

**Notes on “Separability in asymmetric phase-covariant cloning”
(by Rezakhani, Siadatnejad, and Ghaderi)**

These notes provide detailed calculations pertaining to the paper of Rezakhani, Siadatnejad, and Ghaderi [11].

1. UNIVERSAL ASYMMETRIC CLONING MACHINE

The goal of this section is to develop an asymmetric cloning machine which is universal. In this context, ‘universal’ means that it treats all inputs in the same way. We will begin by setting up a symmetric universal cloning machine and reviewing some of its properties. Following this we will develop an asymmetric cloner and explore some of its properties.

1.1. The d -dimensional universal cloning machine of Buzek and Hillery [3]. Consider the unitary transformation

$$|i\rangle_A |0\rangle_B |\Sigma\rangle_X \mapsto \mu |i\rangle_A |i\rangle_B |i\rangle_X + \nu \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \left(|i\rangle_A |j\rangle_B + |j\rangle_A |i\rangle_B \right) |j\rangle_X, \quad (1) \quad \{\text{universal}\}$$

where A is the input qudit, B is the blank qudit, and X is an ancilla which is initially in a fixed state, $|\Sigma\rangle$. (RSG comment that the ancilla can be regarded as the cloning machine; I would like to understand this comment.) The set $\{|i\rangle_A\}_{i=0}^{d-1}$ is a set of orthonormal basis vectors for the Hilbert space of the input, \mathcal{H}_A . We make similar definitions for the Hilbert space of the ancilla \mathcal{H}_X .

Fact 1. We can, without loss of generality, take the parameters μ, ν to be real numbers.

We require the following conditions of the mapping (1):

- (1) the mapping is unitary;
- (2) the mapping is *universal*, meaning that the quality of the clones does not depend on the input state (measured in terms of the fidelity $F = \langle \psi | \rho^{(\text{out})} | \psi \rangle$);
- (3) the outputs are symmetric, i.e., $\rho_A^{(\text{out})} = \rho_B^{(\text{out})}$.
- (4) While this isn’t required in the RSG paper, there are authors who require the mapping for the cloning transformation to be completely positive (in the operator-theoretic sense), see for instance [7] or, in the phase-covariant case, [1].

Remark 1. There is a fairly good set of definitions provided in the paper of Horodecki, Horodecki, and Horodecki [6] that has to do with all of the various notions of positivity for operators. If \mathcal{A}_i is the algebra of operators on the Hilbert space \mathcal{H}_i ($i = 1, 2$), then \mathcal{A}_i is a Hilbert space with inner product $\langle A, B \rangle = \text{tr}(B^\dagger A)$. We denote by $\mathcal{L}(\mathcal{A}_1, \mathcal{A}_2)$ the space of linear maps from \mathcal{A}_1 to \mathcal{A}_2 . A map $\Lambda \in \mathcal{L}(\mathcal{A}_1, \mathcal{A}_2)$ is *positive* if it maps positive operators in \mathcal{A}_1 to positive operators in \mathcal{A}_2 . We say that Λ is a *completely positive map* if the induced map

$$\Lambda_n := \Lambda \otimes I : \mathcal{A}_1 \otimes \mathcal{M}_n \rightarrow \mathcal{A}_2 \otimes \mathcal{M}_n$$

is positive for all n . (Here \mathcal{M}_n is the set of all complex $n \times n$ matrices and I is the identity map.)

Lemma 1.1. $\text{tr}_{BX} \left(|ijk\rangle \langle pqr| \right) = \begin{cases} |i\rangle \langle p| & j = q \text{ and } k = r \\ 0 & \text{otherwise.} \end{cases}$

Proof. We recall that the partial trace (over B) of a composite system AB is defined as follows (see [9]):

$$\text{tr}_B \left(|a_1\rangle \langle a_2| \otimes |b_1\rangle \langle b_2| \right) := |a_1\rangle \langle a_2| \text{tr} \left(|b_1\rangle \langle b_2| \right),$$

where $|a_1\rangle, |a_2\rangle$ are vectors in the state space of A and $|b_1\rangle, |b_2\rangle$ are vectors in the state space of B . We note that the trace functional appearing on the right side of this equation is the usual trace:

$$\text{tr} \left(|b_1\rangle \langle b_2| \right) = \langle b_1 | b_2 \rangle.$$

In our case we have

$$\text{tr}_{BX} \left(|ijk\rangle \langle pqr| \right) = |i\rangle \langle p| \text{tr} \left(|jk\rangle \langle qr| \right).$$

Note that

$$\begin{aligned}\mathrm{tr}\left(|jk\rangle\langle qr|\right) &= \mathrm{tr}\left(|j\rangle\langle q|\right)\mathrm{tr}\left(|k\rangle\langle r|\right) \\ &= \langle j|q\rangle\langle k|r\rangle \\ &= \delta_{jq}\delta_{kr},\end{aligned}$$

where

$$\delta_{mn} = \begin{cases} 1 & m = n \\ 0 & m \neq n. \end{cases}$$

□

partialtrace}

Lemma 1.2. $\mathrm{tr}_{AX}\left(|ijk\rangle\langle pqr|\right) = \begin{cases} |j\rangle\langle q| & i = p \text{ and } k = r \\ 0 & \text{otherwise.} \end{cases}$

Proof. The proof of this result is almost identical to that of Lemma 1.1.

□

CMidentities}

Proposition 1.1. *In light of (1)-(3) we have the following identities:*

iversal_rhoA}

$$\rho_A^{(\mathrm{out})} = \eta\rho_A^{(\mathrm{id})} + \frac{1-\eta}{d}1_A \quad (2)$$

iversal_rhoB}

$$\rho_B^{(\mathrm{out})} = \eta\rho_B^{(\mathrm{id})} + \frac{1-\eta}{d}1_B \quad (3)$$

{mu_nu_d}

$$\mu^2 = 2\mu\nu \quad \mu^2 = \frac{2}{d+1} \quad \nu^2 = \frac{1}{2(d+1)} \quad (4)$$

_universal_d}

$$\eta = \mu^2 + (d-2)\nu^2 = \frac{d+2}{2(d+1)}. \quad (5)$$

I would like to know why the clone has to have the form in (2) and (3).

The reason for this is developed in [1], where complete positivity of the cloning transformation allows the authors to use the Kraus decomposition of the reduced density operator. There is also (allegedly) a proof contained in [7].

Alternatively, in [3] it is claimed that the particular form of the output density operator minimizes the Bures distance:

$$d(\rho_1, \rho_2) = \sqrt{2}\left(1 - \mathrm{tr}\left(\sqrt{\rho_1^{1/2}\rho_2\rho_1^{1/2}}\right)\right)^{1/2}$$

Unfortunately there is no argument accompanying this claim.

Proof. The key component of this proof (and for some of the proofs that will come later in these notes) is the calculation of the density operator $\rho^{(\mathrm{out})}$ and the subsequent calculation of the reduced density operators $\rho_A^{(\mathrm{out})}$ and $\rho_B^{(\mathrm{out})}$.

We begin with an input of the form $|\psi\rangle = \sum_{i=0}^{d-1} \alpha_i |i\rangle$. Thanks to linearity we find that the output state of the operator defined by (1) is

$$\sum_{i=0}^{d-1} \alpha_i \left(\mu |iii\rangle + \nu \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \left(|ijj\rangle + |jij\rangle \right) \right)$$

(Note that we are now writing $|ijk\rangle = |i\rangle_A |j\rangle_B |k\rangle_X$ to simplify the notation.) The density operator of the output is

$$\begin{aligned}
\rho^{(\text{out})} &= \left[\sum_{i=0}^{d-1} \alpha_i \left(\mu |iii\rangle + \nu \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} (|ikk\rangle + |kik\rangle) \right) \right] \\
&\quad \left[\sum_{j=0}^{d-1} \alpha_j^* \left(\mu \langle jjj| + \nu \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} (\langle jll| + \langle ljl|) \right) \right] \\
&= \mu^2 \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* |iii\rangle \langle jjj| \\
&\quad + \mu\nu \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} (|ikk\rangle \langle jjj| + |kik\rangle \langle jjj|) \\
&\quad + \mu\nu \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} (|iii\rangle \langle jll| + |iii\rangle \langle ljl|) \\
&\quad + \nu^2 \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} (|ikk\rangle \langle jll| + |ikk\rangle \langle ljl| + |kik\rangle \langle jll| + |kik\rangle \langle ljl|) \\
&=: \text{Term I} + \text{Term II} + \text{Term III} + \text{Term IV}.
\end{aligned}$$

To compute the reduced density operators we consider the partial trace of Terms I-IV.

Calculation of $\rho_A^{(\text{out})}$. To calculate the reduced density operator we now calculate the partial trace over the B and X qudits:

$$\rho_A^{(\text{out})} = \text{tr}_{BX}(\rho^{(\text{out})}).$$

Term I.

$$\begin{aligned}
\text{tr}_{BX}(\text{Term I}) &= \mu^2 \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \text{tr}_{BX}(|iii\rangle \langle jjj|) \\
&= \mu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i|,
\end{aligned}$$

where we've used the fact that

$$\text{tr}_{BX}(|iii\rangle \langle jjj|) = |i\rangle \langle i|$$

in the case where $i = j$ and is zero otherwise (this follows from Lemma 1.1).

Term II.

$$\begin{aligned}
\text{tr}_{BX}(\text{Term II}) &= \mu\nu \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} \text{tr}_{BX}(|ikk\rangle \langle jjj|) \\
&\quad + \mu\nu \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} \text{tr}_{BX}(|kik\rangle \langle jjj|) \\
&=: \text{Term II}_1 + \text{Term II}_2.
\end{aligned}$$

In Term II_1 we note that we only obtain a nontrivial contribution to the sum if $j = k \neq i$, in which case we have $\text{tr}_{BX}(|ikk\rangle\langle jjj|) = |i\rangle\langle j|$. Thus

$$\text{Term II}_1 = \mu\nu \sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* |i\rangle\langle j|.$$

In Term II_2 we have a nontrivial contribution to the sum if $j = k$ and $j = i$. However, since $k \neq i$ in the internal most sum, we find that $\text{Term II}_2 = 0$.

In summary, then, we have

$$\text{tr}_{BX}(\text{Term II}) = \mu\nu \sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* |i\rangle\langle j|.$$

Term III.

$$\begin{aligned} \text{tr}_{BX}(\text{Term III}) &= \mu\nu \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} \text{tr}_{BX}(|iii\rangle\langle jll|) \\ &\quad + \mu\nu \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} \text{tr}_{BX}(|iii\rangle\langle ljl|) \\ &=: \text{Term III}_1 + \text{Term III}_2. \end{aligned}$$

Note that $\text{tr}_{BX}(|iii\rangle\langle jll|) = 0$ unless $i = l \neq j$. As in the case of Term II_1 above we find that

$$\text{Term III}_1 = \sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* |i\rangle\langle j|.$$

(Here it helps to switch the order of summation and then to use symmetry in reordering the sum.)

In order for the terms in Term III_2 to be nonzero we require that $i = l = j$. However, since the internal sum requires $l \neq j$ we find that $\text{Term III}_2 = 0$.

We thus find that

$$\text{tr}_{BX}(\text{Term III}) = \sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* |i\rangle\langle j|.$$

Term IV. As we did in the cases of Term II and Term III, we decompose the partial trace of Term IV into four terms:

$$\text{tr}_{BX}(\text{Term IV}) = \text{Term IV}_1 + \text{Term IV}_2 + \text{Term IV}_3 + \text{Term IV}_4,$$

where

$$\begin{aligned}
\text{Term IV}_1 &= \nu^2 \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ d \neq i}} \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} \text{tr}_{BX} \left(|ikk\rangle \langle jll| \right), \\
\text{Term IV}_2 &= \nu^2 \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ d \neq i}} \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} \text{tr}_{BX} \left(|ikk\rangle \langle ljl| \right), \\
\text{Term IV}_3 &= \nu^2 \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ d \neq i}} \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} \text{tr}_{BX} \left(|kik\rangle \langle jll| \right), \\
\text{Term IV}_4 &= \nu^2 \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ d \neq i}} \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} \text{tr}_{BX} \left(|kik\rangle \langle ljl| \right).
\end{aligned}$$

Calculation of Term IV₁. In light of Lemma 1.1 we find that

$$\text{tr}_{BX} \left(|ikk\rangle \langle jll| \right) = \begin{cases} |i\rangle \langle j| & k = l \\ 0 & k \neq l. \end{cases}$$

It is helpful to further subdivide the resulting expression into the two cases in which $i = j$ and $i \neq j$. If $i = j$, then

$$\begin{aligned}
\sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} \text{tr}_{BX} \left(|ikk\rangle \langle jll| \right) &= \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} \sum_{\substack{0 \leq l \leq d-1 \\ l \neq i}} \text{tr}_{BX} \left(|ikk\rangle \langle ill| \right) \\
&= \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} |i\rangle \langle i| \\
&= (d-1) |i\rangle \langle i|.
\end{aligned}$$

On the other hand, if $i \neq j$, then

$$\sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} \text{tr}_{BX} \left(|ikk\rangle \langle jll| \right) = \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i \\ k \neq j}} |i\rangle \langle j| = (d-2) |i\rangle \langle j|$$

Together this yields

$$\text{Term IV}_1 = (d-1) \nu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i| + (d-2) \nu^2 \sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* |i\rangle \langle j|$$

Calculation of Term IV₂. Using Lemma 1.1, we find that $\text{tr}_{BX} \left(|ikk\rangle \langle ljl| \right) = |i\rangle \langle j|$ whenever $j = k = l$. However the sum over l specifies the $l \neq j$, meaning that there are no nontrivial terms in this sum. We conclude that Term IV₂ = 0.

Calculation of Term IV₃. As in our calculation of Term IV₂ we have that $\text{tr}_{BX} \left(|kik\rangle \langle jll| \right) \neq 0$ provided $k = l = i$. However this is forbidden in the sum over k (we require $k \neq i$), whence Term IV₃ = 0.

Calculation of Term IV₄. Using Lemma 1.1 here gives

$$\text{tr}_{BX} \left(|kik\rangle \langle ljl| \right) = \begin{cases} |k\rangle \langle l| & i = j \text{ and } k = l \\ 0 & \text{otherwise} \end{cases} = \begin{cases} |k\rangle \langle k| & i = j, k = l \\ 0 & \text{otherwise.} \end{cases}$$

It follows that

$$\begin{aligned}
\text{Term IV}_4 &= \nu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} |k\rangle \langle k| \\
&= \nu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 \left(1_A - |i\rangle \langle i| \right) \\
&= \nu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 1_A - \nu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i| \\
&= \nu^2 1_A - \nu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i|.
\end{aligned}$$

Here we've used that $\sum_{i=0}^{d-1} |\alpha_i|^2 = 1$ along with the simplifying calculation

$$\sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} |i\rangle \langle i| = 1_A - |i\rangle \langle i|,$$

where 1_A denotes the identity operator on the state space of A .

To summarize, we have

$$\begin{aligned}
\text{tr}_{BX}(\text{Term IV}) &= (d-1)\nu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i| + (d-2)\nu^2 \sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* |i\rangle \langle j| + \nu^2 1_A - \nu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i| \\
&= (d-2)\nu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i| + (d-2)\nu^2 \sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* |i\rangle \langle j| + \nu^2 1_A \\
&= (d-2)\nu^2 |\psi\rangle \langle \psi| + \nu^2 1_A.
\end{aligned}$$

Putting these calculations together, we find that

$$\begin{aligned}
\rho_A^{(\text{out})} &= \mu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i| + 2\mu\nu \sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* |i\rangle \langle j| + (d-2)\nu^2 |\psi\rangle \langle \psi| + \nu^2 1_A \\
&= \left(\mu^2 + (d-2)\nu^2 \right) |\psi\rangle \langle \psi| + \nu^2 1_A,
\end{aligned}$$

where we require that $\mu^2 = 2\mu\nu$. Here we are using the fact that the reduced density operator has the form (2) to identify μ^2 and $2\mu\nu$. Our calculation reveals that $\eta = \mu^2 + (d-2)\nu^2$, as in (5).

To see that we now have the remaining identities in (4), note that the requirement that the cloning map (1) be unitary implies that

$$\mu^2 + 2(d-1)\nu^2 = 1. \tag{6}$$

From $\mu^2 = 2\mu\nu$ we see that $\mu = 0$ or $\mu = 2\nu$. If $\mu = 2\nu$, then (6) reads

$$4\nu^2 + 2(d-1)\nu^2 = 1$$

which yields $\nu^2 = 1/2(d+1)$. Returning this to (6) and solving for μ^2 yields the final identity of (4).

Why are we forbidden from taking $\mu = 0$?

Calculation of $\rho_B^{(\text{out})}$. The calculation of the reduced density operator $\rho_B^{(\text{out})}$ closely follows the corresponding calculation for $\rho_A^{(\text{out})}$ but relies on Lemma 1.2. We summarize the calculations below:

$$\begin{aligned}
\text{tr}_{AX}(\text{Term I}) &= \mu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i| \\
\text{tr}_{AX}(\text{Term II}) &= \mu\nu \sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* |i\rangle \langle j| \\
\text{tr}_{AX}(\text{Term III}) &= \mu\nu \sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* |i\rangle \langle j| \\
\text{tr}_{AX}(\text{Term IV}) &= \nu^2 1_B - \nu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i| + (d-1)\nu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i| + (d-2)\nu^2 \sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* |i\rangle \langle j| \\
&= (d-2)\nu^2 |\psi\rangle \langle \psi| + \nu^2 1_B.
\end{aligned}$$

We conclude that $\rho_B^{(\text{out})} = \rho_A^{(\text{out})}$ automatically.

Why do RSG insist on this condition? It seems it's automatic once one has the form (1).

Some clarification is provided in [3] on this point. They have a nicer way of explaining exactly what we require of a UQCM. \square

Note. For ease of reference we have included the main calculations in a table at the end of these notes.

Remark 2. Using the same techniques as those used in the proof of Proposition 1.1 we can also calculate the reduced operator for the ancilla.

$$\begin{aligned}
\rho_X^{(\text{out})} &= \text{tr}_{AB}(\rho^{(\text{out})}) \\
&= \text{tr}_{AB}(\text{Term I}) + \text{tr}_{AB}(\text{Term II}) + \text{tr}_{AB}(\text{Term III}) + \text{tr}_{AB}(\text{Term IV}),
\end{aligned}$$

where Term I through Term IV are as above. In the case of the ancilla we find that

$$\text{tr}_{AB}(\text{Term I}) = \mu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i|.$$

Interestingly, we find that

$$\text{tr}_{AB}(\text{Term II}) = 0,$$

since we encounter terms of the form

$$\text{tr}_{AB}(|ikk\rangle \langle jjj|) \quad \text{and} \quad \text{tr}_{AB}(|kik\rangle \langle jjj|).$$

In each of these cases the partial trace is zero unless $i = j = k$. However, the internal sum is only taken over $k \neq i$, meaning that the partial trace is, in this case, trivial. Similarly, we find that

$$\text{tr}_{AB}(\text{Term III}) = 0.$$

Term IV is, as always, more complicated. We consider four cases, corresponding to the four terms in the sum for Term IV:

$$\begin{aligned}
\text{Term IV}_1 &= \nu^2 \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} \text{tr}_{AB}(|ikk\rangle \langle jll|) \\
\text{Term IV}_2 &= \nu^2 \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} \text{tr}_{AB}(|ikk\rangle \langle ljl|) \\
\text{Term IV}_3 &= \nu^2 \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} \text{tr}_{AB}(|kik\rangle \langle jll|) \\
\text{Term IV}_4 &= \nu^2 \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} \text{tr}_{AB}(|kik\rangle \langle ljl|).
\end{aligned}$$

For Term IV₁ we find that the partial trace is zero unless $i = j$ and $k = l$. Thus

$$\begin{aligned}
\text{Term IV}_1 &= \nu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} |k\rangle \langle k| \\
&= \nu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 (1 - |i\rangle \langle i|) \\
&= \nu^2 1 - \nu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i|.
\end{aligned}$$

For Term IV₂ we have

$$\text{tr}_{AB}(|ikk\rangle \langle ljl|) = \begin{cases} |j\rangle \langle i| & \text{if } i = l \text{ and } j = k \\ 0 & \text{otherwise.} \end{cases}$$

This means that

$$\text{Term IV}_2 = \nu^2 \sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* |j\rangle \langle i|.$$

Term IV₃ is similar to Term IV₂ and also works out to be

$$\text{Term IV}_3 = \nu^2 \sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* |j\rangle \langle i|.$$

Finally, Term IV₄ is similar to Term IV₁ and is given by

$$\text{Term IV}_4 = \nu^2 1 - \nu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i|.$$

Altogether this yields

$$\rho_X^{(\text{out})} = (\mu^2 - 2\nu^2) \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i| + 2\nu^2 \sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* |j\rangle \langle i| + 2\nu^2 1.$$

There is a formula for the output in the ancilla given in Buzek and Hillery [3] which seems to be different. It may be that they have simplified their calculation somewhat using the identities developed in the proposition. Or it could just be that there is an error in the calculation that is developed in these notes.

We note that

$$\left(\sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* |j\rangle \langle i| \right) = \left(\sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* |i\rangle \langle j| \right)^\dagger,$$

meaning that we have an adjoint term involved.

1.2. Asymmetric Cloning Machine. In tracing through the calculations for the universal cloning machine given in (1) we observe that the symmetry $\rho_A^{(\text{out})} = \rho_B^{(\text{out})}$ is a result of the symmetry in the summation in (1). That is, the symmetry results from the fact that terms of the form

$$|i\rangle_A |j\rangle_B |j\rangle_X \quad \text{and} \quad |j\rangle_A |i\rangle_B |j\rangle_X$$

are equally weighted by the coefficient ν . To break this symmetry, then, it is reasonable to allow for different coefficients on these two terms. With that in mind, we consider a unitary transformation defined by

$$|i\rangle_A |0\rangle_B |\Sigma\rangle_X \mapsto \mu |i\rangle_A |i\rangle_B |i\rangle_X + \nu \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} |i\rangle_A |j\rangle_B |j\rangle_X + \xi \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} |j\rangle_A |i\rangle_B |j\rangle_X. \quad (7)$$

Proposition 1.2. Suppose that the initial state we wish to clone is given by $|\psi\rangle = \sum_{i=0}^{d-1} \alpha_i |i\rangle$. The reduced density operators associated with the transformation (7) are given by

$$\rho_A^{(\text{out})} = \left((d-2)\nu^2 + 2\mu\nu \right) |\psi\rangle \langle \psi| + \xi^2 1 + \left(\mu^2 + \nu^2 - 2\mu\nu - \xi^2 \right) \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i|,$$

and

$$\rho_B^{(\text{out})} = \left((d-2)\xi^2 + 2\mu\xi \right) |\psi\rangle \langle \psi| + \nu^2 1 + \left(\mu^2 + \xi^2 - 2\mu\xi - \nu^2 \right) \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i|.$$

Proof. The reduced density operators that are the conclusion of Proposition 1.2 are computed as the partial traces of the density operator

$$\begin{aligned} \rho^{(\text{out})} = & \mu^2 \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* |iii\rangle \langle jjj| + \mu\nu \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} |iii\rangle \langle jll| \\ & + \mu\xi \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} |iii\rangle \langle ljl| + \mu\nu \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} |ikk\rangle \langle jjj| \\ & + \nu^2 \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} |ikk\rangle \langle jll| + \nu\xi \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} |ikk\rangle \langle ljl| \\ & + \mu\xi \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} |kik\rangle \langle jjj| + \nu\xi \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} |kik\rangle \langle jll| \\ & + \xi^2 \sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} |kik\rangle \langle ljl| \end{aligned}$$

We now compute $\rho_A^{(\text{out})} = \text{tr}_{BX}(\rho^{(\text{out})})$. The calculations involved are the same as those that were used in the proof of Proposition 1.1. For convenience we have collected those calculations into a table located on

the last page of these notes. We find that

$$\begin{aligned}
\rho_A^{(\text{out})} &= \text{tr}_{BX}(\rho^{(\text{out})}) \\
&= \mu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i| + 2\mu\nu \sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* |i\rangle \langle j| \\
&\quad + (d-1)\nu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i| + (d-2)\nu^2 \sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* |i\rangle \langle j| \\
&\quad + \xi^2 1 - \xi^2 \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i|.
\end{aligned}$$

We use the identity

$$\sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* |i\rangle \langle j| = |\psi\rangle \langle \psi| - \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i|,$$

which is obtained by calculating $|\psi\rangle \langle \psi|$ and separating the resulting sum into diagonal terms and off-diagonal terms. This identity enables us to simplify our calculation, yielding

$$\begin{aligned}
\rho_A^{(\text{out})} &= \mu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i| + 2\mu\nu |\psi\rangle \langle \psi| - 2\mu\nu \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i| \\
&\quad + (d-2)\nu^2 |\psi\rangle \langle \psi| + \nu^2 \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i| + \xi^2 1 - \xi^2 \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i| \\
&= \left((d-2)\nu^2 + 2\mu\nu \right) |\psi\rangle \langle \psi| + \xi^2 1 + \left(\mu^2 + \nu^2 - 2\mu\nu - \xi^2 \right) \sum_{i=0}^{d-1} |\alpha_i|^2 |i\rangle \langle i|.
\end{aligned}$$

An analogous calculation yields the result for $\rho_B^{(\text{out})} = \text{tr}_{AX}(\rho^{(\text{out})})$. □

As a result of Proposition 1.2 we have the following identities:

$$\begin{aligned}
\eta_A &= (d-2)\nu^2 + 2\mu\nu, & \mu^2 + \nu^2 - 2\mu\nu - \xi^2 &= 0, \\
\eta_B &= (d-2)\xi^2 + 2\mu\xi, & \mu^2 + \xi^2 - 2\mu\xi - \nu^2 &= 0.
\end{aligned}$$

These inequalities follow from the calculations in Proposition 1.2 together with the form of the density operators specified in (2), (3). We note that the sum that occurs in our formulas for $\rho_A^{(\text{out})}$ and $\rho_B^{(\text{out})}$ is state-dependent, so its coefficient must be zero.

The fidelities for these clones are given by

$$F_A = (d-2)\nu^2 + 2\mu\nu + \xi^2 \quad \text{and} \quad F_B = (d-2)\xi^2 + 2\mu\xi + \nu^2.$$

2. PHASE-COVARIANT CLONING

In this section we introduce a cloning machine of the form (7) which acts on a restricted class of qudits, namely those that lie in the equatorial plane:

$$|\psi\rangle = \frac{1}{\sqrt{d}} \sum_{k=0}^{d-1} e^{i\phi_k} |k\rangle, \tag{8}$$

where $\phi_k \in [0, 2\pi)$ for each $k \in \{0, \dots, d-1\}$. Notice that in this case we have $\alpha_k = e^{i\phi_k}/\sqrt{d}$, whence $|\alpha_k|^2 = 1/d$ for each $k \in \{0, \dots, d-1\}$. Returning to the calculations from Proposition 1.2 we find that the sum in the last term is

$$\sum_{k=0}^{d-1} |\alpha_k|^2 |k\rangle \langle k| = \frac{1}{d} \sum_{k=0}^{d-1} |k\rangle \langle k| = \frac{1}{d} 1.$$

In particular, we observe that the density operator is no longer state dependent. Indeed,

$$\rho_A^{(\text{out})} = \left((d-2)\nu^2 + 2\mu\nu \right) |\psi\rangle \langle\psi| + \left(\xi^2 + \frac{\mu^2 + \nu^2 - 2\mu\nu - \xi^2}{d} \right) 1,$$

and

$$\rho_B^{(\text{out})} = \left((d-2)\xi^2 + 2\mu\xi \right) |\psi\rangle \langle\psi| + \left(\nu^2 + \frac{\mu^2 + \xi^2 - 2\mu\xi - \nu^2}{d} \right) 1.$$

This means that the fidelity of the clones satisfy

$$F_A = \frac{1 + (d-1)\left((d-2)\xi^2 + 2\mu\xi\right)}{d} \quad \text{and} \quad F_B = \frac{1 + (d-1)\left((d-2)\xi^2 + 2\mu\xi\right)}{d}.$$

2.1. Optimization. We say that the clones developed above are *optimal* if, whenever the fidelity of one of the clones is fixed, the fidelity of the other clone is as large as possible. Thus, for instance, if we fix the value of F_A , the value of F_B should be as large as possible. {optimization}

We will focus here on the case in which $d = 2$. The general case appears to be treated in [8], though the discussion there doesn't appear to be particularly edifying.

Suppose that F_A is fixed. Our goal is to determine the values of the coefficients μ, ν, ξ so that the fidelity of the A -clone is F_A while the fidelity of the B clone is as large as possible. Since $F_A = (1 + \eta_A)/2$ and $F_B = (1 + \eta_B)/2$, we see that it is sufficient to treat the corresponding (constrained) optimization for η_A and η_B . Namely, we regard η_A as being a fixed constant and our goal is to determine the largest possible value for η_B . We are thus lead to formulate the following constrained optimization problem:

$$\text{maximize } \eta_B = 2\mu\xi \quad \text{given that} \quad 2\mu\nu = \eta_A \quad \text{and} \quad \mu^2 + \nu^2 + \xi^2 = 1. \quad (9) \quad \text{{opt_prob}}$$

Again, for clarity, we point out that η_A is viewed here as being constant, meaning that the formula $2\mu\nu = \eta_A$ provides a constraint on the possible values of μ and ν .

To solve the optimization problem (9) we use Lagrange multipliers; we seek $\lambda_1, \lambda_2, \mu, \nu, \xi$ satisfying

$$\nabla \eta_B = \lambda_1 \nabla g_1(\mu, \nu, \xi) + \lambda_2 \nabla g_2(\mu, \nu, \xi) \quad (10) \quad \text{{grad_cond}}$$

$$g_1(\mu, \nu, \xi) = 0 \quad (11) \quad \text{{g1zero}}$$

$$g_2(\mu, \nu, \xi) = 0 \quad (12) \quad \text{{g2zero}}$$

where

$$g_1(\mu, \nu, \xi) = 2\mu\nu - \eta_A \quad \text{and} \quad g_2(\mu, \nu, \xi) = \mu^2 + \nu^2 + \xi^2 - 1.$$

Calculating the gradients and writing equation (10) in components leaves us to solve the following system of (nonlinear) equations:

$$\xi = \lambda_1 \nu + \lambda_2 \mu \quad (13) \quad \text{{systemeq1}}$$

$$0 = \lambda_1 \mu + \lambda_2 \nu \quad (14) \quad \text{{systemeq2}}$$

$$\mu = \lambda_2 \xi \quad (15) \quad \text{{systemeq3}}$$

$$2\mu\nu = \eta_A \quad (16) \quad \text{{systemeq4}}$$

$$\mu^2 + \nu^2 + \xi^2 = 1. \quad (17) \quad \text{{systemeq5}}$$

We substitute (15) into (14) and cancel a factor of λ_2 to find that $\nu = -\lambda_1 \xi$. (Note that if $\lambda_2 = 0$, then $\mu = 0$, whence $\eta_A = 0$.) Inserting these new relationships into (13) we find that $\lambda_2^2 - \lambda_1^2 = 1$. (Here we have canceled a factor of ξ ; if $\xi = 0$ we find that we must have $\eta_A = 0$.) Further, using (17), we obtain

$$\lambda_2^2 \xi^2 + \lambda_1^2 \xi + \xi^2 = 1. \quad (18) \quad \text{{opt_eq_1}}$$

Using the identity $\lambda_2^2 - \lambda_1^2 = 1$ and solving (18) for ξ yields

$$\xi = \pm \frac{1}{\sqrt{2}\lambda_2}.$$

Recalling that $\mu = \lambda_2 \xi$, we thus have

$$\mu = \pm \frac{1}{\sqrt{2}}.$$

As $2\mu\nu = \eta_A$ (which we recall is a constant specified in advance) we have

$$\nu = \pm \frac{1}{\sqrt{2}} \eta_A.$$

Finally we return to the normalization condition (17) to see that

$$\frac{1}{2} + \frac{1}{2} \eta_A^2 + \xi^2 = 1.$$

Solving for ξ yields

$$\xi = \pm \sqrt{\frac{1 - \eta_A^2}{2}}.$$

These calculations indicate that the optimal values of η_B are given by

$$\eta_B = 2\mu\xi = \pm \sqrt{1 - \eta_A^2}.$$

In particular, we see that for the optimal phase-covariant clones, the shrinking factors η_A, η_B satisfy the circle relation

$$\eta_A^2 + \eta_B^2 = 1.$$

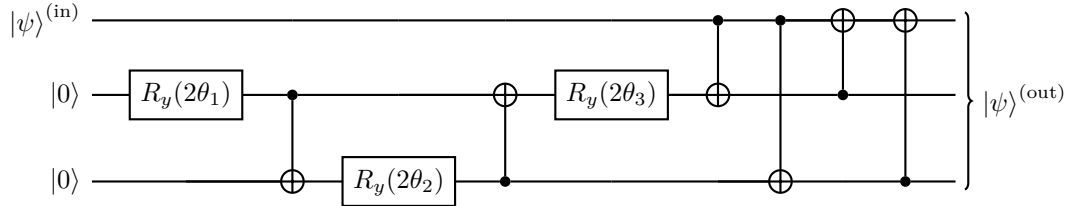
2.1.1. The Symmetric Case. In the case where the clones are symmetric, we find that $\eta_A = \eta_B = 1/\sqrt{2}$. In this case we find that the fidelities are given by

$$F_A = F_B = \frac{1 + \frac{1}{\sqrt{2}}}{2} = \frac{1}{2} + \frac{1}{2\sqrt{2}}.$$

This agrees with previous calculations for the case of the symmetric phase-covariant cloning machines given in [4].

3. IMPLEMENTATION

In this section we will focus on a circuit that can be used to implement phase-covariant cloning on equatorial qubits as described above. The circuit that we consider for this implementation is as follows:



This circuit was originally developed in [2]. Note that the R_y gate has the matrix representation

$$R_y(\theta) = \begin{pmatrix} \cos\left(\frac{\theta}{2}\right) & -\sin\left(\frac{\theta}{2}\right) \\ \sin\left(\frac{\theta}{2}\right) & \cos\left(\frac{\theta}{2}\right) \end{pmatrix},$$

meaning that by doubling the angles entered into the gate we remove the denominators in the trigonometric functions that are the entries of the matrix.

Our goal in this section is to calculate the output for this circuit when $|\psi\rangle^{(\text{in})} = \frac{1}{\sqrt{2}}(|0\rangle + e^{i\phi}|1\rangle)$ and to select angles $\theta_1, \theta_2, \theta_3$ that realize the optimal cloning described in the preceding section.

3.1. Determining Circuit Output. To simplify the exposition we will focus on the circuit output for the qubits $|0\rangle$ and $|1\rangle$ and then use linearity to express the output for the input state $|\psi\rangle^{(\text{in})}$.

3.1.1. *Input State* $|0\rangle$. In the case when the input state is $|0\rangle$ the initial state of the (composite) system is $|000\rangle$. Upon application of the first rotation (on the 1-qubit) our system is of the form

$$\cos(\theta_1) |000\rangle + \sin(\theta_1) |010\rangle.$$

Next we apply our first CNOT gate to get

$$\cos(\theta_1) |000\rangle + \sin(\theta_1) |011\rangle.$$

The rotation through $2\theta_2$ in the 2-qubit yields

$$\begin{aligned} & \cos(\theta_1) \left(\cos(\theta_2) |000\rangle + \sin(\theta_2) |001\rangle \right) + \sin(\theta_1) \left(-\sin(\theta_2) |010\rangle + \cos(\theta_2) |011\rangle \right) \\ &= \cos(\theta_1) \cos(\theta_2) |000\rangle + \cos(\theta_1) \sin(\theta_2) |001\rangle - \sin(\theta_1) \sin(\theta_2) |010\rangle + \sin(\theta_1) \cos(\theta_2) |011\rangle. \end{aligned}$$

The second CNOT gate applied to this state leaves us with

$$\cos(\theta_1) \cos(\theta_2) |000\rangle + \cos(\theta_1) \sin(\theta_2) |011\rangle - \sin(\theta_1) \sin(\theta_2) |010\rangle + \sin(\theta_1) \cos(\theta_2) |001\rangle.$$

The final rotation (through an angle of $2\theta_3$ in the 1-qubit) gives

$$\begin{aligned} & \cos(\theta_1) \cos(\theta_2) \left(\cos(\theta_3) |000\rangle + \sin(\theta_3) |010\rangle \right) + \cos(\theta_1) \sin(\theta_2) \left(-\sin(\theta_3) |001\rangle + \cos(\theta_3) |011\rangle \right) \\ & - \sin(\theta_1) \sin(\theta_2) \left(-\sin(\theta_3) |000\rangle + \cos(\theta_3) |010\rangle \right) + \sin(\theta_1) \cos(\theta_2) \left(\cos(\theta_3) |001\rangle + \sin(\theta_3) |011\rangle \right) \\ &= \cos(\theta_1) \cos(\theta_2) \cos(\theta_3) |000\rangle + \cos(\theta_1) \cos(\theta_2) \sin(\theta_3) |010\rangle - \cos(\theta_1) \sin(\theta_2) \sin(\theta_3) |001\rangle + \cos(\theta_1) \sin(\theta_2) \cos(\theta_3) |011\rangle \\ & + \sin(\theta_1) \sin(\theta_2) \sin(\theta_3) |000\rangle - \sin(\theta_1) \sin(\theta_2) \cos(\theta_3) |010\rangle + \sin(\theta_1) \cos(\theta_2) \cos(\theta_3) |001\rangle + \sin(\theta_1) \cos(\theta_2) \sin(\theta_3) |011\rangle \\ &= \left(\cos(\theta_1) \cos(\theta_2) \cos(\theta_3) + \sin(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) |000\rangle \\ & + \left(\cos(\theta_1) \cos(\theta_2) \sin(\theta_3) - \sin(\theta_1) \sin(\theta_2) \cos(\theta_3) \right) |010\rangle \\ & + \left(\sin(\theta_1) \cos(\theta_2) \cos(\theta_3) - \cos(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) |001\rangle \\ & + \left(\cos(\theta_1) \sin(\theta_2) \cos(\theta_3) + \sin(\theta_1) \cos(\theta_2) \sin(\theta_3) \right) |011\rangle. \end{aligned}$$

This concludes the preparatory phase of the circuit. We point out that the input $|\psi\rangle^{(\text{in})}$ has not yet been used in our calculations. That is, only the 1- and 2- qubits have been affected. This fact will be important when we consider the case when the input is $|1\rangle$.

Turning to the remainder of the circuit we have four CNOT gates that we must now pass through. Since the 0-qubit is in the state $|0\rangle$, the first two CNOT gates do not affect the state. The third CNOT gate (with the control on the 1-qubit) yields

$$\begin{aligned} & \left(\cos(\theta_1) \cos(\theta_2) \cos(\theta_3) + \sin(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) |000\rangle \\ & + \left(\cos(\theta_1) \cos(\theta_2) \sin(\theta_3) - \sin(\theta_1) \sin(\theta_2) \cos(\theta_3) \right) |110\rangle \\ & + \left(\sin(\theta_1) \cos(\theta_2) \cos(\theta_3) - \cos(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) |001\rangle \\ & + \left(\cos(\theta_1) \sin(\theta_2) \cos(\theta_3) + \sin(\theta_1) \cos(\theta_2) \sin(\theta_3) \right) |111\rangle, \end{aligned}$$

and the final CNOT gate (with the control on the 2-qubit) gives

$$\begin{aligned} |\psi\rangle^{(\text{out})} &= \left(\cos(\theta_1) \cos(\theta_2) \cos(\theta_3) + \sin(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) |000\rangle \\ & + \left(\cos(\theta_1) \cos(\theta_2) \sin(\theta_3) - \sin(\theta_1) \sin(\theta_2) \cos(\theta_3) \right) |110\rangle \\ & + \left(\sin(\theta_1) \cos(\theta_2) \cos(\theta_3) - \cos(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) |101\rangle \\ & + \left(\cos(\theta_1) \sin(\theta_2) \cos(\theta_3) + \sin(\theta_1) \cos(\theta_2) \sin(\theta_3) \right) |011\rangle, \end{aligned}$$

3.1.2. *Input State* $|1\rangle$. In this case the initial state of the system is $|100\rangle$. As pointed out above, the preparatory phase of the circuit does not affect the value in the first qubit, so at the end of this phase our system is in the state

$$\begin{aligned} & \left(\cos(\theta_1) \cos(\theta_2) \cos(\theta_3) + \sin(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) |100\rangle \\ & + \left(\cos(\theta_1) \cos(\theta_2) \sin(\theta_3) - \sin(\theta_1) \sin(\theta_2) \cos(\theta_3) \right) |110\rangle \\ & + \left(\sin(\theta_1) \cos(\theta_2) \cos(\theta_3) - \cos(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) |101\rangle \\ & + \left(\cos(\theta_1) \sin(\theta_2) \cos(\theta_3) + \sin(\theta_1) \cos(\theta_2) \sin(\theta_3) \right) |111\rangle. \end{aligned}$$

We now consider the four CNOT gates at the end of the circuit, all of which must now be considered. The first CNOT gate yields

$$\begin{aligned} & \left(\cos(\theta_1) \cos(\theta_2) \cos(\theta_3) + \sin(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) |110\rangle \\ & + \left(\cos(\theta_1) \cos(\theta_2) \sin(\theta_3) - \sin(\theta_1) \sin(\theta_2) \cos(\theta_3) \right) |100\rangle \\ & + \left(\sin(\theta_1) \cos(\theta_2) \cos(\theta_3) - \cos(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) |111\rangle \\ & + \left(\cos(\theta_1) \sin(\theta_2) \cos(\theta_3) + \sin(\theta_1) \cos(\theta_2) \sin(\theta_3) \right) |101\rangle. \end{aligned}$$

The second CNOT affects the value in the 2-qubit, leaving us with

$$\begin{aligned} & \left(\cos(\theta_1) \cos(\theta_2) \cos(\theta_3) + \sin(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) |111\rangle \\ & + \left(\cos(\theta_1) \cos(\theta_2) \sin(\theta_3) - \sin(\theta_1) \sin(\theta_2) \cos(\theta_3) \right) |101\rangle \\ & + \left(\sin(\theta_1) \cos(\theta_2) \cos(\theta_3) - \cos(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) |110\rangle \\ & + \left(\cos(\theta_1) \sin(\theta_2) \cos(\theta_3) + \sin(\theta_1) \cos(\theta_2) \sin(\theta_3) \right) |100\rangle. \end{aligned}$$

The third CNOT gate has its control on the 1-qubit and target on the 0-qubit:

$$\begin{aligned} & \left(\cos(\theta_1) \cos(\theta_2) \cos(\theta_3) + \sin(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) |011\rangle \\ & + \left(\cos(\theta_1) \cos(\theta_2) \sin(\theta_3) - \sin(\theta_1) \sin(\theta_2) \cos(\theta_3) \right) |101\rangle \\ & + \left(\sin(\theta_1) \cos(\theta_2) \cos(\theta_3) - \cos(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) |010\rangle \\ & + \left(\cos(\theta_1) \sin(\theta_2) \cos(\theta_3) + \sin(\theta_1) \cos(\theta_2) \sin(\theta_3) \right) |100\rangle. \end{aligned}$$

Applying the final CNOT gate yields

$$\begin{aligned} |\psi\rangle^{(\text{out})} &= \left(\cos(\theta_1) \cos(\theta_2) \cos(\theta_3) + \sin(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) |111\rangle \\ &+ \left(\cos(\theta_1) \cos(\theta_2) \sin(\theta_3) - \sin(\theta_1) \sin(\theta_2) \cos(\theta_3) \right) |001\rangle \\ &+ \left(\sin(\theta_1) \cos(\theta_2) \cos(\theta_3) - \cos(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) |010\rangle \\ &+ \left(\cos(\theta_1) \sin(\theta_2) \cos(\theta_3) + \sin(\theta_1) \cos(\theta_2) \sin(\theta_3) \right) |100\rangle. \end{aligned}$$

3.1.3. *Output for Equatorial Qubits.* We can now use linearity to see that if

$$|\psi\rangle^{(\text{in})} = \frac{1}{\sqrt{2}} \left(|0\rangle + e^{i\phi} |1\rangle \right),$$

then

$$\begin{aligned}
|\psi\rangle^{(\text{out})} = & \frac{1}{\sqrt{2}} \left(\cos(\theta_1) \cos(\theta_2) \cos(\theta_3) + \sin(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) |000\rangle \\
& + \frac{1}{\sqrt{2}} \left(\cos(\theta_1) \cos(\theta_2) \sin(\theta_3) - \sin(\theta_1) \sin(\theta_2) \cos(\theta_3) \right) |110\rangle \\
& + \frac{1}{\sqrt{2}} \left(\sin(\theta_1) \cos(\theta_2) \cos(\theta_3) - \cos(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) |101\rangle \\
& + \frac{1}{\sqrt{2}} \left(\cos(\theta_1) \sin(\theta_2) \cos(\theta_3) + \sin(\theta_1) \cos(\theta_2) \sin(\theta_3) \right) |011\rangle \\
& + \frac{e^{i\phi}}{\sqrt{2}} \left(\cos(\theta_1) \cos(\theta_2) \cos(\theta_3) + \sin(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) |111\rangle \\
& + \frac{e^{i\phi}}{\sqrt{2}} \left(\cos(\theta_1) \cos(\theta_2) \sin(\theta_3) - \sin(\theta_1) \sin(\theta_2) \cos(\theta_3) \right) |001\rangle \\
& + \frac{e^{i\phi}}{\sqrt{2}} \left(\sin(\theta_1) \cos(\theta_2) \cos(\theta_3) - \cos(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) |010\rangle \\
& + \frac{e^{i\phi}}{\sqrt{2}} \left(\cos(\theta_1) \sin(\theta_2) \cos(\theta_3) + \sin(\theta_1) \cos(\theta_2) \sin(\theta_3) \right) |100\rangle.
\end{aligned} \tag{19}$$

3.2. Identifying Coefficients. We recall that the asymmetric phase-covariant cloning machine for equatorial qubits (meaning that $d = 2$) has the form

$$\begin{aligned}
|0\rangle & \mapsto \mu |000\rangle + \nu |011\rangle + \xi |101\rangle \\
|1\rangle & \mapsto \mu |111\rangle + \nu |100\rangle + \xi |010\rangle.
\end{aligned}$$

By comparing coefficients with those from (19) we see that

$$\begin{aligned}
\mu &= \cos(\theta_1) \cos(\theta_2) \sin(\theta_3) - \sin(\theta_1) \sin(\theta_2) \sin(\theta_3) \\
\nu &= \cos(\theta_1) \sin(\theta_2) \cos(\theta_3) + \sin(\theta_1) \cos(\theta_2) \sin(\theta_3) \\
\xi &= \sin(\theta_1) \cos(\theta_2) \cos(\theta_3) - \cos(\theta_1) \sin(\theta_2) \sin(\theta_3),
\end{aligned} \tag{20}$$

together with the requirement that

$$\cos(\theta_1) \cos(\theta_2) \sin(\theta_3) - \sin(\theta_1) \sin(\theta_2) \cos(\theta_3) = 0. \tag{21}$$

3.3. Parameter Selection. We now turn to the problem of determining values for $\theta_1, \theta_2, \theta_3$ that realize the optimal cloning described in Section 2.1.

To begin we recall that the factors η_A and η_B are given by

$$\eta_A = 2\mu\nu \quad \text{and} \quad \eta_B = 2\mu\xi$$

and satisfy the circle identity

$$\eta_A^2 + \eta_B^2 = 1.$$

Using the identifications (20) we are able to express η_A and η_B in terms of θ_1, θ_2 , and θ_3 . Indeed,

$$\begin{aligned}
\eta_A &= 2\mu\nu \\
&= 2 \left(\cos(\theta_1) \cos(\theta_2) \cos(\theta_3) + \sin(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) \left(\cos(\theta_1) \sin(\theta_2) \cos(\theta_3) + \sin(\theta_1) \cos(\theta_2) \sin(\theta_3) \right) \\
&= 2 \left(\cos^2(\theta_1) \sin(\theta_2) \cos(\theta_2) \cos^2(\theta_3) + \sin(\theta_1) \cos(\theta_1) \cos^2(\theta_2) \sin(\theta_3) \cos(\theta_3) \right. \\
&\quad \left. + \sin(\theta_1) \cos(\theta_1) \sin^2(\theta_2) \sin(\theta_3) \cos(\theta_3) + \sin^2(\theta_1) \sin(\theta_2) \cos(\theta_2) \sin^2(\theta_3) \right)
\end{aligned}$$

Observe that we can use the identity (21) to rewrite the two middle terms in this expression:

$$\begin{aligned}\sin(\theta_1) \cos(\theta_1) \cos^2(\theta_2) \sin(\theta_3) \cos(\theta_3) &= \sin(\theta_1) \cos(\theta_2) \cos(\theta_3) \left(\cos(\theta_1) \cos(\theta_2) \sin(\theta_3) \right) \\ &= \sin(\theta_1) \cos(\theta_2) \cos(\theta_3) \left(\sin(\theta_1) \sin(\theta_2) \cos(\theta_3) \right) \\ &= \sin^2(\theta_1) \sin(\theta_2) \cos(\theta_2) \cos^2(\theta_3),\end{aligned}$$

and

$$\begin{aligned}\sin(\theta_1) \cos(\theta_1) \sin^2(\theta_2) \sin(\theta_3) \cos(\theta_3) &= \cos(\theta_1) \sin(\theta_2) \sin(\theta_3) \left(\sin(\theta_1) \sin(\theta_2) \cos(\theta_3) \right) \\ &= \cos(\theta_1) \sin(\theta_2) \sin(\theta_3) \left(\cos(\theta_1) \cos(\theta_2) \sin(\theta_3) \right) \\ &= \cos^2(\theta_1) \sin(\theta_2) \cos(\theta_2) \sin^2(\theta_3).\end{aligned}$$

Returning to our calculation of η_A , we have

$$\begin{aligned}\eta_A &= 2 \left(\cos^2(\theta_1) \sin(\theta_2) \cos(\theta_2) \cos^2(\theta_3) + \sin^2(\theta_1) \sin(\theta_2) \cos(\theta_2) \cos^2(\theta_3) \right. \\ &\quad \left. + \cos^2(\theta_1) \sin(\theta_2) \cos(\theta_2) \sin^2(\theta_3) + \sin^2(\theta_1) \sin(\theta_2) \cos(\theta_2) \sin^2(\theta_3) \right) \\ &= 2 \left(\sin(\theta_2) \cos(\theta_2) \cos^2(\theta_3) + \sin(\theta_2) \cos(\theta_2) \sin^2(\theta_3) \right) \\ &= 2 \sin(\theta_2) \cos(\theta_2) \\ &= \sin(2\theta_2).\end{aligned}$$

In the case of the η_B factor we proceed similarly. We calculate

$$\begin{aligned}\eta_B &= 2\mu\xi \\ &= 2 \left(\cos(\theta_1) \cos(\theta_2) \cos(\theta_3) + \sin(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) \left(\sin(\theta_1) \cos(\theta_2) \cos(\theta_3) - \cos(\theta_1) \sin(\theta_2) \sin(\theta_3) \right) \\ &= 2 \left(\sin(\theta_1) \cos(\theta_1) \cos^2(\theta_2) \cos^2(\theta_3) - \cos^2(\theta_1) \sin(\theta_2) \cos(\theta_2) \sin(\theta_3) \cos(\theta_3) \right. \\ &\quad \left. + \sin^2(\theta_1) \sin(\theta_2) \cos(\theta_2) \sin(\theta_3) \cos(\theta_3) - \sin(\theta_1) \cos(\theta_1) \sin^2(\theta_2) \sin^2(\theta_3) \right)\end{aligned}$$

Using (21) we simplify the second and third terms in our sum, as we did in the case of η_A above:

$$\begin{aligned}\cos^2(\theta_1) \sin(\theta_2) \cos(\theta_2) \sin(\theta_3) \cos(\theta_3) &= \cos(\theta_1) \sin(\theta_2) \cos(\theta_3) \left(\cos(\theta_1) \cos(\theta_2) \sin(\theta_3) \right) \\ &= \cos(\theta_1) \sin(\theta_2) \cos(\theta_3) \left(\sin(\theta_1) \sin(\theta_2) \cos(\theta_3) \right) \\ &= \sin(\theta_1) \cos(\theta_1) \sin^2(\theta_2) \cos^2(\theta_3),\end{aligned}$$

and

$$\begin{aligned}\sin^2(\theta_1) \sin(\theta_2) \cos(\theta_2) \sin(\theta_3) \cos(\theta_3) &= \sin(\theta_1) \cos(\theta_2) \sin(\theta_3) \left(\sin(\theta_1) \sin(\theta_2) \cos(\theta_3) \right) \\ &= \sin(\theta_1) \cos(\theta_2) \sin(\theta_3) \left(\cos(\theta_1) \cos(\theta_2) \sin(\theta_3) \right) \\ &= \sin(\theta_1) \cos(\theta_1) \cos^2(\theta_2) \sin^2(\theta_3)\end{aligned}$$

Thus

$$\begin{aligned}\eta_B &= 2 \left(\sin(\theta_1) \cos(\theta_1) \cos^2(\theta_2) \cos^2(\theta_3) - \sin(\theta_1) \cos(\theta_1) \sin^2(\theta_2) \cos^2(\theta_3) \right. \\ &\quad \left. + \sin(\theta_1) \cos(\theta_1) \cos^2(\theta_2) \sin^2(\theta_3) - \sin(\theta_1) \cos(\theta_1) \sin^2(\theta_2) \sin^2(\theta_3) \right) \\ &= 2 \left(\sin(\theta_1) \cos(\theta_1) \cos^2(\theta_2) - \sin(\theta_1) \cos(\theta_1) \sin^2(\theta_2) \right) \\ &= \sin(2\theta_1) \cos(2\theta_2).\end{aligned}$$

In case the cloning transformation is optimal, η_A and η_B satisfy the circle relation

$$\eta_A^2 + \eta_B^2 = 1.$$

In terms of our angles $\theta_1, \theta_2, \theta_3$, this now reads

$$\sin^2(2\theta_2) + \sin^2(2\theta_1) \cos^2(2\theta_2) = 1.$$

In light of the standard Pythagorean identity, $\cos^2(2\theta_2) + \sin^2(2\theta_2) = 1$, we find that

$$\sin^2(2\theta_1) = 1 \quad \text{or} \quad \cos^2(2\theta_2) = 0.$$

Notice that if $\cos(2\theta_2) = 0$, then $\eta_B = 0$ and $\eta_A = 1$. While this presents itself as a possibility, it isn't a particularly interesting one and we will not dwell on it. Of greater interest is the other restriction which requires $\sin(2\theta_1) = \pm 1$.

- If $\sin(2\theta_1) = 1$, then $\theta_1 = \pi/4$ (up to integer multiples of π). Returning this to (21) yields

$$\cos(\theta_2) \sin(\theta_3) - \sin(\theta_2) \cos(\theta_3) = 0.$$

That is, we require that $\sin(\theta_3 - \theta_2) = 0$, meaning that $\theta_3 - \theta_2 = k\pi$ for some $k \in \mathbf{Z}$. In what follows we will take $k = 0$, yielding $\theta_2 = \theta_3$. This choice is also consistent with the requirement that

$$\mu = \cos(\theta_1) \cos(\theta_2) \cos(\theta_3) + \sin(\theta_1) \sin(\theta_2) \sin(\theta_3) = \frac{1}{\sqrt{2}},$$

which emerges as a requirement of the optimal transformation.

- If $\sin(2\theta_1) = -1$, then $\theta_1 = 3\pi/4$ (again up to an integer multiple of π). Returning this to (21) leaves us with

$$-\cos(\theta_2) \sin(\theta_3) - \sin(\theta_2) \cos(\theta_3) = 0,$$

which simplifies to $\sin(\theta_3 + \theta_2) = 0$, meaning that $\theta_2 + \theta_3 = k\pi$ for some $k \in \mathbf{Z}$. **I haven't thought much about this case. It is still consistent with the requirement that $\mu = 1/\sqrt{2}$.**

Summary. The circuit presented at the beginning of this section implements a phase-covariant cloning machine if the angles $\theta_1, \theta_2, \theta_3$ are selected so that

$$\theta_1 = \frac{\pi}{4}, \quad \theta_2 = \theta_3.$$

There is a requirement that $\eta_A, \eta_B \in [0, 1]$, which may restrict our choices, but I've not thought further about this.

4. SEPARABILITY

The chief objective of RSG in [11] is to show that the output states of the $1 \rightarrow 2$ phase-covariant cloner are separable in the optimal cases. This means that if we choose our cloner to be optimal, then the output state $\rho_{AB}^{(\text{out})}$ (where we compute the partial trace over the ancilla) is separable.

A state (described by a density operator) ρ defined on a composite system whose state space is given by $\mathcal{H}_A \otimes \mathcal{H}_B$ is said to be separable if it can be written as

$$\rho = \sum_{i=1}^k c_i \rho_i^A \otimes \rho_i^B \tag{22} \quad \{\text{separable_st.}\}$$

where $c_i \in \mathbf{C}$ and ρ_i^A, ρ_i^B are density operators for the A and B systems, respectively. (Note that $\rho_i^A \neq \rho_i^B$ in general.) The decomposition (22) is referred to as the Schmidt decomposition of the state ρ . It is, in general, a very difficult problem to determine if a state is separable or entangled. There is a criterion developed in Peres [10] and in Horodecki, Horodecki, and Horodecki in [6] called the positive partial trace condition. This is the condition that RSG use to determine the separability of states in their paper. Unfortunately, the positive partial trace is only sufficient for this purpose in the case of qubits (where $d = 2$), so RSG can only guarantee the separability of the outputs in the qubit case.

Open Problem. It appears that the separability of output states in the case of optimal phase-covariant clones is unknown in high dimensions ($d > 2$).

4.1. Recent results on separability. There are some more recent results on separability. In particular, [5] appears to present a complete classification of entanglement states in d -dimensional systems. The tools that are involved are largely algebro-geometric in nature.

Partial Trace Calculations

Term	tr_{BX}	tr_{AX}
$\sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* iii\rangle \langle jjj $	$\sum_{i=0}^{d-1} \alpha_i ^2 i\rangle \langle i $	$\sum_{i=0}^{d-1} \alpha_i ^2 i\rangle \langle i $
$\sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} iii\rangle \langle jll $	$\sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* i\rangle \langle j $	0
$\sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} iii\rangle \langle ljl $	0	$\sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* i\rangle \langle j $
$\sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} ikk\rangle \langle jjj $	$\sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* i\rangle \langle j $	0
$\sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} ikk\rangle \langle jll $	$(d-1) \sum_{i=0}^{d-1} \alpha_i ^2 i\rangle \langle i $ $+ (d-2) \sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* i\rangle \langle j $	$1 - \sum_{i=0}^{d-1} \alpha_i ^2 i\rangle \langle i $
$\sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} ikk\rangle \langle ljl $	0	0
$\sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} kik\rangle \langle jjj $	0	$\sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* i\rangle \langle j $
$\sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} kik\rangle \langle jll $	0	0
$\sum_{i=0}^{d-1} \sum_{j=0}^{d-1} \alpha_i \alpha_j^* \sum_{\substack{0 \leq k \leq d-1 \\ k \neq i}} \sum_{\substack{0 \leq l \leq d-1 \\ l \neq j}} kik\rangle \langle ljl $	$1 - \sum_{i=0}^{d-1} \alpha_i ^2 i\rangle \langle i $	$(d-1) \sum_{i=0}^{d-1} \alpha_i ^2 i\rangle \langle i $ $+ (d-2) \sum_{i=0}^{d-1} \sum_{\substack{0 \leq j \leq d-1 \\ j \neq i}} \alpha_i \alpha_j^* i\rangle \langle j $

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