rangle\langle\psi'_j|)^\dagger |\psi'_j\rangle\langle\psi'_j| | \psi_i \rangle && \text{(definition of } E_j) \\
 &= \langle \psi_i | \psi'_j \rangle \langle \psi'_j | \psi'_j \rangle \langle \psi'_j | \psi_i \rangle && \text{(Exercise 2.13: } (|\phi\rangle\langle\phi|)^\dagger = |\phi\rangle\langle\phi|) \\
 &= \delta_{i,j} \cdot 1 \cdot \delta_{i,j} && \text{(construction)} \\
 &= \delta_{i,j} && \text{(simplify)}
 \end{aligned}$$

That is, E_1, \dots, E_{m+1} is a POVM such that state $|\psi_i\rangle$ and outcome E_i correspond exactly, for $1 \leq i \leq m$. Outcome E_{m+1} will result from measuring any state outside of the span of S , that is, in their orthogonal complement, with probability 1 as well. The first author included links to several relevant journal articles. The second author has not evaluated them, but they are include below.

- Lu-Ming Duan, Guang-Can Guo. Probabilistic cloning and identification of linearly independent quantum states. Phys. Rev. Lett.,80:4999-5002, 1998. arXiv:quant-ph/9804064
<https://journals.aps.org/prl/abstract/10.1103/PhysRevLett.80.4999>
<https://arxiv.org/abs/quant-ph/9804064>
- Stephen M. Barnett, Sarah Croke, Quantum state discrimination, arXiv:0810.1970 [quant-ph]
<https://arxiv.org/abs/0810.1970>
https://www.osapublishing.org/DirectPDFAccess/67EF4200-CBD2-8E68-1979E37886263936_176580/aop-1-2-238.pdf

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

Soln: Note that $(|0\rangle + |1\rangle)/\sqrt{2} = |+\rangle$ and $(|0\rangle - |1\rangle)/\sqrt{2} = |-\rangle$ are an orthonormal basis. Since the amplitude of $|+\rangle$ in the expression $(|0\rangle + |1\rangle)/\sqrt{2} = |+\rangle$ is 1, and that in $(|0\rangle - |1\rangle)/\sqrt{2} = |-\rangle$ is 0, there is no relative phase θ such that $e^{i\theta} \cdot 1 = 0$, so $|+\rangle$ and $|-\rangle$ do not differ by a relative phase in the basis they comprise.

2.66) Show that the average value of the observable $X_1 Z_2$ for a two qubit system measured in the state $(|00\rangle + |11\rangle)/\sqrt{2}$ is zero.

Soln: Let $|\Phi^+\rangle = (|00\rangle + |11\rangle)/\sqrt{2}$.

$$\begin{aligned}
 \mathbf{E}(X_1 Z_2) &= \langle X_1 Z_2 \rangle = \langle \Phi^+ | X_1 Z_2 | \Phi^+ \rangle \\
 &= \frac{\langle 00 | + \langle 11 |}{\sqrt{2}} \cdot \frac{X_1 Z_2 (|00\rangle + |11\rangle)}{\sqrt{2}} \\
 &= \frac{\langle 00 | + \langle 11 |}{\sqrt{2}} \cdot \frac{|10\rangle - |01\rangle}{\sqrt{2}} \\
 &= \frac{\langle 00 | 10 \rangle - \langle 00 | 01 \rangle + \langle 11 | 10 \rangle - \langle 11 | 01 \rangle}{2} \\
 &= \frac{0 - 0 + 0 - 0}{2} \\
 &= 0
 \end{aligned}$$

2.67) Suppose V is a Hilbert space with a subspace W . Suppose $U : W \rightarrow V$ 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 $U' : V \rightarrow V$ 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.

Soln: Let $|w_i\rangle$ be an orthonormal basis for W . Since U preserves inner products in W , $|u_i\rangle \equiv U|w_i\rangle$ is an orthonormal basis for $\text{image}(U)$, hence $\langle u_i|u_j\rangle = \delta_{i,j}$. Consider the orthogonal complement, W^\perp . By definition $V = W \oplus W^\perp$. Let $|w'_j\rangle$ and $|u'_j\rangle$ be orthonormal bases for W^\perp and $(\text{image}(U))^\perp$. Note that, as provided, U is not defined on the $|w'_j\rangle$, so we may not state that $|u'_j\rangle = U|w'_j\rangle$, where here we strike the statement not because it is necessarily false, but because it is not assumed. We can, however, state that $\langle u'_i|u'_j\rangle = \delta_{i,j}$, since the $|u_j\rangle$ are orthonormal, and that $\langle u_i|u'_j\rangle = 0$, since the $|u'_j\rangle$ are in a space orthogonal to the $|u_i\rangle$. Define $U' : V \rightarrow V$ as $U' = \sum_i |u_i\rangle\langle w_i| + \sum_j |u'_j\rangle\langle w'_j|$. First, we prove unitarity:

$$\begin{aligned}
 (U')^\dagger U' &= \left(\sum_i |w_i\rangle\langle u_i| + \sum_j |w'_j\rangle\langle u'_j| \right) \left(\sum_k |u_k\rangle\langle w_k| + \sum_\ell |u'_\ell\rangle\langle w'_\ell| \right) \quad (\text{definition, linearity, Exercise 2.13}) \\
 &= \sum_{i,k} |w_i\rangle\langle u_i|u_k\rangle\langle w_k| + \sum_{i,\ell} |w_i\rangle\langle u_i|u'_\ell\rangle\langle w'_\ell| + \sum_{j,k} |w'_j\rangle\langle u'_j|u_k\rangle\langle w_k| + \sum_{j,\ell} |w'_j\rangle\langle u'_j|u'_\ell\rangle\langle w'_\ell| \\
 &\quad (\text{F.O.I.L., linearity}) \\
 &= \sum_{i,k} \delta_{i,k} |w_i\rangle\langle w_k| + \sum_{j,\ell} \delta_{j,\ell} |w'_j\rangle\langle w'_\ell| \quad (\text{orthonormality, orthogonality}) \\
 &= \sum_i |w_i\rangle\langle w_i| + \sum_j |w'_j\rangle\langle w'_j| \quad (\text{collect non-zero term}) \\
 &= I \quad (\text{completeness})
 \end{aligned}$$

where the last equality holds because $\{|w_i\rangle\} \cup \{|w'_j\rangle\}$ forms a basis for the entire space V . Similarly $U'(U')^\dagger = I$, so U' is unitary. It is left only to show that U' is an extension of U , that is $U'|w\rangle = U|w\rangle$ for all w in W .

$$\begin{aligned}
 U'|w\rangle &= \left(\sum_i |u_i\rangle\langle w_i| + \sum_j |u'_j\rangle\langle w'_j| \right) |w\rangle \quad (\text{definition of } U') \\
 &= \left(\sum_i |u_i\rangle\langle w_i| \right) |w\rangle + \sum_j |u'_j\rangle\langle w'_j|w\rangle \quad (\text{linearity}) \\
 &= \left(\sum_i |u_i\rangle\langle w_i| \right) |w\rangle \quad (|w\rangle \in W, |w'_j\rangle \in W^\perp) \\
 &= \left(\sum_i U|w_i\rangle\langle w_i| \right) |w\rangle \quad (\text{definition of } u_i) \\
 &= U \left(\sum_i |w_i\rangle\langle w_i| \right) |w\rangle \quad (\text{linearity}) \\
 &= U|w\rangle. \quad (\text{completeness})
 \end{aligned}$$

where the last inequality holds only when applied to vectors in the subspace W , as is being done.

2.68) Prove that $(|\Psi^+\rangle) = |\psi\rangle \equiv (|00\rangle + |11\rangle)/\sqrt{2} \neq |a\rangle|b\rangle$ for all single qubit states $|a\rangle$ and $|b\rangle$.

Soln: Suppose $|a\rangle = a_0|0\rangle + a_1|1\rangle$ and $|b\rangle = b_0|0\rangle + b_1|1\rangle$. Then

$$|a\rangle|b\rangle = |a\rangle \otimes |b\rangle = a_0b_0|00\rangle + a_0b_1|01\rangle + a_1b_0|10\rangle + a_1b_1|11\rangle.$$

If $|\psi\rangle = |a\rangle|b\rangle$, then $a_0b_0 = 1$, $a_0b_1 = 0$, $a_1b_0 = 0$, $a_1b_1 = 1$ since $\{|ij\rangle\}$ is an orthonormal basis. Since $a_0b_1 = 0$, either $a_0 = 0$ or $b_1 = 0$, however, $a_0 = 0$ contradicts $a_0b_0 = 1$, and $b_1 = 0$ contradicts $a_1b_1 = 1$. Thus $|\psi\rangle \neq |a\rangle|b\rangle$.

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

states and the Bell Matrix as follows.

$$\begin{aligned} |\Phi^+\rangle = |\psi_0\rangle &\equiv \frac{|00\rangle + |11\rangle}{\sqrt{2}} = \frac{1}{\sqrt{2}} \begin{bmatrix} 1 \\ 0 \\ 0 \\ 1 \end{bmatrix} & |\Phi^-\rangle = |\psi_1\rangle &\equiv \frac{|00\rangle - |11\rangle}{\sqrt{2}} = \frac{1}{\sqrt{2}} \begin{bmatrix} 1 \\ 0 \\ 0 \\ -1 \end{bmatrix} \\ |\Psi^+\rangle = |\psi_2\rangle &\equiv \frac{|01\rangle + |10\rangle}{\sqrt{2}} = \frac{1}{\sqrt{2}} \begin{bmatrix} 0 \\ 1 \\ 1 \\ 0 \end{bmatrix} & |\Psi^-\rangle = |\psi_3\rangle &\equiv \frac{|01\rangle - |10\rangle}{\sqrt{2}} = \frac{1}{\sqrt{2}} \begin{bmatrix} 0 \\ 1 \\ -1 \\ 0 \end{bmatrix} \end{aligned}$$

$$B = \begin{bmatrix} |\Phi^+\rangle & |\Phi^-\rangle & |\Psi^+\rangle & |\Psi^-\rangle \end{bmatrix}$$

Note that $\langle\psi_i|\psi_j\rangle$ is the i, j -entry in $B^\dagger B$, so orthonormality will follow if $B^\dagger B = I$.

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

Linear independence follows from orthogonality. Being 4 linearly independent vectors in a 4-dimensional vector space, the Bell states form a basis.

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?

Soln: For any of the Bell states we get $\langle\psi_i|E \otimes I|\psi_i\rangle = \frac{1}{2}(\langle 0|E|0\rangle + \langle 1|E|1\rangle)$. We exhibit the calculation for $|\psi_0\rangle$. Others are similar.

$$\begin{aligned} \langle\psi_0|E \otimes I|\psi_0\rangle &= \frac{1}{2} \left((\langle 00| + \langle 11|) \cdot (E \otimes I)(|00\rangle + |11\rangle) \right) && \text{(definition)} \\ &= \frac{1}{2} \left((\langle 00| + \langle 11|) \cdot ((E|0\rangle) \otimes |0\rangle + (E|1\rangle) \otimes |1\rangle) \right) && (E \text{ and } I \text{ act independently}) \\ &= \frac{1}{2} \left(\langle 0|E|0\rangle \cdot \langle 0|0\rangle + \langle 0|E|1\rangle \cdot \langle 0|1\rangle + \langle 1|E|0\rangle \cdot \langle 1|0\rangle + \langle 1|E|1\rangle \cdot \langle 1|1\rangle \right) && \text{(F.O.I.L.)} \\ &= \frac{1}{2} (\langle 0|E|0\rangle + \langle 1|E|1\rangle) && \text{(collect non-zero terms)} \end{aligned}$$

Suppose Eve measures the qubit Alice sent by measurement operators M_m . The probability that Eve gets result m is $p_i(m) = \langle\psi_i|M_m^\dagger M_m \otimes I|\psi_i\rangle$. Since $M_m^\dagger M_m$ is positive, the result above applies, in which case the $p_i(m)$ take on the same values for all $|\psi_i\rangle$, that is, for all i and m , $p_i(m) = 1/4$. In other words, no matter the outcome m , the probability that ψ was in any of the Bell states is uniform. So Eve can't distinguish Bell states given only access to a single qubit. [Note, the exercise above only proves Eve can't distinguish Bell states given only access to the first qubit, but the only difference from Bob's perspective is that Eve's Bell basis uses $-|\psi_3\rangle = -|\Psi^-\rangle$. This negative does not change the result of any of the calculations, so Eve can't distinguish the Bell states given access to Bob's qubit either.]

2.71) (Criterion to decide if a state is mixed or pure) Let ρ be a density operator. Show that $\text{tr}(\rho^2) \leq 1$, with equality if and only if ρ is a pure state.

Soln: By definition, there exists an ensemble of pure states $\{p_i, |\psi_i\rangle\}$, where $0 \leq p_i \leq 1$ and $\sum_i p_i = 1$ such that $\rho = \sum_i p_i |\psi_i\rangle\langle\psi_i|$. Note that on this range $0 \leq p_i^2 \leq p_i$. Also, $\rho = \sum_i \lambda_i |i\rangle\langle i|$ for some orthonormal

basis $|i\rangle\langle i|$, by the spectral decomposition theorem. Express $|\psi_i\rangle = \sum_j \alpha_j |i\rangle \dots$ I believe the first author's proof below assumes that $|\psi_i\rangle = |i\rangle$, i.e. that the $|\psi_i\rangle$ are orthonormal, but this is not guaranteed, and in fact, very much not assumed.

$$\begin{aligned}\rho^2 &= \sum_{i,j} p_i p_j |i\rangle\langle i|j\rangle\langle j| \\ &= \sum_{i,j} p_i p_j |i\rangle\langle j| \delta_{ij} \\ &= \sum_i p_i^2 |i\rangle\langle i|\end{aligned}$$

$$\text{tr}(\rho^2) = \text{tr}\left(\sum_i p_i^2 |i\rangle\langle i|\right) = \sum_i p_i^2 \text{tr}(|i\rangle\langle i|) = \sum_i p_i^2 \langle i|i\rangle = \sum_i p_i^2 \leq \sum_i p_i = 1 \quad (\because p_i^2 \leq p_i)$$

Suppose $\text{tr}(\rho^2) = 1$. Then $\sum_i p_i^2 = 1$. Since $p_i^2 < p_i$ for $0 < p_i < 1$, only single p_i should be 1 and otherwise have to vanish. Therefore $\rho = |\psi_i\rangle\langle\psi_i|$. It is a pure state.

Conversely if ρ is pure, then $\rho = |\psi\rangle\langle\psi|$.

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

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

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

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

(2) What is the Bloch vector representation for the state $\rho = I/2$?

(3) Show that a state ρ is pure if and only if $\|\vec{r}\| = 1$.

(4) Show that for pure states the description of the Bloch vector we have given coincides with that in Section 1.2.

Soln: Note, even though the topic of this problem includes mixed states, the representation we are constructing is a representation of a single qubit. That qubit could be entangled with others, but these other qubits are not explicitly represented in the Bloch sphere representation of the qubit of interest. The level of entanglement of the qubit of interest with other qubits is represented, though.

(1) Let ρ be an arbitrary density matrix for a single complex-dimensional state. ρ is Hermitian, so we

may set $\rho = \begin{bmatrix} a & b \\ b^* & d \end{bmatrix}$, where $a, d \in \mathbb{R}$ and $b \in \mathbb{C}$. Because ρ is density matrix, $\text{tr}(\rho) = a + d = 1$. Define $r_1 = \Re(b)/2$, $r_2 = -\Im(b)/2$, $r_3 = a - d$, and finally $\vec{r} = (r_1, r_2, r_3)$. Expressing a, b , and d in terms of \vec{r} , we have $a = \frac{1+r_3}{2}$, $b = \frac{r_1 - ir_2}{2}$, and $d = \frac{1-r_3}{2}$. Now

$$\rho = \begin{bmatrix} a & b \\ b^* & d \end{bmatrix} = \frac{1}{2} \begin{bmatrix} 1+r_3 & r_1 - ir_2 \\ r_1 + ir_2 & 1-r_3 \end{bmatrix} = \frac{1}{2}(I + \vec{r} \cdot \vec{\sigma}).$$

Thus an arbitrary density matrix ρ can be written as $\rho = \frac{1}{2}(I + \vec{r} \cdot \vec{\sigma})$ for some real-valued three-dimensional vector \vec{r} . It remains to show that $\|\vec{r}\| \leq 1$. To do so, note that Theorem 2.5 gives that ρ is positive. Being

positive, the eigenvalues of ρ must be non-negative. So, let's find the eigenvalues:

$$\begin{aligned}
 \det(\rho - \lambda I) &= (a - \lambda)(d - \lambda) - \|b\|^2 && \text{(characteristic equation)} \\
 &= \lambda^2 - (a + d)\lambda + ad - \|b\|^2 = 0 && \text{(simplify)} \\
 &= \lambda^2 - \lambda + \left(\frac{1 - r_3^2}{4} - \frac{r_1^2 + r_2^2}{4} \right) && (\text{tr}(\rho) = 1, \text{ express in terms of } \vec{r}) \\
 &= \lambda^2 - \lambda + \left(\frac{1 - \|\vec{r}\|^2}{4} \right) && \text{(definition of } \|\cdot\| \text{ in } \mathbb{R}^3) \\
 \lambda &= \frac{1 \pm \sqrt{1 - (1 - \|\vec{r}\|^2)}}{2} && \text{(quadratic formula)} \\
 &= \frac{1 \pm \|\vec{r}\|^2}{2} && \text{(simplify)}
 \end{aligned}$$

Now $\frac{1 - \|\vec{r}\|^2}{2} \geq 0 \rightarrow \|\vec{r}\| \leq 1$.

(2) If $\rho = I/2$, then $a = d = 1/2$ and $b = 0$. So $\vec{v} = (0, 0, 0)$, and $\rho = I/2$ corresponds to the origin of Bloch sphere.

(3) By exercise 2.71, ρ is a pure state if and only if $\text{tr}(\rho^2) = 1$. Let's calculate $\text{tr}(\rho^2)$.

$$\begin{aligned}
 \rho^2 &= \frac{1}{2}(I + \vec{r} \cdot \vec{\sigma}) \frac{1}{2}(I + \vec{r} \cdot \vec{\sigma}) && \text{(part (1) above)} \\
 &= \frac{1}{4} \left[I + 2\vec{r} \cdot \vec{\sigma} + \|\vec{r}\|^2 \left(\frac{\vec{r} \cdot \vec{\sigma}}{\|\vec{r}\|} \right)^2 \right] && \text{(F.O.I.L., normalize)} \\
 &= \frac{1}{4} (I + 2\vec{r} \cdot \vec{\sigma} + \|\vec{r}\|^2 I) && \left(\begin{array}{l} \text{see the first author's attempted resolution} \\ \text{of the special case of Exercise 2.60} \end{array} \right) \\
 \text{tr}(\rho^2) &= \frac{1}{4} [\text{tr}(I) + 2\text{tr}(\vec{r} \cdot \vec{\sigma}) + \|\vec{r}\|^2 \text{tr}(I)] && \text{(linearity of tr)} \\
 &= \frac{2 + 2\|\vec{r}\|^2}{4}. && (\text{tr}(I) = 2 \text{ in } \mathbb{C}^2, \text{tr}(\vec{r} \cdot \vec{\sigma}) = r_3 - r_3 = 0)
 \end{aligned}$$

Now ρ is a pure state if and only if $\text{tr}(\rho^2) = \frac{2 + 2\|\vec{r}\|^2}{4} = 1$, which occurs if and only if $\|\vec{r}\| = 1$.

(4) TODO

2.73)

Theorem 2.6)

$$\rho = \sum_i p_i |\psi_i\rangle\langle\psi_i| = \sum_i |\tilde{\psi}_i\rangle\langle\tilde{\psi}_i| = \sum_j |\tilde{\varphi}_j\rangle\langle\tilde{\varphi}_j| = \sum_j q_j |\varphi_j\rangle\langle\varphi_j| \Leftrightarrow |\tilde{\psi}_i\rangle = \sum_j u_{ij} |\tilde{\varphi}_j\rangle$$

where u is unitary.

The-transformation in theorem 2.6, $|\tilde{\psi}_i\rangle = \sum_j u_{ij} |\tilde{\varphi}_j\rangle$, corresponds to

$$\left[|\tilde{\psi}_1\rangle \cdots |\tilde{\psi}_k\rangle \right] = \left[|\tilde{\varphi}_1\rangle \cdots |\tilde{\varphi}_k\rangle \right] U^T$$

where $k = \text{rank}(\rho)$.

$$\sum_i |\tilde{\psi}_i\rangle\langle\tilde{\psi}_i| = \left[|\tilde{\psi}_1\rangle \cdots |\tilde{\psi}_k\rangle \right] \begin{bmatrix} \langle\tilde{\psi}_1| \\ \vdots \\ \langle\tilde{\psi}_k| \end{bmatrix} \quad (2.2)$$

$$= \left[|\tilde{\varphi}_1\rangle \cdots |\tilde{\varphi}_k\rangle \right] U^T U^* \begin{bmatrix} \langle\tilde{\varphi}_1| \\ \vdots \\ \langle\tilde{\varphi}_k| \end{bmatrix} \quad (2.3)$$

$$= \left[|\tilde{\varphi}_1\rangle \cdots |\tilde{\varphi}_k\rangle \right] \begin{bmatrix} \langle\tilde{\varphi}_1| \\ \vdots \\ \langle\tilde{\varphi}_k| \end{bmatrix} \quad (2.4)$$

$$= \sum_j |\tilde{\varphi}_j\rangle\langle\tilde{\varphi}_j|. \quad (2.5)$$

From spectral theorem, density matrix ρ is decomposed as $\rho = \sum_{k=1}^d \lambda_k |k\rangle\langle k|$ where $d = \dim \mathcal{H}$. Without loss of generality, we can assume $p_k > 0$ for $k = 1 \cdots l$ where $l = \text{rank}(\rho)$ and $p_k = 0$ for $k = l+1, \cdots, d$. Thus $\rho = \sum_{k=1}^l p_k |k\rangle\langle k| = \sum_{k=1}^l |\tilde{k}\rangle\langle\tilde{k}|$, where $|\tilde{k}\rangle = \sqrt{\lambda_k} |k\rangle$.

Suppose $|\psi_i\rangle$ is a state in support ρ . Then

$$|\psi_i\rangle = \sum_{k=1}^l c_{ik} |k\rangle, \quad \sum_k |c_{ik}|^2 = 1.$$

Define $p_i = \frac{1}{\sum_k \frac{|c_{ik}|^2}{\lambda_k}}$ and $u_{ik} = \frac{\sqrt{p_i} c_{ik}}{\sqrt{\lambda_k}}$.

Now

$$\sum_k |u_{ik}|^2 = \sum_k \frac{p_i |c_{ik}|^2}{\lambda_k} = p_i \sum_k \frac{|c_{ik}|^2}{\lambda_k} = 1.$$

Next prepare an unitary operator ¹ such that i th row of U is $[u_{i1} \cdots u_{ik} \cdots u_{il}]$. Then we can define another ensemble such that

$$\left[|\tilde{\psi}_1\rangle \cdots |\tilde{\psi}_i\rangle \cdots |\tilde{\psi}_l\rangle \right] = \left[|\tilde{k}_1\rangle \cdots |\tilde{k}_l\rangle \right] U^T$$

¹By Gram-Schmidt procedure construct an orthonormal basis $\{\mathbf{u}_j\}$ (row vector) with $\mathbf{u}_i = [u_{i1} \cdots u_{ik} \cdots u_{il}]$. Then define

unitary $U = \begin{bmatrix} \mathbf{u}_1 \\ \vdots \\ \mathbf{u}_i \\ \vdots \\ \mathbf{u}_l \end{bmatrix}.$

where $|\tilde{\psi}_i\rangle = \sqrt{p_i}|\psi_i\rangle$. From theorem 2.6,

$$\rho = \sum_k |\tilde{k}\rangle\langle\tilde{k}| = \sum_k |\tilde{\psi}_k\rangle\langle\tilde{\psi}_k|.$$

Therefore we can obtain a minimal ensemble for ρ that contains $|\psi_i\rangle$.

Moreover since $\rho^{-1} = \sum_k \frac{1}{\lambda_k} |k\rangle\langle k|$,

$$\langle\psi_i|\rho^{-1}|\psi_i\rangle = \sum_k \frac{1}{\lambda_k} \langle\psi_i|k\rangle\langle k|\psi_i\rangle = \sum_k \frac{|c_{ik}|^2}{\lambda_k} = \frac{1}{p_i}.$$

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

2.74)

$$\begin{aligned}\rho_{AB} &= |a\rangle\langle a|_A \otimes |b\rangle\langle b|_B \\ \rho_A &= \text{tr}_B \rho_{AB} = |a\rangle\langle a| \text{tr}(|b\rangle\langle b|) = |a\rangle\langle a| \\ \text{tr}(\rho_A^2) &= 1\end{aligned}$$

Thus ρ_A is pure.

2.75) Define $|\Phi_{\pm}\rangle = \frac{1}{\sqrt{2}}(|00\rangle \pm |11\rangle)$ and $|\Psi_{\pm}\rangle = \frac{1}{\sqrt{2}}(|01\rangle \pm |10\rangle)$.

$$\begin{aligned}|\Phi_{\pm}\rangle\langle\Phi_{\pm}|_{AB} &= \frac{1}{2}(|00\rangle\langle 00| \pm |00\rangle\langle 11| \pm |11\rangle\langle 00| + |11\rangle\langle 11|) \\ \text{tr}_B(|\Phi_{\pm}\rangle\langle\Phi_{\pm}|_{AB}) &= \frac{1}{2}(|0\rangle\langle 0| + |1\rangle\langle 1|) = \frac{I}{2} \\ |\Psi_{\pm}\rangle\langle\Psi_{\pm}| &= \frac{1}{2}(|01\rangle\langle 01| \pm |01\rangle\langle 10| \pm |10\rangle\langle 01| + |10\rangle\langle 10|) \\ \text{tr}_B(|\Psi_{\pm}\rangle\langle\Psi_{\pm}|) &= \frac{1}{2}(|0\rangle\langle 0| + |1\rangle\langle 1|) = \frac{I}{2}\end{aligned}$$

2.76)

Unsolved. I think the polar decomposition can only apply to square matrix A , not arbitrary linear operators. Suppose A is $m \times n$ matrix. Then size of $A^\dagger A$ is $n \times n$. Thus the size of U should be $m \times n$. Maybe U is isometry, but I think it is not unitary.

I misunderstand linear operator.

Quoted from "Advanced Linear Algebra" by Steven Roman, ISBN 0387247661.

A linear transformation $\tau : V \rightarrow V$ is called a **linear operator** on V .²

Thus coordinate matrices of linear operator are square matrices. And Nielsen and Chaung say at Theorem 2.3, "Let A be a linear operator on a vector space V ." Therefore A is a linear transformation such that $A : V \rightarrow V$.

2.77)

$$\begin{aligned}|\psi\rangle &= |0\rangle |\Phi_+\rangle \\ &= |0\rangle \left[\frac{1}{\sqrt{2}}(|00\rangle + |11\rangle) \right] \\ &= (\alpha |\phi_0\rangle + \beta |\phi_1\rangle) \left[\frac{1}{\sqrt{2}}(|\phi_0\phi_0\rangle + |\phi_1\phi_1\rangle) \right]\end{aligned}$$

²According to Roman, some authors use the term linear operator for any linear transformation from V to W .

where $|\phi_i\rangle$ are arbitrary orthonormal states and $\alpha, \beta \in \mathbb{C}$. We cannot vanish cross term. Therefore $|\psi\rangle$ cannot be written as $|\psi\rangle = \sum_i \lambda_i |i\rangle_A |i\rangle_B |i\rangle_C$.

2.78)

Proof. Former part.

If $|\psi\rangle$ is product, then there exist a state $|\phi_A\rangle$ for system A , and a state $|\phi_B\rangle$ for system B such that $|\psi\rangle = |\phi_A\rangle |\phi_B\rangle$.

Obviously, this Schmidt number is 1.

Conversely, if Schmidt number is 1, the state is written as $|\psi\rangle = |\phi_A\rangle |\phi_B\rangle$. Hence this is a product state. \square

Proof. Later part.

(\Rightarrow) Proved by exercise 2.74.

(\Leftarrow) Let a pure state be $|\psi\rangle = \sum_i \lambda_i |i_A\rangle |i_B\rangle$. Then $\rho_A = \text{tr}_B(|\psi\rangle\langle\psi|) = \sum_i \lambda_i^2 |i\rangle\langle i|$. If ρ_A is a pure state, then $\lambda_j = 1$ and otherwise 0 for some j . It follows that $|\psi_j\rangle = |j_A\rangle |j_B\rangle$. Thus $|\psi\rangle$ is a product state. \square

2.79)

Procedure of Schmidt decomposition.

Goal: $|\psi\rangle = \sum_i \sqrt{\lambda_i} |i_A\rangle |i_B\rangle$

- Diagonalize reduced density matrix $\rho_A = \sum_i \lambda_i |i_A\rangle\langle i_A|$.
- Derive $|i_B\rangle$, $|i_B\rangle = \frac{(I \otimes \langle i_A|) |\psi\rangle}{\sqrt{\lambda_i}}$
- Construct $|\psi\rangle$.

(i)

$$\frac{1}{\sqrt{2}}(|00\rangle + |11\rangle) \text{ This is already decomposed.}$$

(ii)

$$\frac{|00\rangle + |01\rangle + |10\rangle + |11\rangle}{2} = \left(\frac{|0\rangle + |1\rangle}{\sqrt{2}} \right) \otimes \left(\frac{|0\rangle + |1\rangle}{\sqrt{2}} \right) = |\psi\rangle |\psi\rangle \text{ where } |\psi\rangle = \frac{|0\rangle + |1\rangle}{\sqrt{2}}$$

(iii)

$$|\psi\rangle_{AB} = \frac{1}{\sqrt{3}}(|00\rangle + |01\rangle + |10\rangle)$$

$$\rho_{AB} = |\psi\rangle\langle\psi|_{AB}$$

$$\rho_A = \text{tr}_B(\rho_{AB}) = \frac{1}{3} (2|0\rangle\langle 0| + |0\rangle\langle 1| + |1\rangle\langle 0| + |1\rangle\langle 1|)$$

$$\det(\rho_A - \lambda I) = \left(\frac{2}{3} - \lambda \right) \left(\frac{1}{3} - \lambda \right) - \frac{1}{9} = 0$$

$$\lambda^2 - \lambda + \frac{1}{9} = 0$$

$$\lambda = \frac{1 \pm \sqrt{5}/3}{2} = \frac{3 \pm \sqrt{5}}{6}$$

Eigenvector with eigenvalue $\lambda_0 \equiv \frac{3 + \sqrt{5}}{6}$ is $|\lambda_0\rangle \equiv \frac{1}{\sqrt{\frac{5+\sqrt{5}}{2}}} \begin{bmatrix} \frac{1+\sqrt{5}}{2} \\ 1 \end{bmatrix}$.

Eigenvector with eigenvalue $\lambda_1 \equiv \frac{3 - \sqrt{5}}{6}$ is $|\lambda_1\rangle \equiv \frac{1}{\sqrt{\frac{5-\sqrt{5}}{2}}} \begin{bmatrix} \frac{1-\sqrt{5}}{2} \\ 1 \end{bmatrix}$.

$$\rho_A = \lambda_0 |\lambda_0\rangle\langle\lambda_0| + \lambda_1 |\lambda_1\rangle\langle\lambda_1|.$$

$$|a_0\rangle \equiv \frac{(I \otimes \langle\lambda_0|) |\psi\rangle}{\sqrt{\lambda_0}}$$

$$|a_1\rangle \equiv \frac{(I \otimes \langle\lambda_1|) |\psi\rangle}{\sqrt{\lambda_1}}$$

Then

$$|\psi\rangle = \sum_{i=0}^1 \sqrt{\lambda_i} |a_i\rangle |\lambda_i\rangle.$$

(It's too tiresome to calculate $|a_i\rangle$)

2.80)

Let $|\psi\rangle = \sum_i \lambda_i |\psi_i\rangle_A |\psi_i\rangle_B$ and $|\varphi\rangle = \sum_i \lambda_i |\varphi_i\rangle_A |\varphi_i\rangle_B$.

Define $U = \sum_i |\psi_j\rangle\langle\varphi_j|_A$ and $V = \sum_j |\psi_j\rangle\langle\varphi_j|_B$.

Then

$$\begin{aligned} (U \otimes V) |\varphi\rangle &= \sum_i \lambda_i U |\varphi_i\rangle_A V |\varphi_i\rangle_B \\ &= \sum_i \lambda_i |\psi_i\rangle_A |\psi_i\rangle_B \\ &= |\psi\rangle. \end{aligned}$$

2.81)

Let the Schmidt decomposition of $|AR_1\rangle$ be $|AR_1\rangle = \sum_i \sqrt{p_i} |\psi_i^A\rangle |\psi_i^R\rangle$ and let $|AR_2\rangle = \sum_i \sqrt{q_i} |\phi_i^A\rangle |\phi_i^R\rangle$.

Suppose ρ^A has orthonormal decomposition $\rho^A = \sum_i p_i |i\rangle\langle i|$.

Since $|AR_1\rangle$ and $|AR_2\rangle$ are purifications of the ρ^A , we have

$$\begin{aligned} \text{tr}_R(|AR_1\rangle\langle AR_1|) &= \text{tr}_R(|AR_2\rangle\langle AR_2|) = \rho^A \\ \therefore \sum_i p_i |\psi_i^A\rangle\langle\psi_i^A| &= \sum_i q_i |\phi_i^A\rangle\langle\phi_i^A| = \sum_i \lambda_i |i\rangle\langle i|. \end{aligned}$$

The $|i\rangle$, $|\psi_i^A\rangle$, and $|\psi_i^R\rangle$ are orthonormal bases and they are eigenvectors of ρ^A . Hence without loss of generality, we can consider

$$\lambda_i = p_i = q_i \text{ and } |i\rangle = |\psi_i^A\rangle = |\phi_i^A\rangle.$$

Then

$$|AR_1\rangle = \sum_i \lambda_i |i\rangle |\psi_i^R\rangle$$

$$|AR_2\rangle = \sum_i \lambda_i |i\rangle |\phi_i^R\rangle$$

Since $|AR_1\rangle$ and $|AR_2\rangle$ have same Schmidt numbers, there are two unitary operators U and V such that $|AR_1\rangle = (U \otimes V) |AR_2\rangle$ from exercise 2.80.

Suppose $U = I$ and $V = \sum_i |\psi_i^R\rangle\langle\phi_i^R|$. Then

$$\begin{aligned} \left(I \otimes \sum_j |\psi_j^R\rangle\langle\phi_j^R| \right) |AR_2\rangle &= \sum_i \lambda_i |i\rangle \left(\sum_j |\psi_j^R\rangle \langle\phi_j^R|\phi_i^R\rangle \right) \\ &= \sum_i \lambda_i |i\rangle |\psi_i^R\rangle \\ &= |AR_1\rangle. \end{aligned}$$

Therefore there exists a unitary transformation U_R acting on system R such that $|AR_1\rangle = (I \otimes U_R) |AR_2\rangle$.

2.82)

(1)

Let $|\psi\rangle = \sum_i \sqrt{p_i} |\psi_i\rangle |i\rangle$.

$$\begin{aligned} \text{tr}_R(|\psi\rangle\langle\psi|) &= \sum_{i,j} \sqrt{p_i} \sqrt{p_j} |\psi_i\rangle\langle\psi_j| \text{tr}_R(|i\rangle\langle j|) \\ &= \sum_{i,j} \sqrt{p_i} \sqrt{p_j} |\psi_i\rangle\langle\psi_j| \delta_{ij} \\ &= \sum_i p_i |\psi_i\rangle\langle\psi_i| = \rho. \end{aligned}$$

Thus $|\psi\rangle$ is a purification of ρ .

(2)

Define the projector P by $P = I \otimes |i\rangle\langle i|$. The probability we get the result i is

$$\text{tr}[P|\psi\rangle\langle\psi|] = \langle\psi|P|\psi\rangle = \langle\psi|(I \otimes |i\rangle\langle i|)|\psi\rangle = p_i \langle\psi_i|\psi_i\rangle = p_i.$$

The post-measurement state is

$$\frac{P|\psi\rangle}{\sqrt{p_i}} = \frac{(I \otimes |i\rangle\langle i|)|\psi\rangle}{\sqrt{p_i}} = \frac{\sqrt{p_i} |\psi_i\rangle |i\rangle}{\sqrt{p_i}} = |\psi_i\rangle |i\rangle.$$

If we only focus on the state on system A ,

$$\text{tr}_R(|\psi_i\rangle |i\rangle) = |\psi_i\rangle.$$

(3)

($\{|\psi_i\rangle\}$ is not necessary an orthonormal basis.)

Suppose $|AR\rangle$ is a purification of ρ and its Schmidt decomposition is $|AR\rangle = \sum_i \sqrt{\lambda_i} |\phi_i^A\rangle |\phi_i^R\rangle$.

From assumption

$$\text{tr}_R(|AR\rangle\langle AR|) = \sum_i \lambda_i |\phi_i^A\rangle\langle\phi_i^A| = \sum_i p_i |\psi_i\rangle\langle\psi_i|.$$

By theorem 2.6, there exists a unitary matrix u_{ij} such that $\sqrt{\lambda_i} |\phi_i^A\rangle = \sum_j u_{ij} \sqrt{p_j} |\psi_j\rangle$. Then

$$\begin{aligned} |AR\rangle &= \sum_i \left(\sum_j u_{ij} \sqrt{p_j} |\psi_j\rangle \right) |\phi_i^R\rangle \\ &= \sum_j \sqrt{p_j} |\psi_j\rangle \otimes \left(\sum_i u_{ij} |\phi_i^R\rangle \right) \\ &= \sum_j \sqrt{p_j} |\psi_j\rangle |j\rangle \\ &= \sum_i \sqrt{p_i} |\psi_i\rangle |i\rangle \end{aligned}$$

where $|i\rangle = \sum_k u_{ki} |\phi_k^R\rangle$.

About $|i\rangle$,

$$\begin{aligned} \langle k|l\rangle &= \sum_{m,n} u_{mk}^* u_{nl} \langle \phi_m^R | \phi_n^R \rangle \\ &= \sum_{m,n} u_{mk}^* u_{nl} \delta_{mn} \\ &= \sum_m u_{mk}^* u_{ml} \\ &= \delta_{kl}, \quad (\because u_{ij} \text{ is unitary.}) \end{aligned}$$

which implies $|j\rangle$ is an orthonormal basis for system R .

Therefore if we measure system R w.r.t $|j\rangle$, we obtain j with probability p_j and post-measurement state for A is $|\psi_j\rangle$ from (2). Thus for any purification $|AR\rangle$, there exists an orthonormal basis $|i\rangle$ which satisfies the assertion.

Problem 2.1)

From Exercise 2.35, $\vec{n} \cdot \vec{\sigma}$ is decomposed as

$$\vec{n} \cdot \vec{\sigma} = |\lambda_1\rangle\langle\lambda_1| - |\lambda_{-1}\rangle\langle\lambda_{-1}|$$

where $|\lambda_{\pm 1}\rangle$ are eigenvector of $\vec{n} \cdot \vec{\sigma}$ with eigenvalues ± 1 .

Thus

$$\begin{aligned} f(\theta \vec{n} \cdot \vec{\sigma}) &= f(\theta) |\lambda_1\rangle\langle\lambda_1| + f(-\theta) |\lambda_{-1}\rangle\langle\lambda_{-1}| \\ &= \left(\frac{f(\theta) + f(-\theta)}{2} + \frac{f(\theta) - f(-\theta)}{2} \right) |\lambda_1\rangle\langle\lambda_1| + \left(\frac{f(\theta) + f(-\theta)}{2} - \frac{f(\theta) - f(-\theta)}{2} \right) |\lambda_{-1}\rangle\langle\lambda_{-1}| \\ &= \frac{f(\theta) + f(-\theta)}{2} (|\lambda_1\rangle\langle\lambda_1| + |\lambda_{-1}\rangle\langle\lambda_{-1}|) + \frac{f(\theta) - f(-\theta)}{2} (|\lambda_1\rangle\langle\lambda_1| - |\lambda_{-1}\rangle\langle\lambda_{-1}|) \\ &= \frac{f(\theta) + f(-\theta)}{2} I + \frac{f(\theta) - f(-\theta)}{2} \vec{n} \cdot \vec{\sigma} \end{aligned}$$

Problem 2.2) Unsolved

Problem 2.3) Unsolved

Chapter 3

Introduction to computer science

3.1) (Non-computable processes in Nature) How might we recognize that a process in Nature computes a function non computable by a Turing machine?

Soln: There are several well known non-Turing-computable functions which if identified to be computable by a process in nature would provide examples. For instance, the Halting problem: https://en.wikipedia.org/wiki/Halting_problem. More specifically, since Turing machines map non-negative integers to non-negative integers, their input and output spaces are countable (https://en.wikipedia.org/wiki/Countable_set). If any process in nature was found to compute a function taking input or providing output from an uncountable space, this could not be computed using a Turing machine. Note, Turing machines could compute the function within any desired level of approximation, but could not compute the function exactly.

3.2) (Turing numbers) Show that single-tape Turing machines can each be given a number from a list 1,2,3,... in such a way that the number uniquely specifies the corresponding machine. We call this number the *Turing number* of the corresponding machine. (*Hint:* Every positive integer has a unique prime factorization $p_1^{a_1} p_2^{a_2} \dots p_k^{a_k}$, where p_i are distinct prime numbers and a_1, \dots, a_k are non-negative integers.)

Soln: Per the hint, we show that a Turing machine can be encoded uniquely by a finite ordered list of integer values $[a_1, a_2, \dots, a_k]$. Unique prime factorization can be used to encode the Turing machine as the non-negative integer $\prod_i p_i^{a_i}$, where p_i is the i -th prime, starting with $p_1 = 2, p_2 = 3, \dots$. A non-negative integer corresponding to a Turing machine can then be decoded to reproduce the unique Turing machine whose encoding gives rise to it via the exponents in its unique prime factorization. What follows is likely overly detailed for some. The basic idea is that each part of the Turing machine can be encoded in a finite sequence of non-negative integers and decoded from that sequence. Concatenating those sequences (carefully) is then enough to specify the Turing machine. Note, it will not be the case that all non-negative integers correspond to valid Turing machines, but this is not required. We'll extend this encoding to encode Turing machines in operation, and explain how operation of a Turing machine can be simulated by multiplication of its Turing number by a rational number determined conditionally by the Turing number itself. [Note: it is unclear how useful this extension will be. The idea is relatively simple, but its formal specification is intricate and very much not necessary to understand. Feel free to skip it.]

To produce an encoding of a Turing machine, we encode each of its elements separately. We start with the finite state control. The finite state control consists of a finite set of $m+2$ states $Q = \{q_s, q_1, \dots, q_m, q_h\}$. Individually, it doesn't matter what form the q_i take, only that they are distinguishable. The integers $0, 1, \dots, m, m+1$ are distinguishable, so all that is required to encode a finite state machine is a single integer. So, setting $a_1 = m$ allows a_1 to track the size of the finite state machine and is enough to encode it.

To encode the tape, we let a_2 be the size of the alphabet Γ : $a_2 = |\Gamma|$, where here Γ includes the starting character \triangleright , corresponding to tape value 0, and blank character b corresponding to tape value

$a_2 - 1$. The other states may be assumed to be non-negative integers $1, \dots, a_2 - 2$. To encode the entirety of the tape, note that only a finite number of squares are non-blank. Let β be the largest index of a non-blank tape square. Then set $a_3 = \beta$, and for each tape square with index i , for $1 \leq i \leq \beta$, set a_{3+i} equal to the non-negative integer value assigned to the alphabet character occupying tape square i . All tape squares with index more than β are blank and need not be encoded. Note that by construction $a_4 = 0$ for all Turing machines, since tape square 1 always contains \triangleright , which was assigned value 0.

Next, we encode the program. The program contains a finite ordered list of program lines, say π of them. Set $a_{3+\beta+1} = \pi$. For $1 \leq i \leq \pi$, we encode program line i with a second prime factorization. Program line i consists of 5 elements: $\langle q_i, x_i, q'_i, x'_i, s_i \rangle$. Here, q_i and q'_i are states in Q which can be indexed with non-negative integers, say $\ell_{i,1}$ and $\ell_{i,3}$, with $0 \leq \ell_{i,1}, \ell_{i,3} \leq m + 1$. x_i and x'_i are characters in the alphabet Γ which can be indexed with non-negative integers, say $\ell_{i,2}$ and $\ell_{i,4}$, with $0 \leq \ell_{i,2}, \ell_{i,4} < |\Gamma| (= a_2)$. s_i is an integer value that is either -1, 0, or 1. Setting $\ell_{i,5} = s_i$ directly leaves open the possibility that $\ell_{i,5} = -1$, which in turn will yield non-integer encodings of the program line. There are several ways to circumvent this, the likely easiest of which is to set $\ell_{i,5} = s_i + 1$. However, the author prefers setting $\ell_{i,5} = s_i \% 3$, the remainder of s_i when divided by 3 (its residue modulo 3). This allows $s_i = 0$ and $s_i = 1$ to be encoded as $\ell_{i,5} = 0$ and $\ell_{i,5} = 1$, which are natural Boolean indicators that the tape-head should advance to the right, but requires $s_i = -1$ be encoded as $\ell_{i,5} = 2$, indicating that the tape-head should move to the left. Now, to encode program line i , for $1 \leq i \leq \pi$, set $a_{3+\beta+1+i} = 2^{\ell_{i,1}} \cdot 3^{\ell_{i,2}} \cdot 5^{\ell_{i,3}} \cdot 7^{\ell_{i,4}} \cdot 11^{\ell_{i,5}}$.

Now, for a Turing machine M , assigning Turing number $\tau(M) = \prod_{i=1}^{3+\beta+1+\pi} p_i^{a_i}$ produces an integer encoding. To show that it is unique, we reverse the encoding process and argue that all pieces of the Turing machine can be recovered uniquely from this integer value. Let an encoding of a Turing machine, $\tau(M)$, be given and begin with its unique prime factorization $\tau(M) = \prod_{i=1}^{\omega(\tau(M))} p_i^{a_i}$, where here ω is a function that returns the largest index of a prime that divides input integer. [Note, $a_4 = 0$ will mean that this isn't the number of distinct prime factors]. Immediately, we recover the size of the finite state machine, *i.e.* m . It contains a_1 states indexed by integers, along with the special starting and halting state q_s and q_h . Next, the size of the alphabet Γ is given by a_2 , where here Γ includes the starting character \triangleright and the blank character b . Next, the encoding of the tape starts with $a_3 = \beta$ which indicates the maximum index of a non-blank tape square. β encodings of tape squares follow, starting with $a_4 = 0$, indicating that tape square 1 contains the starting character \triangleright , which was assigned character value 0. If $a_4 \neq 0$, the integer provided could not be an encoding of a Turing machine, violating the assumption that $\tau(M)$ was such an integer. a_{3+i} encodes the value stored on the tape at index i , for $i \leq \beta$, where $a_{3+i} = a_2 - 1$ indicates the tape square i is blank. All tape squares with index i , for $i > \beta$, are assumed to be blank. It is left only to decode the program. We start with its length $\pi = a_{3+\beta+1}$. π encodings of individual program lines should follow, each of which should be of the form $a_{3+\beta+1+i} = 2^{\ell_{i,1}} \cdot 3^{\ell_{i,2}} \cdot 5^{\ell_{i,3}} \cdot 7^{\ell_{i,4}} \cdot 11^{\ell_{i,5}}$, from which we can recover $q_i = \ell_{i,1}$, $x_i = \ell_{i,2}$, $q'_i = \ell_{i,3}$, $x'_i = \ell_{i,4}$, and $s_i = \ell_{i,5} \% 3$, where here $\% 3$ is modular reduction on to the set of residues $-1, 0$, and 1 , instead of the standard set of residues $0, 1, 2$. Note that $r \% 3 = ((r + 1) \% 3) - 1$. It is easy to see that the program line encoded by $a_{3+\beta+1+i}$ is uniquely determined, as is the initial state of the tape from $a_3, \dots, a_{3+\beta}$. The alphabet Γ is uniquely determined by a_2 , and the finite state machine is uniquely determined by its size, given by a_1 . So, the entirety of the Turing machine M can be uniquely recovered from its Turing number $\tau(M)$, so $\tau(M)$ is unique.

(Extension): Note that, as defined, our encoding uniquely encodes Turing machines in their initial state q_s , with read-write tape-head positioned on tape square 1 holding the starting character \triangleright . The encoding scheme could be extended to encode Turing machines in operation by adding an encoding of the current state in the finite state control and current position of the read-write tape-head which will require only two additional prime factors and exponents. For compatibility, to encode the current state of the finite state machine, we use $a_{3+\beta+1+\pi+1}$, where $a_{3+\beta+1+\pi+1} = 0$ indicates the state machine is in starting state q_s , and $a_{3+\beta+1+\pi+1} = m + 1$ indicates the state machine is in state q_h and has halted. To encode the position of the read-write tape-head we require one more additional prime factor and exponent. For compatibility, we use $a_{3+\beta+1+\pi+2}$. Full compatibility of encoding will require $a_{3+\beta+1+\pi+2} = 0$ to indicate that the machine

is not yet operating and that the read-write tape-head has not yet been positioned on a tape square, neither tape square 1 holding \triangleright , as encoded by $a_4 = 0$, or another subsequent tape square holding any other value. Having $a_{3+\beta+1+\pi+2} > 0$ indicates that the Turing machine M is in operation in state specified by $a_{3+\beta+1+\pi+1}$, which we'll call σ , and read-write tape-head on the tape square specified by $a_{3+\beta+1+\pi+2}$, which we'll call σ . Note then that the tape-square pointed to by the read-write tape-head would contain the value specified by a_{3+a_σ} , which we'll call ν .

Now, to simulate execution of the Turing machine, note that in each step the program list is searched for a pattern matching its current state and the character in the tape square being read by the read-write tape-head, that is, for $\langle \sigma, \nu, \cdot, \cdot, \cdot \rangle$. This is equivalent to searching $a_{4+\beta+1}, \dots, a_{4+\beta+\pi}$ for an integer divisible by $2^\sigma \cdot 3^\nu$, but no more powers of 2 or 3. Once a matching $a_{4+\beta+i}$ is found, the multiplicities of 5, 7, and 11 in its factorization will give values for $\ell_{i,3}, \ell_{i,4}$, and $\ell_{i,5}$. The finite state machine can then be updated by multiplying by $p_{3+\beta+1+\pi+1}^{\ell_{i,3}-\sigma}$. The contents of the tape can be updated by multiplying by $p_{3+a_\sigma}^{\ell_{i,4}-\nu}$. The read-write tape-head can be moved by multiplying by $p_{3+\beta+1+\pi+2}^{\ell_{i,5} \cdot 3}$. Doing so will change the Turing number of the machine in operation M to the number encoding M after a single step.

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 and the next exercise it may help to use a mutli-tape Turing machine and/or symbols other than \triangleright , 0, 1, and the blanks.)

Soln: By “takes a binary number x as input”, what is meant is the non-blank portion of the tape contains the value x , encoded somehow. In general, the tap can hold more than just function input, but for this problem that won't be necessary. We'll use a two-tape machine, with both tapes containing symbols from the alphabet $\triangleright, 0, 1, b$. Tape 1 will contain \triangleright , followed by the input value x in binary, followed by blanks indicated with bs . The second tape will contain \triangleright and blanks. The Turing machine will populate the second tape with the reversed binary value of x , followed by blanks. It will not clear the first tape (although that could be done without too much trouble). Before we define the program, we specify that the finite state machine will contain 4 states, the starting state q_s , the halted state q_h , a search state s , and a write state w . Now, consider the program

$$P = \begin{cases} 1: \langle q_s, \triangleright, \triangleright, s, \triangleright, \triangleright, +1, 0 \rangle \\ 2: \langle s, 0, \triangleright, s, 0, \triangleright, +1, 0 \rangle \\ 3: \langle s, 1, \triangleright, s, 1, \triangleright, +1, 0 \rangle \\ 4: \langle s, b, \triangleright, w, b, \triangleright, -1, +1 \rangle \\ 5: \langle w, 0, b, w, 0, 0, -1, +1 \rangle \\ 6: \langle w, 1, b, w, 1, 1, -1, +1 \rangle \\ 7: \langle w, \triangleright, b, q_h, \triangleright, b, 0, 0 \rangle \end{cases}$$

Execution of the Turing machine begins by executing line 1 of P , which moves tape-head 1 forward, leaves tape-head 2 in place, and sets the state of the finite state machine to the search state. While in the search state, the program operates by executing lines 2 and 3, advancing tape-head 1 leaving the content of tape 1 unchanged, until it reaches a blank indicating that the end of the input has been reached, finally matching line 4. Once the blank on tape 1 is reached, line 4 changes the finite state machine to the write state, shifts tape-head 1 to the last bit of input, and advancing tape-head 2 to the first position in tape 2. Until the start of tape 1 is encountered, the program operates by executing lines 5 and 6, each of which copies the character on tape 1 pointed to by tape-head 1 onto tape 2 in the position pointed to by tape-head 2. It then moves the tape-heads in opposite directions so that tape-head 1 points to the preceeded bit of input and tape-head 2 points to the next bit of output. When the start of tape 1 is encountered, line 7 explicitly halts the program. Explicitly halting is not necessary in this case. Note, tape 1 could be cleared by replacing x'_1 in lines 5 and 6 with bs . Then, the output could be moved to tape 1 while simultaneously

clearing tape 2 by replacing line 7 with:

$$\begin{aligned}
 7 : & \langle w, \triangleright, b, w, \triangleright, b, 0, -1 \rangle \\
 8 : & \langle w, \triangleright, 0, w, \triangleright, 0, 0, -1 \rangle \\
 9 : & \langle w, \triangleright, 1, w, \triangleright, 1, 0, -1 \rangle \\
 10 : & \langle w, \triangleright, \triangleright, w, \triangleright, \triangleright, +1, +1 \rangle \\
 11 : & \langle w, b, 1, w, 1, b, +1, +1 \rangle \\
 12 : & \langle w, b, 0, w, 0, b, +1, +1 \rangle \\
 13 : & \langle w, b, b, q_h, b, b, 0, 0 \rangle
 \end{aligned}$$

Here, line 7 reverses the direction of tape-head 2. Lines 8 and 9 allow it to retreat to the start of tape 2 in line 10, at which point tape-head 1 and 2 are advanced in tandem and the blank in tape 1 is swapped with the bit in tape 2, one bit at a time, by executing lines 11 and 12. Once tape-head 2 is at the end of the reversed bit-string, line 13 is reached, explicitly halting the program.

Chapter 8

Quantum noise and quantum operations

8.1) Density operator of initial state is written by $|\psi\rangle\langle\psi|$ and final state is written by $U|\psi\rangle\langle\psi|U^\dagger$. Thus time development of $\rho = |\psi\rangle\langle\psi|$ can be written by $\mathcal{E}(\rho) = U\rho U^\dagger$.

8.2) From eqn (2.147) (on page 100),

$$\rho_m = \frac{M_m \rho M_m^\dagger}{\text{tr}(M_m^\dagger M_m \rho)} = \frac{M_m \rho M_m^\dagger}{\text{tr}(M_m \rho M_m^\dagger)} = \frac{\mathcal{E}_m(\rho)}{\text{tr} \mathcal{E}_m(\rho)}.$$

And from eqn (2.143) (on page 99), $p(m) = \text{tr}(M_m^\dagger M_m \rho) = \text{tr}(M_m \rho M_m^\dagger) = \text{tr} \mathcal{E}_m(\rho)$.

8.3)

8.4)

8.5)

8.6)

8.7)

8.8)

8.9)

8.10)

8.11)

8.12)

8.13)

8.14)

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8.30)

8.31)

8.32)

8.33)

8.34)

8.35)

Chapter 9

Distance measures for quantum information

9.1)

$$\begin{aligned} D((1, 0), (1/2, 1/2)) &= \frac{1}{2} (|1 - 1/2| + |0 - 1/2|) \\ &= \frac{1}{2} \left(\frac{1}{2} + \frac{1}{2} \right) \\ &= \frac{1}{2} \end{aligned}$$

$$\begin{aligned} D((1/2, 1/3, 1/6), (3/4, 1/8, 1/8)) &= \frac{1}{2} (|1/2 - 3/4| + |1/3 - 1/8| + |1/6 - 1/8|) \\ &= \frac{1}{2} (1/4 + 5/24 + 1/24) \\ &= \frac{1}{4} \end{aligned}$$

9.2)

$$\begin{aligned} D((p, 1-p), (q, 1-q)) &= \frac{1}{2} (|p - q| + |(1-p) - (1-q)|) \\ &= \frac{1}{2} (|p - q| + |-p + q|) \\ &= |p - q| \end{aligned}$$

9.3)

$$F((1, 0), (1/2, 1/2)) = \sqrt{1 \cdot 1/2} + \sqrt{0 \cdot 1/2} = \frac{1}{\sqrt{2}}$$

$$\begin{aligned} F((1/2, 1/3, 1/6), (3/4, 1/8, 1/8)) &= \sqrt{1/2 \cdot 3/4} + \sqrt{1/3 \cdot 1/8} + \sqrt{1/6 \cdot 1/8} \\ &= \frac{4\sqrt{6} + \sqrt{3}}{12} \end{aligned}$$

9.4)

Define $r_x = p_x - q_x$. Let U be the whole index set.

$$\begin{aligned} \max_S |p(S) - q(S)| &= \max_S \left| \sum_{x \in S} p_x - \sum_{x \in S} q_x \right| \\ &= \max_S \left| \sum_{x \in S} (p_x - q_x) \right| \\ &= \max_S \left| \sum_{x \in S} r_x \right| \end{aligned}$$

Since $\sum_{x \in S} r_x$ is written as

$$\sum_{x \in S} r_x = \sum_{\substack{x \in S \\ r_x \geq 0}} r_x + \sum_{\substack{x \in S \\ r_x < 0}} r_x, \quad (9.1)$$

$|\sum_{x \in S} r_x|$ is maximized when $S = \{x \in U | r_x \geq 0\}$ or $S = \{x \in U | r_x < 0\}$.

Define $S_+ = \{x \in U | r_x \geq 0\}$ and $S_- = \{x \in U | r_x < 0\}$.

Now the sum of all r_x is 0,

$$\begin{aligned} \sum_{x \in U} r_x &= \sum_{x \in S_+} r_x + \sum_{x \in S_-} r_x = 0 \\ \therefore \sum_{x \in S_+} r_x &= - \sum_{x \in S_-} r_x. \end{aligned}$$

Thus

$$\max_S \left| \sum_{x \in S} r_x \right| = \sum_{x \in S_+} r_x = - \sum_{x \in S_-} r_x. \quad (9.2)$$

On the other hand,

$$\begin{aligned} D(p_x, q_x) &= \frac{1}{2} \sum_{x \in U} |p_x - q_x| \\ &= \frac{1}{2} \sum_{x \in U} |r_x| \\ &= \frac{1}{2} \sum_{x \in S_+} |r_x| + \frac{1}{2} \sum_{x \in S_-} |r_x| \\ &= \frac{1}{2} \sum_{x \in S_+} r_x - \frac{1}{2} \sum_{x \in S_-} r_x \\ &= \frac{1}{2} \sum_{x \in S_+} r_x + \frac{1}{2} \sum_{x \in S_+} r_x \quad (\because \text{eqn(9.2)}) \\ &= \sum_{x \in S_+} r_x \\ &= \max_S \left| \sum_{x \in S} r_x \right|. \end{aligned}$$

Therefore $D(p_x, q_x) = \max_S |\sum_{x \in S} p_x - \sum_{x \in S} q_x| = \max_S |p(S) - q(S)|$.

9.5) From eqn (9.1) and (9.2), maximizing $|\sum_{x \in S} r_x|$ is equivalent to maximizing $\sum_{x \in S} r_x$.

Hence

$$D(p_x, q_x) = \max_S (p(S) - q(S)) = \max_S \left(\sum_{x \in S} p_x - \sum_{x \in S} q_x \right).$$

9.6)

Define $\rho = \frac{3}{4} |0\rangle\langle 0| + \frac{1}{4} |1\rangle\langle 1|$, $\sigma = \frac{2}{3} |1\rangle\langle 1| + \frac{1}{3} |1\rangle\langle 1|$.

$$\begin{aligned} D(\rho, \sigma) &= \frac{1}{2} \text{tr} |\rho - \sigma| \\ &= D((3/4, 1/4), (2/3, 1/3)) \\ &= \frac{1}{2} \left(\left| \frac{3}{4} - \frac{2}{3} \right| + \left| \frac{1}{4} - \frac{1}{3} \right| \right) \\ &= \frac{1}{2} \left(\frac{1}{12} + \frac{1}{12} \right) \\ &= \frac{1}{12} \end{aligned}$$

Define $\rho = \frac{3}{4} |0\rangle\langle 0| + \frac{1}{4} |1\rangle\langle 1|$, $\sigma = \frac{2}{3} |+\rangle\langle +| + \frac{1}{3} |-\rangle\langle -|$.

$$\begin{aligned} |+\rangle\langle +| &= \frac{1}{2} (|0\rangle\langle 0| + |0\rangle\langle 1| + |1\rangle\langle 0| + |1\rangle\langle 1|) \\ |-\rangle\langle -| &= \frac{1}{2} (|0\rangle\langle 0| - |0\rangle\langle 1| - |1\rangle\langle 0| + |1\rangle\langle 1|) \end{aligned}$$

$$\begin{aligned} \rho - \sigma &= \left(\frac{3}{4} - \frac{1}{2} \right) |0\rangle\langle 0| - \frac{1}{6} (|0\rangle\langle 1| + |1\rangle\langle 0|) + \left(\frac{1}{4} - \frac{1}{2} \right) |1\rangle\langle 1| \\ &= \frac{1}{4} |0\rangle\langle 0| - \frac{1}{6} (|0\rangle\langle 1| + |1\rangle\langle 0|) - \frac{1}{4} |1\rangle\langle 1| \end{aligned}$$

$$\begin{aligned} (\rho - \sigma)^\dagger (\rho - \sigma) &= \frac{1}{4^2} |0\rangle\langle 0| - \frac{1}{4 \cdot 6} |0\rangle\langle 1| + \frac{1}{6^2} |0\rangle\langle 0| + \frac{1}{6 \cdot 4} |0\rangle\langle 1| - \frac{1}{4 \cdot 6} |1\rangle\langle 0| + \frac{1}{6^2} |1\rangle\langle 1| + \frac{1}{4 \cdot 6} |1\rangle\langle 0| + \frac{1}{4^2} |1\rangle\langle 1| \\ &= \left(\frac{1}{4^2} + \frac{1}{6^2} \right) (|0\rangle\langle 0| + |1\rangle\langle 1|) \end{aligned}$$

$$\begin{aligned} D(\rho, \sigma) &= \frac{1}{2} \text{tr} |\rho - \sigma| \\ &= \sqrt{\frac{1}{4^2} + \frac{1}{6^2}} \end{aligned}$$

9.7)

Since $\rho - \sigma$ is Hermitian, we can apply spectral decomposition. Then $\rho - \sigma$ is written as

$$\rho - \sigma = \sum_{i=1}^k \lambda_i |i\rangle\langle i| + \sum_{i=k+1}^n \lambda_i |i\rangle\langle i|$$

where λ_i are positive eigenvalues for $i = 1, \dots, k$ and negative eigenvalues for $i = k+1, \dots, n$.

Define $Q = \sum_{i=1}^k \lambda_i |i\rangle\langle i|$ and $S = -\sum_{i=k+1}^n \lambda_i |i\rangle\langle i|$. Then P and S are positive operator. Therefore $\rho - \sigma = P - S$.

Proof of $|\rho - \sigma| = Q + S$.

$$\begin{aligned}
 |\rho - \sigma| &= |Q - S| \\
 &= \sqrt{(Q - S)^\dagger (Q - S)} \\
 &= \sqrt{(Q - S)^2} \\
 &= \sqrt{Q^2 - QS - SQ + S^2} \\
 &= \sqrt{Q^2 + S^2} \\
 &= \sqrt{\sum_i \lambda_i^2 |i\rangle\langle i|} \\
 &= \sum_i |\lambda_i| |i\rangle\langle i| \\
 &= Q + S
 \end{aligned}$$

9.8)

Suppose $\sigma = \sigma_i$. Then $\sigma = \sum_i p_i \sigma_i$.

$$D\left(\sum_i p_i \rho_i, \sigma\right) = D\left(\sum_i p_i \rho_i, \sum_i p_i \sigma_i\right) \quad (9.3)$$

$$\leq \sum_i p_i D(\rho_i, \sigma_i) \quad (\because \text{eqn(9.50)}) \quad (9.4)$$

$$= \sum_i p_i D(\rho_i, \sigma). \quad (\because \text{assumption}). \quad (9.5)$$

9.9)

9.10)

9.11)

9.12)

Suppose $\rho = \frac{1}{2}(I + \vec{r} \cdot \vec{\sigma})$ and $\sigma = \frac{1}{2}(I + \vec{s} \cdot \vec{\sigma})$ where \vec{r} and \vec{s} are real vectors s.t. $|\vec{r}|, |\vec{s}| \leq 1$.

$$\mathcal{E}(\rho) = p \frac{I}{2} + (1-p)\rho, \quad \mathcal{E}(\sigma) = p \frac{I}{2} + (1-p)\sigma.$$

$$\begin{aligned}
 D(\mathcal{E}(\rho), \mathcal{E}(\sigma)) &= \frac{1}{2} \text{tr} |\mathcal{E}(\rho) - \mathcal{E}(\sigma)| \\
 &= \frac{1}{2} \text{tr} |(1-p)(\rho - \sigma)| \\
 &= \frac{1}{2} (1-p) \text{tr} |\rho - \sigma| \\
 &= (1-p) D(\rho, \sigma) \\
 &= (1-p) \frac{|\vec{r} - \vec{s}|}{2}
 \end{aligned}$$

Is this strictly contractive?

9.13)

Bit flip channel $E_0 = \sqrt{p}I$, $E_1 = \sqrt{1-p}\sigma_x$.

$$\begin{aligned}\mathcal{E}(\rho) &= E_0\rho E_0^\dagger + E_1\rho E_1^\dagger \\ &= p\rho + (1-p)\sigma_x\rho\sigma_x.\end{aligned}$$

Since $\sigma_x\sigma_x\sigma_x = \sigma_x$, $\sigma_x\sigma_y\sigma_x = -\sigma_y$ and $\sigma_x\sigma_z\sigma_x = -\sigma_z$, then $\sigma_x(\vec{r} \cdot \vec{\sigma}) = r_1\sigma_x - r_2\sigma_y - r_3\sigma_z$.

Thus

$$\begin{aligned}D(\mathcal{E}(\rho), \mathcal{E}(\sigma)) &= \frac{1}{2} \text{tr} |\mathcal{E}(\rho) - \mathcal{E}(\sigma)| \\ &= \frac{1}{2} \text{tr} |p(\rho - \sigma) + (1-p)(\sigma_x\rho\sigma_x - \sigma_x\sigma\sigma_x)| \\ &\leq \frac{1}{2}p \text{tr} |\rho - \sigma| + \frac{1}{2}(1-p) \text{tr} |\sigma_x(\rho - \sigma)\sigma_x| \\ &= pD(\rho, \sigma) + (1-p)D(\sigma_x\rho\sigma_x, \sigma_x\sigma\sigma_x) \\ &= D(\rho, \sigma) \quad (\because \text{eqn(9.21)}).\end{aligned}$$

Suppose $\rho_0 = \frac{1}{2}(I + \vec{r} \cdot \vec{\sigma})$ is a fixed point. Then

$$\begin{aligned}\rho_0 &= \mathcal{E}(\rho_0) = p\rho_0 + (1-p)\sigma_x\rho_0\sigma_x \\ \therefore (1-p)\rho_0 - (1-p)\sigma_x\rho_0\sigma_x &= 0 \\ \therefore (1-p)(\rho - \sigma_x\rho_0\sigma_x) &= 0 \\ \therefore \rho_0 &= \sigma_x\rho_0\sigma_x \\ \therefore \frac{1}{2}(I + r_1\sigma_x + r_2\sigma_y + r_3\sigma_z) &= \frac{1}{2}(I + r_1\sigma_x - r_2\sigma_y - r_3\sigma_z)\end{aligned}$$

Since $\{I, \sigma_x, \sigma_y, \sigma_z\}$ are linearly independent, $r_2 = -r_2$ and $r_3 = -r_3$. Thus $r_2 = r_3 = 0$.

Therefore the set of fixed points for the bit flip channel is $\{\rho \mid \rho = \frac{1}{2}(I + r\sigma_x), |r| \leq 1, r \in \mathbb{R}\}$

9.14)

$$\begin{aligned}F(U\rho U^\dagger, U\sigma U^\dagger) &= \text{tr} \sqrt{(U\rho U^\dagger)^{1/2}\sigma(U\rho U^\dagger)} \\ &= \text{tr} \sqrt{U\rho^{1/2}\sigma\rho^{1/2}U^\dagger} \\ &= \text{tr}(U\sqrt{\rho^{1/2}\sigma\rho^{1/2}}U^\dagger) \\ &= \text{tr}(\sqrt{\rho^{1/2}\sigma\rho^{1/2}}U^\dagger U) \\ &= \text{tr} \sqrt{\rho^{1/2}\sigma\rho^{1/2}} \\ &= F(\rho, \sigma)\end{aligned}$$

I think the fact $\sqrt{UAU^\dagger} = U\sqrt{A}U^\dagger$ is not restricted for positive operator. Suppose A is a normal matrix. From spectral theorem, it is decomposed as

$$A = \sum_i a_i |i\rangle\langle i|.$$

Let f be a function. Then

$$\begin{aligned} f(UAU^\dagger) &= f\left(\sum_i a_i U|i\rangle\langle i|U^\dagger\right) \\ &= \sum_i f(a_i) U|i\rangle\langle i|U^\dagger \\ &= U\left(\sum_i f(a_i) |i\rangle\langle i|\right)U^\dagger \\ &= Uf(A)U^\dagger \end{aligned}$$

9.15) $|\psi\rangle = (U_R \otimes \sqrt{\rho}U_Q)|m\rangle$ is any fixed purification of ρ , and $|\phi\rangle = (V_R \otimes \sqrt{\sigma}V_Q)|m\rangle$ is purification of σ . Suppose $\sqrt{\rho}\sqrt{\sigma} = |\sqrt{\rho}\sqrt{\sigma}|V$ is the polar decomposition of $\sqrt{\rho}\sqrt{\sigma}$. Then

$$\begin{aligned} |\langle\psi|\phi\rangle| &= \left| \langle m| \left(U_R^\dagger V_R \otimes U_Q^\dagger \sqrt{\rho}\sqrt{\sigma} V_Q \right) |m\rangle \right| \\ &= \left| \text{tr} \left((U_R^\dagger V_R)^T U_Q^\dagger \sqrt{\rho}\sqrt{\sigma} V_Q \right) \right| \\ &= \left| \text{tr} \left(V_R^T U_R^* U_Q^\dagger \sqrt{\rho}\sqrt{\sigma} V_Q \right) \right| \\ &= \left| \text{tr} \left(V_Q V_R^T U_R^* U_Q^\dagger \sqrt{\rho}\sqrt{\sigma} \right) \right| \\ &= \left| \text{tr} \left(V_Q V_R^T U_R^* U_Q^\dagger |\sqrt{\rho}\sqrt{\sigma}| V \right) \right| \\ &= \left| \text{tr} \left(V V_Q V_R^T U_R^* U_Q^\dagger |\sqrt{\rho}\sqrt{\sigma}| \right) \right| \\ &\leq \text{tr} |\sqrt{\rho}\sqrt{\sigma}| \\ &= F(\rho, \sigma) \end{aligned}$$

Choosing $V_Q = V^\dagger$, $V_R^T = (U_Q^* U_R^\dagger)^\dagger$ we see that equality is attained.

9.16) I think eq (9.73) has a typo. $\text{tr}(A^\dagger B) = \langle m|A \otimes B|m\rangle$ should be $\text{tr}(A^T B) = \langle m|A \otimes B|m\rangle$. See errata list.

In order to show that this exercise, I will prove following two properties,

$$\text{tr}(A) = \langle m|(I \otimes A)|m\rangle, \quad (I \otimes A)|m\rangle = (A^T \otimes I)|m\rangle$$

where A is a linear operator and $|m\rangle$ is unnormalized maximally entangled state, $|m\rangle = \sum_i |ii\rangle$.

$$\begin{aligned} \langle m|I \otimes A|m\rangle &= \sum_{ij} \langle ii|(I \otimes A)|jj\rangle \\ &= \sum_{ij} \langle i|I|j\rangle \langle i|A|j\rangle \\ &= \sum_{ij} \delta_{ij} \langle i|A|j\rangle \\ &= \sum_i \langle i|A|i\rangle \\ &= \text{tr}(A) \end{aligned}$$

Suppose $A = \sum_{ij} a_{ij} |i\rangle\langle j|$.

$$\begin{aligned}
 (I \otimes A) |m\rangle &= \left(I \otimes \sum_{ij} a_{ij} |i\rangle\langle j| \right) \sum_k |kk\rangle \\
 &= \sum_{ijk} a_{ij} |k\rangle \otimes |i\rangle \langle j|k\rangle \\
 &= \sum_{ijk} a_{ij} |k\rangle \otimes |i\rangle \delta_{jk} \\
 &= \sum_{ij} a_{ij} |j\rangle \otimes |i\rangle \\
 &= \sum_{ij} a_{ji} |i\rangle \otimes |j\rangle
 \end{aligned}$$

$$\begin{aligned}
 (A^T \otimes I) |m\rangle &= \left(\sum_{ij} a_{ji} |i\rangle\langle j| \otimes I \right) \sum_k |kk\rangle \\
 &= \sum_{ij} a_{ji} |i\rangle \langle j|k\rangle \otimes |k\rangle \\
 &= \sum_{ij} a_{ji} |i\rangle \delta_{jk} \otimes |k\rangle \\
 &= \sum_{ij} a_{ji} |ij\rangle \\
 &= (I \otimes A) |m\rangle
 \end{aligned}$$

Thus

$$\begin{aligned}
 \text{tr}(A^T B) &= \text{tr}(B A^T) = \langle m | I \otimes B A^T | m \rangle \\
 &= \langle m | (I \otimes B) (I \otimes A^T) | m \rangle \\
 &= \langle m | (I \otimes B) (A \otimes I) | m \rangle \\
 &= \langle m | A \otimes B | m \rangle.
 \end{aligned}$$

9.17) If $\rho = \sigma$, then $F(\rho, \sigma) = 1$. Thus $A(\rho, \sigma) = \arccos F(\rho, \sigma) = \arccos 1 = 0$.

If $A(\rho, \sigma) = 0$, then $\arccos F(\rho, \sigma) = 0 \Rightarrow \cos(\arccos F(\rho, \sigma)) = \cos(0) \Rightarrow F(\rho, \sigma) = 1$ (\because text p.411, the fifth line from bottom).

9.18) For $0 \leq x \leq y \leq 1$, $\arccos(x) \geq \arccos(y)$. From $F(\mathcal{E}(\rho), \mathcal{E}(\sigma)) \geq F(\rho, \sigma)$ and $0 \leq F(\mathcal{E}(\rho), \mathcal{E}(\sigma)), F(\rho, \sigma) \leq 1$,

$$\begin{aligned}
 \arccos F(\mathcal{E}(\rho), \mathcal{E}(\sigma)) &\geq \arccos F(\rho, \sigma) \\
 \therefore A(\mathcal{E}(\rho), \mathcal{E}(\sigma)) &\geq A(\rho, \sigma)
 \end{aligned}$$

9.19) From eq (9.92)

$$\begin{aligned}
 F\left(\sum_i p_i \rho_i, \sum_i p_i \sigma_i\right) &\geq \sum_i \sqrt{p_i p_i} F(\rho_i, \sigma_i) \\
 &= \sum_i p_i F(\rho_i, \sigma_i).
 \end{aligned}$$

9.20) Suppose $\sigma_i = \sigma$. Then

$$\begin{aligned}
 F\left(\sum_i p_i \rho_i, \sigma\right) &= F\left(\sum_i p_i \rho_i, \sum_i p_i \sigma\right) \\
 &= F\left(\sum_i p_i \rho_i, \sum_i p_i \sigma_i\right) \\
 &\geq \sum_i p_i F(\rho_i, \sigma_i) \quad (\because \text{Exercise 9.19}) \\
 &= \sum_i p_i F(\rho_i, \sigma)
 \end{aligned}$$

9.21)

$$1 - F(|\psi\rangle, \sigma)^2 = 1 - \langle\psi|\sigma|\psi\rangle \quad (\because \text{eq(9.60)})$$

$$\begin{aligned}
 D(|\psi\rangle, \sigma) &= \max_P \text{tr}(P(\rho - \sigma)) \quad (\text{where } P \text{ is projector.}) \\
 &\geq \text{tr}(|\psi\rangle\langle\psi|(\rho - \sigma)) \\
 &= \langle\psi|(|\psi\rangle\langle\psi| - \sigma)|\psi\rangle \\
 &= 1 - \langle\psi|\sigma|\psi\rangle \\
 &= 1 - F(|\psi\rangle, \sigma)^2.
 \end{aligned}$$

9.22) (ref: QCQI Exercise Solutions (Chapter 9) - めもめも

<http://enakai00.hatenablog.com/entry/2018/04/12/134722>)

For all ρ , following inequality is satisfied,

$$\begin{aligned}
 d(VU\rho U^\dagger V^\dagger, \mathcal{F} \circ \mathcal{E}(\rho)) &\leq d(VU\rho U^\dagger V^\dagger, \mathcal{F}(U\rho U^\dagger)) + d(\mathcal{F}(U\rho U^\dagger), \mathcal{F} \circ \mathcal{E}(\rho)) \\
 &\leq d(VU\rho U^\dagger V^\dagger) + d(U\rho U^\dagger, \mathcal{E}(\rho)) \\
 &\leq E(V, \mathcal{F}) + E(U, \mathcal{E}).
 \end{aligned}$$

First inequality is triangular inequality, second is contractivity of the metric¹ and third is from definition of E .

Above inequality is hold for all ρ . Thus $E(VU, \mathcal{F} \circ \mathcal{E}) \leq E(V, \mathcal{F}) + E(U, \mathcal{E})$.

9.23) (\Leftarrow) If $\mathcal{E}(\rho_j) = \rho_j$ for all j such that $p_j > 0$, then

$$\bar{F} = \sum_j p_j F(\rho_j, \mathcal{E}(\rho_j))^2 = \sum_j p_j F(\rho_j, \rho_j)^2 = \sum_j p_j 1^2 = \sum_j p_j = 1.$$

(\Rightarrow) Suppose $\mathcal{E}(\rho_j) \neq \rho_j$. Then $F(\rho_j, \mathcal{E}(\rho_j)) < 1$ (\because text p.411, the fifth line from bottom). Thus

$$\bar{F} = \sum_j p_j F(\rho_j, \mathcal{E}(\rho_j))^2 < \sum_j p_j = 1.$$

Therefore if $\bar{F} = 1$, then $\mathcal{E}(\rho_j) = \rho_j$.

¹Trace distance and angle are satisfied with contractive (eq (9.35), eq (9.91)), but I don't assure that arbitrary metric satisfied with contractive.

Problem 1)

Problem 2)

Problem 3) Theorem 5.3 of "Theory of Quantum Error Correction for General Noise", Emanuel Knill, Raymond Laflamme, and Lorenza Viola, Phys. Rev. Lett. 84, 2525 – Published 13 March 2000. arXiv:quant-ph/9604034 <https://arxiv.org/abs/quant-ph/9604034>

Chapter 11

Entropy and information

11.1) Fair coin:

$$H(1/2, 1/2) = \left(-\frac{1}{2} \log \frac{1}{2}\right) \times 2 = 1 \quad (11.1)$$

Fair die:

$$H(p) = \left(-\frac{1}{6} \log \frac{1}{6}\right) \times 6 = \log 6. \quad (11.2)$$

The entropy decreases if the coin or die is unfair.

11.2)

From assumption $I(pq) = I(p) + I(q)$.

$$\frac{\partial I(pq)}{\partial p} = \frac{\partial I(p)}{\partial p} + 0 = \frac{\partial I(p)}{\partial p} \quad (11.3)$$

$$\frac{\partial I(pq)}{\partial q} = 0 + \frac{\partial I(q)}{\partial q} = \frac{\partial I(q)}{\partial q} \quad (11.4)$$

$$\frac{\partial I(pq)}{\partial p} = \frac{\partial I(pq)}{\partial(pq)} \frac{\partial(pq)}{\partial p} = q \frac{\partial I(pq)}{\partial(pq)} \Rightarrow \frac{\partial I(pq)}{\partial(pq)} = \frac{1}{q} \frac{\partial I(p)}{\partial p} \quad (11.5)$$

$$\frac{\partial I(pq)}{\partial q} = \frac{\partial I(pq)}{\partial(pq)} \frac{\partial(pq)}{\partial q} = p \frac{\partial I(pq)}{\partial(pq)} \Rightarrow \frac{\partial I(pq)}{\partial(pq)} = \frac{1}{p} \frac{\partial I(q)}{\partial q} \quad (11.6)$$

Thus

$$\frac{1}{q} \frac{\partial I(p)}{\partial p} = \frac{1}{p} \frac{\partial I(q)}{\partial q} \quad (11.7)$$

$$\therefore p \frac{dI(p)}{dp} = q \frac{dI(q)}{dq} \quad \text{for all } p, q \in [0, 1]. \quad (11.8)$$

$$(11.9)$$

Then $p(dI(p)/dp)$ is constant.

If $p(dI(p)/dp) = k$, $k \in \mathbb{R}$. Then $I(p) = k \ln p = k' \log p$ where $k' = k / \log e$.

11.3) $H_{\text{bin}}(p) = -p \log p - (1 - p) \log(1 - p)$.

$$\frac{dH_{\text{bin}}(p)}{dp} = \frac{1}{\ln 2} (-\log p - 1 + \log(1-p) + 1) \quad (11.10)$$

$$= \frac{1}{\ln 2} \ln \frac{1-p}{p} = 0 \quad (11.11)$$

$$\Rightarrow \frac{1-p}{p} = 1 \quad (11.12)$$

$$\Rightarrow p = 1/2. \quad (11.13)$$

11.4)

11.5)

$$H(p(x, y) || p(x)p(y)) = \sum_{x,y} p(x, y) \log \frac{p(x)p(y)}{p(x, y)} \quad (11.14)$$

$$= -H(p(x, y)) - \sum_{x,y} p(x, y) \log [p(x)p(y)] \quad (11.15)$$

$$= -H(p(x, y)) - \sum_{x,y} p(x, y) [\log p(x) + \log p(y)] \quad (11.16)$$

$$= -H(p(x, y)) - \sum_{x,y} p(x, y) \log p(x) - \sum_{x,y} p(x, y) \log p(y) \quad (11.17)$$

$$= -H(p(x, y)) - \sum_x p(x) \log p(x) - \sum_y p(y) \log p(y) \quad (11.18)$$

$$= -H(p(x, y)) + H(p(x)) + H(p(y)) \quad (11.19)$$

$$= -H(X, Y) + H(X) + H(Y). \quad (11.20)$$

From the non-negativity of the relative entropy,

$$H(X) + H(Y) - H(X, Y) \geq 0 \quad (11.21)$$

$$\therefore H(X) + H(Y) \geq H(X, Y). \quad (11.22)$$

11.6)

$$H(Y) + H(X, Y, Z) - H(X, Y) - H(Y, Z) = \sum_{x,y,z} p(x, y, z) \log (p(x, y)p(y, z)/p(y)p(x, y, z)) \quad (11.23)$$

$$\geq \frac{1}{\ln 2} \sum_{x,y,z} p(x, y, z) [1 - p(x, y)p(y, z)/p(y)p(x, y, z)] \quad (11.24)$$

$$= \frac{1-1}{\ln 2} = 0 \quad (11.25)$$

The equality occurs if and only if $p(x, y)p(y, z)/p(y)p(x, y, z) = 1$, which means a Markov chain condition of $Z \rightarrow Y \rightarrow X$; $p(x|y) = p(x|y, z)$

11.7)

11.8)

11.9)

11.10)

11.11)

11.12)

11.13)

11.14)
 11.15)
 11.16)
 11.17)
 11.18)
 11.19)
 11.20)
 11.21)
 11.22)
 11.23)
 11.24)
 11.25)
 11.26)

Problem 11.1)
 Problem 11.2)
 Problem 11.3)
 Problem 11.4)
 Problem 11.5)

12.31) Eve makes her qubits entangled with $|\beta_{00}\rangle$, and gets ρ^E .

$$|ABE\rangle = U |\beta_{00}^{\otimes n}\rangle |0\rangle_E \quad (11.26)$$

$$\rho^E = \text{tr}_{AB}(|ABE\rangle \langle ABE|) \quad (11.27)$$

Note that Eve's mutual information with Alice and Bob measurements does not depend on whether Eve measures ρ^E before Alice and Bob's measurement or after. So we can assume that Eve measures ρ^E after Alice and Bob's measurement. Alice and Bob measure their Bell state, getting binary string \vec{k} as an outcome. Let ρ_k^E and p_k are the corresponding Eve's states and probabilities. Note,

$$\rho_E = \sum_k p_k \rho_k^E. \quad (11.28)$$

Let K is a variable of \vec{k} and e is an outcom of a measurement of ρ^E , and E is its variable. From Holevo bound,

$$H(K : E) \leq S(\rho^E) - \sum_k p_k S(\rho_k^E) \leq S(\rho^E) = S(\rho). \quad (11.29)$$

Chapter 1

Fundamental Concepts

1.1) Probabilistic Classical Deutsch-Jozsa 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?

Soln: To a mathematician, this problem is (*slightly*) under-specified. Missing is the probability that the function f in question is balanced, vice constant. We assume that both are **equally** likely, a priori. The results when all balanced or constant functions are chosen from randomly are significantly different, and likely less interesting. We describe *an* algorithm and analyze the error rate, but make no effort to show that it is the *best* algorithm, nor that this is the most effective analysis. Let C be the event that f is constant, and B be the event that it is balanced. By hypothesis $P(C) = P(B) = \frac{1}{2}$, a priori. Evaluating f provides information which can be used to update these prior probabilities. Classically evaluating the function once, say at x_0 , provides no useful information, since comparison of values is at the heart of this problem. Evaluating f twice, say at x_0 and x_1 , can unambiguously determine if f is balanced when their values disagree. So, let's assume they agree. We use Bayesian inference to iteratively update the probability that f is constant, given k successive measurements that agree. In a convenient abuse of notation, let $P(E | k) = P(E | f(x_0) = \dots = f(x_{k-1}))$, $P(k | E) = P(f(x_0) = \dots = f(x_{k-1}) | E)$, and $P(k) = P(f(x_0) = \dots = f(x_{k-1}))$, for $E = B, C$, and $k \in \mathbb{N}$. We have $P(C | 0) = P(C | 1) = P(B | 0) = P(B | 1) = 1/2$. Note also that $P(k | C) = 1$, since if f is constant all evaluations (including the k in question) will agree. By Baye's theorem and the Law of Total Probability:

$$\begin{aligned} P(C | k) &= \frac{P(k | C) \cdot P(C | k-1)}{P(k)} \\ &= \frac{P(k | C) \cdot P(C | k-1)}{P(C | k-1) \cdot P(k | C) + P(B | k-1) \cdot P(k | B)} \end{aligned}$$

The formula above can be used to iteratively update $P(C, k)$, and hence $P(B, k) = 1 - P(C, k)$, but first we must discuss $P(k | B)$. It is important to note that when this quantity is used to update $P(C | k)$, it is already known with certainty that $f(x_0) = \dots = f(x_{k-2})$, *i.e.* $P(k-1) = 1$. $P(k | B)$ is the probability that, given this information, evaluating f one more time, at x_{k-1} , yields another value in agreement with $f(x_0), \dots, f(x_{k-2})$. We evaluate this by separating the two possible outcomes of evaluation and counting the number of balanced functions satisfying the hypotheses that would produce them. If $f(x_{k-1}) = f(x_0)$, then x_{k-1} is the k -th value on which f agrees. There are $\binom{n-k}{n/2-k}$ balanced functions which would produce this result, corresponding to the selections of $n/2 - k$ more of the remaining $n - k$ values on which f can agree. If $f(x_{k-1}) \neq f(x_0)$, then f must still agree on $n/2 - k + 1$ of the remaining $n - k$ values. There are $\binom{n-k}{n/2-k+1}$ balanced functions that would produce this result. So:

$$P(k, B) = \frac{\binom{n-k}{n/2-k}}{\binom{n-k}{n/2-k} + \binom{n-k}{n/2-k+1}} = \frac{\binom{n-k}{n/2-k}}{\binom{n-k+1}{n/2-k+1}} = \frac{n/2 - k + 1}{n - k + 1} = \frac{n - 2k + 2}{2n - 2k + 2}$$

We are finally in a position to calculate $P(C | k)$. Unfortunately, for fixed n , the machinery above does not produce formulas of bounded complexity as k grows. Each formula will be a rational function with equal degree in numerator and denominator, but those degrees seem to be $\lfloor k/2 \rfloor$. The coefficients of the leading terms show some structure that can be used for asymptotic analysis, which we do below. We illustrate the calculation of $P(C | 2)$, $P(C | 3)$, and $P(C | 4)$, and list formulas for $P(C | 5)$ through $P(C | 7)$, then discuss the results and some experimental confirmation.

$$\begin{aligned}
P(C | 2) &= \frac{P(2 | C) \cdot P(C | 1)}{P(C | 1) \cdot P(2 | C) + P(B | 1) \cdot P(2 | B)} \\
&= \frac{1 \cdot \frac{1}{2}}{\frac{1}{2} \cdot 1 + \frac{1}{2} \cdot \frac{n-2}{2n-2}} \\
&= \frac{1}{1 + \frac{n-2}{2n-2}} \\
&= \frac{2n-2}{3n-4} \\
P(C | 3) &= \frac{P(3 | C) \cdot P(C | 2)}{P(C | 2) \cdot P(3 | C) + P(B | 2) \cdot P(3 | B)} \\
&= \frac{1 \cdot \frac{2n-2}{3n-4}}{\frac{2n-2}{3n-4} + \left(1 - \frac{2n-2}{3n-4}\right) \cdot \frac{n-4}{2n-4}} \\
&= \frac{4n-4}{5n-8} \\
P(C | 4) &= \frac{P(4 | C) \cdot P(C | 3)}{P(C | 3) \cdot P(4 | C) + P(B | 3) \cdot P(4 | B)} \\
&= \frac{1 \cdot \frac{4n-4}{5n-8}}{\frac{4n-4}{5n-8} + \left(1 - \frac{4n-4}{5n-8}\right) \cdot \frac{n-6}{2n-6}} \\
&= \frac{8n^2 - 32n + 24}{9n^2 - 42n + 48} \\
P(C | 5) &= \frac{16n^2 - 64n + 48}{17n^2 - 78n + 96} \\
P(C | 6) &= \frac{32n^3 - 288n^2 + 736n - 480}{33n^3 - 312n^2 + 924n - 960} \\
P(C | 7) &= \frac{64n^3 - 576n^2 + 1472n - 960}{65n^3 - 606n^2 + 1768n - 1920}
\end{aligned}$$

There are clearly patterns, the most striking of which yields $P(C | k) \xrightarrow{n \rightarrow \infty} \frac{2^{k-1}}{2^{k-1}+1}$, that is, given $k \geq 2$ evaluations in agreement, the probability that f is constant is $\sim 1 - \frac{1}{2^{k-1}+1}$, at least for large n . In the quantum context, where n is likely to be exponential in the number of qubits, this asymptotic value would be approached rapidly. To confirm this analysis, a python script is included in the repo which experimentally calculates empirical values of $P(C | k)$ for specified values of n and k . It also calculates the theoretical values, recursing over k , for comparison. See `<git repo>/Python/Problem1.1.py`.

To answer the problem most directly, *i.e.*, “what is the performance of the best classical algorithm for this problem?”, let n be fixed and $0 < \epsilon < \frac{1}{2}$ be specified. The “performance” of classically evaluating the function of n inputs in order to declare it constant with error less than ϵ is equivalent to determining the number k of evaluations in agreement after which the probability that f is constant is greater than $1 - \epsilon$. Note that in no case is this number less than two. The entries in the table below are such values, with rows indexed by n , and columns corresponding to exponentially decreasing values of ϵ . Specifically, column i lists the values of k corresponding to $\epsilon = 1/2^i$. The maximum value of k in each row is $n/2 + 1$, since this implies the function is constant.

$\epsilon = 1/2^i; i =$	1	2	3	4	5	6	7	8	9	10	11	12	13	14	15	16	17	18	19	20
$n = 6$	2	3	3	4	→															
$n = 8$	2	3	4	4	4	5	→													
$n = 10$	2	3	4	4	5	5	6	→												
$n = 12$	2	3	4	4	5	5	6	6	7	→										
$n = 14$	2	3	4	5	5	6	6	7	7	7	8	→								
$n = 16$	2	3	4	5	5	6	6	7	7	8	8	8	9	→						
$n = 18$	2	3	4	5	5	6	7	7	8	8	8	9	9	9	10	→				
$n = 20$	2	3	4	5	6	6	7	7	8	8	9	9	9	10	10	10	11	→		
$n = 22$	2	3	4	5	6	6	7	7	8	9	9	9	10	10	11	11	11	11	12	→
$n = 24$	2	3	4	5	6	6	7	8	8	9	9	10	10	11	11	11	12	12	12	12
$n = 26$	2	3	4	5	6	6	7	8	8	9	9	10	10	11	11	12	12	12	13	13
$n = 28$	2	3	4	5	6	7	7	8	8	9	10	10	11	11	12	12	12	13	13	13
$n = 30$	2	3	4	5	6	7	7	8	9	9	10	10	11	11	12	12	13	13	13	14
$n = 32$	2	3	4	5	6	7	7	8	9	9	10	11	11	12	12	13	13	13	14	14
$n = 34$	2	3	4	5	6	7	7	8	9	9	10	11	11	12	12	13	13	14	14	15
$n = 36$	2	3	4	5	6	7	7	8	9	10	10	11	11	12	12	13	13	14	14	15
$n = 38$	2	3	4	5	6	7	8	8	9	10	10	11	12	12	13	13	14	14	15	15
$n = 40$	2	3	4	5	6	7	8	8	9	10	10	11	12	12	13	13	14	14	15	15
$n = 42$	2	3	4	5	6	7	8	8	9	10	10	11	12	12	13	14	14	15	15	16
$n = 44$	2	3	4	5	6	7	8	8	9	10	11	11	12	13	13	14	14	15	15	16
$n = 46$	2	3	4	5	6	7	8	8	9	10	11	11	12	13	13	14	14	15	16	16
$n = 48$	2	3	4	5	6	7	8	8	9	10	11	11	12	13	13	14	15	15	16	16
$n = 50$	2	3	4	5	6	7	8	9	9	10	11	11	12	13	13	14	15	15	16	16
$n = 52$	2	3	4	5	6	7	8	9	9	10	11	12	12	13	14	14	15	15	16	17
$n = 54$	2	3	4	5	6	7	8	9	9	10	11	12	12	13	14	14	15	16	16	17
$n = 56$	2	3	4	5	6	7	8	9	9	10	11	12	12	13	14	14	15	16	16	17
$n = 58$	2	3	4	5	6	7	8	9	9	10	11	12	12	13	14	15	15	16	16	17
$n = 60$	2	3	4	5	6	7	8	9	9	10	11	12	13	13	14	15	15	16	17	17
$n = 62$	2	3	4	5	6	7	8	9	10	10	11	12	13	13	14	15	15	16	17	17
$n = 64$	2	3	4	5	6	7	8	9	10	10	11	12	13	13	14	15	15	16	17	17

Loosely, for small numbers of evaluations and large n , *i.e.* in the bottom left of the table, each exponential increase in the probability of being constant desired requires an additional evaluation. Eventually, the combinatorial reduction in the number of remaining balanced functions allows additional evaluations to reduce the error with which the function can be declared constant by several powers of 2, as often seen in the top right. That is not to say that the probability is always reduced by at least a factor of 2. In fact, note that $P(C | 1) = 1/2$, and $P(C | 2) \xrightarrow{n \rightarrow \infty} 2/3$, so $\frac{1-P(C | 1)}{1-P(C | 2)} \xrightarrow{n \rightarrow \infty} 3/2$. The second evaluation only reduces the probability that the function is balanced by a factor of ~ 1.5 for large n . Asymptotically, for $k \geq 2$, note that $\frac{1-P(C | k+1)}{1-P(C | k)} \xrightarrow{n \rightarrow \infty} \frac{\frac{1}{2^{k+1}}}{\frac{1}{2^k+1}} = \frac{2^{k-1}+1}{2^k+1} < 2$, so all k -th evaluations eventually reduce the probability of the function being balanced by less than a factor of 2, for large enough n . Theoretically, it is seemingly possible there's a case in which halving the probability of being balanced requires two additional evaluations. That is, there could exist n and an $\epsilon = \frac{1}{2^i}$ requiring k measurements to declare the function constant with error less than ϵ and at least $k+2$ measurements to declare the function constant with error less than $\epsilon/2$. However, attempts to search for such a pathological case have come up empty. The asymptotic short-fallings are overcome by the combinatorial reduction fast enough, before a power of $1/2$ straddles two values of k . It is likely that more careful analysis could refute the possibility rigorously.

We finish discussion of this problem with a (perhaps unnecessary) table of values of $P(C | k)$ for fixed n and k (programmatically constructed with the python script mentioned above, as was the previous table.) Again, once $k = n/2 + 1$, the function must be constant, so all probabilities are 1.

k	2	3	4	5	6	7	8	9	10
$n = 4$	$\frac{3}{4} \simeq 0.7500$	$\frac{1}{1} \simeq 1.0000$	$\frac{1}{1} \simeq 1.0000$	$\frac{1}{1} \simeq 1.0000$	$\frac{1}{1} \simeq 1.0000$	$\frac{1}{1} \simeq 1.0000$	$\frac{1}{1} \simeq 1.0000$	$\frac{1}{1} \simeq 1.0000$	$\frac{1}{1} \simeq 1.0000$
$n = 6$	$\frac{5}{7} \simeq 0.7143$	$\frac{10}{11} \simeq 0.9091$	$\frac{35}{36} \simeq 0.9722$	$\frac{126}{127} \simeq 0.9921$	$\frac{462}{463} \simeq 0.9978$	$\frac{1716}{1717} \simeq 0.9994$	$\frac{6435}{6436} \simeq 0.9998$	$\frac{24310}{24311} \simeq 1.0000$	$\frac{92378}{92379} \simeq 1.0000$
$n = 8$	$\frac{7}{10} \simeq 0.7000$	$\frac{7}{8} \simeq 0.8750$	$\frac{35}{36} \simeq 0.9722$	$\frac{126}{127} \simeq 0.9921$	$\frac{462}{463} \simeq 0.9978$	$\frac{1716}{1717} \simeq 0.9994$	$\frac{6435}{6436} \simeq 0.9998$	$\frac{24310}{24311} \simeq 1.0000$	$\frac{92378}{92379} \simeq 1.0000$
$n = 10$	$\frac{9}{13} \simeq 0.6923$	$\frac{6}{7} \simeq 0.8571$	$\frac{21}{22} \simeq 0.9545$	$\frac{66}{67} \simeq 0.9851$	$\frac{429}{430} \simeq 0.9954$	$\frac{323}{324} \simeq 0.9969$	$\frac{4199}{4200} \simeq 0.9993$	$\frac{8398}{8399} \simeq 0.9999$	$\frac{92378}{92379} \simeq 1.0000$
$n = 12$	$\frac{11}{16} \simeq 0.6875$	$\frac{11}{13} \simeq 0.8462$	$\frac{33}{35} \simeq 0.9429$	$\frac{66}{67} \simeq 0.9851$	$\frac{429}{430} \simeq 0.9954$	$\frac{323}{324} \simeq 0.9969$	$\frac{4199}{4200} \simeq 0.9993$	$\frac{8398}{8399} \simeq 0.9999$	$\frac{92378}{92379} \simeq 1.0000$
$n = 14$	$\frac{13}{19} \simeq 0.6842$	$\frac{26}{31} \simeq 0.8387$	$\frac{143}{153} \simeq 0.9346$	$\frac{146}{147} \simeq 0.9795$	$\frac{431}{432} \simeq 0.9954$	$\frac{323}{324} \simeq 0.9969$	$\frac{4199}{4200} \simeq 0.9993$	$\frac{8398}{8399} \simeq 0.9999$	$\frac{92378}{92379} \simeq 1.0000$
$n = 16$	$\frac{15}{22} \simeq 0.6818$	$\frac{5}{6} \simeq 0.8333$	$\frac{13}{14} \simeq 0.9286$	$\frac{39}{40} \simeq 0.9750$	$\frac{143}{144} \simeq 0.9931$	$\frac{715}{716} \simeq 0.9986$	$\frac{6435}{6436} \simeq 0.9998$	$\frac{24310}{24311} \simeq 1.0000$	$\frac{92378}{92379} \simeq 1.0000$
$n = 18$	$\frac{17}{25} \simeq 0.6800$	$\frac{34}{41} \simeq 0.8293$	$\frac{85}{92} \simeq 0.9239$	$\frac{34}{35} \simeq 0.9714$	$\frac{221}{223} \simeq 0.9910$	$\frac{442}{443} \simeq 0.9977$	$\frac{24310}{24311} \simeq 1.0000$	$\frac{92378}{92379} \simeq 1.0000$	$\frac{92378}{92379} \simeq 1.0000$
$n = 20$	$\frac{19}{28} \simeq 0.6786$	$\frac{19}{23} \simeq 0.8261$	$\frac{323}{353} \simeq 0.9202$	$\frac{646}{667} \simeq 0.9685$	$\frac{646}{653} \simeq 0.9893$	$\frac{323}{324} \simeq 0.9969$	$\frac{4199}{4200} \simeq 0.9993$	$\frac{8398}{8399} \simeq 0.9999$	$\frac{92378}{92379} \simeq 1.0000$
$n = 22$	$\frac{21}{31} \simeq 0.6774$	$\frac{14}{17} \simeq 0.8235$	$\frac{133}{145} \simeq 0.9172$	$\frac{57}{59} \simeq 0.9661$	$\frac{323}{327} \simeq 0.9878$	$\frac{1292}{1297} \simeq 0.9961$	$\frac{969}{970} \simeq 0.9990$	$\frac{4522}{4523} \simeq 0.9998$	$\frac{29393}{29394} \simeq 1.0000$
$n = 24$	$\frac{23}{34} \simeq 0.6765$	$\frac{23}{28} \simeq 0.8214$	$\frac{161}{176} \simeq 0.9148$	$\frac{161}{167} \simeq 0.9641$	$\frac{437}{443} \simeq 0.9865$	$\frac{437}{439} \simeq 0.9954$	$\frac{7429}{7439} \simeq 0.9987$	$\frac{14858}{14863} \simeq 0.9997$	$\frac{14858}{14859} \simeq 0.9999$
$n = 26$	$\frac{25}{37} \simeq 0.6757$	$\frac{50}{61} \simeq 0.8197$	$\frac{116}{126} \simeq 0.9127$	$\frac{230}{239} \simeq 0.9623$	$\frac{805}{817} \simeq 0.9853$	$\frac{575}{578} \simeq 0.9948$	$\frac{10925}{10943} \simeq 0.9984$	$\frac{2185}{2186} \simeq 0.9995$	$\frac{37145}{37149} \simeq 0.9999$
$n = 28$	$\frac{27}{40} \simeq 0.6750$	$\frac{9}{11} \simeq 0.8182$	$\frac{225}{247} \simeq 0.9109$	$\frac{270}{281} \simeq 0.9609$	$\frac{690}{701} \simeq 0.9843$	$\frac{345}{347} \simeq 0.9942$	$\frac{1035}{1037} \simeq 0.9981$	$\frac{1725}{1726} \simeq 0.9994$	$\frac{6555}{6556} \simeq 0.9998$
$n = 30$	$\frac{29}{43} \simeq 0.6744$	$\frac{58}{71} \simeq 0.8169$	$\frac{261}{287} \simeq 0.9094$	$\frac{261}{272} \simeq 0.9596$	$\frac{1305}{1327} \simeq 0.9834$	$\frac{1740}{1751} \simeq 0.9937$	$\frac{10005}{10027} \simeq 0.9978$	$\frac{10005}{10012} \simeq 0.9993$	$\frac{10005}{10007} \simeq 0.9998$
$n = 32$	$\frac{31}{46} \simeq 0.6739$	$\frac{31}{38} \simeq 0.8158$	$\frac{899}{990} \simeq 0.9081$	$\frac{899}{938} \simeq 0.9584$	$\frac{8091}{8234} \simeq 0.9826$	$\frac{8091}{8146} \simeq 0.9932$	$\frac{4495}{4506} \simeq 0.9976$	$\frac{13485}{13496} \simeq 0.9992$	$\frac{310155}{310232} \simeq 0.9998$
$n = 34$	$\frac{33}{49} \simeq 0.6735$	$\frac{22}{27} \simeq 0.8148$	$\frac{341}{376} \simeq 0.9069$	$\frac{2046}{2137} \simeq 0.9574$	$\frac{9889}{10071} \simeq 0.9819$	$\frac{1798}{1811} \simeq 0.9928$	$\frac{24273}{24338} \simeq 0.9973$	$\frac{5394}{5399} \simeq 0.9991$	$\frac{13485}{13489} \simeq 0.9997$
$n = 36$	$\frac{35}{52} \simeq 0.6731$	$\frac{35}{43} \simeq 0.8140$	$\frac{77}{85} \simeq 0.9059$	$\frac{222}{233} \simeq 0.9565$	$\frac{682}{695} \simeq 0.9813$	$\frac{1705}{1718} \simeq 0.9924$	$\frac{4495}{4508} \simeq 0.9971$	$\frac{12586}{12599} \simeq 0.9990$	$\frac{37758}{37771} \simeq 0.9997$
$n = 38$	$\frac{37}{55} \simeq 0.6727$	$\frac{74}{91} \simeq 0.8132$	$\frac{1295}{1431} \simeq 0.9050$	$\frac{259}{271} \simeq 0.9557$	$\frac{407}{415} \simeq 0.9807$	$\frac{1628}{1641} \simeq 0.9921$	$\frac{12617}{12656} \simeq 0.9969$	$\frac{11470}{11483} \simeq 0.9989$	$\frac{33263}{33276} \simeq 0.9996$
$n = 40$	$\frac{39}{58} \simeq 0.6724$	$\frac{13}{16} \simeq 0.8125$	$\frac{481}{532} \simeq 0.9041$	$\frac{1443}{1511} \simeq 0.9550$	$\frac{3367}{3435} \simeq 0.9802$	$\frac{481}{485} \simeq 0.9918$	$\frac{1221}{1225} \simeq 0.9967$	$\frac{814}{815} \simeq 0.9988$	$\frac{2294}{2295} \simeq 0.9996$
$n = 42$	$\frac{41}{61} \simeq 0.6721$	$\frac{82}{101} \simeq 0.8119$	$\frac{533}{590} \simeq 0.9034$	$\frac{1066}{1117} \simeq 0.9543$	$\frac{19721}{20129} \simeq 0.9797$	$\frac{19721}{19891} \simeq 0.9915$	$\frac{19721}{19789} \simeq 0.9966$	$\frac{1517}{1519} \simeq 0.9987$	$\frac{16687}{16695} \simeq 0.9995$
$n = 44$	$\frac{43}{64} \simeq 0.6719$	$\frac{43}{53} \simeq 0.8113$	$\frac{1763}{1953} \simeq 0.9027$	$\frac{3526}{3697} \simeq 0.9537$	$\frac{45838}{46807} \simeq 0.9793$	$\frac{22919}{23123} \simeq 0.9912$	$\frac{848003}{851063} \simeq 0.9964$	$\frac{65231}{65265} \simeq 0.9986$	$\frac{65231}{65265} \simeq 0.9995$
$n = 46$	$\frac{45}{67} \simeq 0.6716$	$\frac{30}{37} \simeq 0.8108$	$\frac{129}{143} \simeq 0.9021$	$\frac{387}{406} \simeq 0.9532$	$\frac{1763}{1801} \simeq 0.9789$	$\frac{35260}{35583} \simeq 0.9909$	$\frac{343785}{345077} \simeq 0.9963$	$\frac{22919}{22953} \simeq 0.9985$	$\frac{848003}{848479} \simeq 0.9994$
$n = 48$	$\frac{47}{70} \simeq 0.6714$	$\frac{47}{58} \simeq 0.8103$	$\frac{705}{782} \simeq 0.9015$	$\frac{141}{148} \simeq 0.9527$	$\frac{6063}{6196} \simeq 0.9785$	$\frac{2021}{2040} \simeq 0.9907$	$\frac{82861}{83184} \simeq 0.9961$	$\frac{414305}{414951} \simeq 0.9984$	$\frac{1077193}{1077839} \simeq 0.9994$
$n = 50$	$\frac{49}{73} \simeq 0.6712$	$\frac{98}{121} \simeq 0.8099$	$\frac{2303}{2556} \simeq 0.9010$	$\frac{658}{691} \simeq 0.9522$	$\frac{987}{1009} \simeq 0.9782$	$\frac{1974}{1993} \simeq 0.9905$	$\frac{14147}{14204} \simeq 0.9960$	$\frac{198058}{198381} \simeq 0.9984$	$\frac{4060189}{4062773} \simeq 0.9994$
$n = 52$	$\frac{51}{76} \simeq 0.6711$	$\frac{17}{21} \simeq 0.8095$	$\frac{833}{925} \simeq 0.9005$	$\frac{4998}{5251} \simeq 0.9518$	$\frac{11186}{11439} \simeq 0.9779$	$\frac{5593}{5648} \simeq 0.9903$	$\frac{50337}{50546} \simeq 0.9959$	$\frac{11186}{11205} \simeq 0.9983$	$\frac{28294}{28313} \simeq 0.9993$
$n = 54$	$\frac{53}{79} \simeq 0.6709$	$\frac{106}{131} \simeq 0.8092$	$\frac{901}{1001} \simeq 0.9001$	$\frac{901}{947} \simeq 0.9514$	$\frac{44149}{45161} \simeq 0.9776$	$\frac{25228}{25481} \simeq 0.9901$	$\frac{296429}{297694} \simeq 0.9958$	$\frac{592858}{593903} \simeq 0.9982$	$\frac{296429}{296638} \simeq 0.9993$
$n = 56$	$\frac{55}{82} \simeq 0.6707$	$\frac{55}{68} \simeq 0.8088$	$\frac{583}{648} \simeq 0.8997$	$\frac{583}{613} \simeq 0.9511$	$\frac{9911}{10141} \simeq 0.9773$	$\frac{4505}{4551} \simeq 0.9899$	$\frac{31535}{31673} \simeq 0.9956$	$\frac{12614}{12659} \simeq 0.9982$	$\frac{592858}{593295} \simeq 0.9993$
$n = 58$	$\frac{57}{85} \simeq 0.6706$	$\frac{38}{47} \simeq 0.8085$	$\frac{1045}{1162} \simeq 0.8993$	$\frac{1254}{1319} \simeq 0.9507$	$\frac{11077}{11337} \simeq 0.9771$	$\frac{11077}{11192} \simeq 0.9897$	$\frac{51357}{51587} \simeq 0.9955$	$\frac{85595}{85756} \simeq 0.9981$	$\frac{119833}{119925} \simeq 0.9992$
$n = 60$	$\frac{59}{88} \simeq 0.6705$	$\frac{59}{73} \simeq 0.8082$	$\frac{1121}{1247} \simeq 0.8990$	$\frac{2242}{2359} \simeq 0.9504$	$\frac{24662}{25247} \simeq 0.9768$	$\frac{12331}{12461} \simeq 0.9896$	$\frac{653543}{656533} \simeq 0.9954$	$\frac{59413}{59528} \simeq 0.9981$	$\frac{1010021}{1010082} \simeq 0.9992$
$n = 62$	$\frac{61}{91} \simeq 0.6703$	$\frac{122}{151} \simeq 0.8079$	$\frac{3599}{4005} \simeq 0.8986$	$\frac{3599}{3788} \simeq 0.9501$	$\frac{68381}{70019} \simeq 0.9766$	$\frac{273524}{276449} \simeq 0.9894$	$\frac{752191}{755701} \simeq 0.9954$	$\frac{752191}{753686} \simeq 0.9980$	$\frac{3624193}{3627183} \simeq 0.9992$
$n = 64$	$\frac{63}{94} \simeq 0.6702$	$\frac{21}{26} \simeq 0.8077$	$\frac{1281}{1426} \simeq 0.8983$	$\frac{549}{578} \simeq 0.9498$	$\frac{3599}{3686} \simeq 0.9764$	$\frac{3599}{3638} \simeq 0.9893$	$\frac{68381}{68706} \simeq 0.9953$	$\frac{478667}{479642} \simeq 0.9980$	$\frac{5265337}{5269822} \simeq 0.9991$