# Quantum Computation

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July 2021

## 1 Lecture 1

# 1.1 Review of Shor's algorithm/quantum period finding algorithm

#### Polynomial time hierarchy:

Computation with input of size n, and we are interested in the number of steps/gates (classical or quantum). When we say O(poly(n)) steps we regard this as an "efficient computation".

Shor's algorithm solves the factoring problem:

Given an integer N needing O(log N) bits, we want to find a non-trival factor in O(polyn) time.

The best known classical algorithm (number sieve):  $e^{O(n^{\frac{1}{3}}(\log n)^{\frac{1}{3}})}$ Shor's alogrithm takes  $O(n^3)$ 

#### 1.1.1 Quantum factoring algorithm (summary)

1. First, convert factoring into periodicity determination. Given N, choose a < N s.t. a is coprime (this is done by picking a at random and then using Euclids algorithm to compute the greatest common divisor b = gcd(a, N), if b > 1 we have found a factor else they are coprime). Consider  $f: \mathbb{Z} \to \mathbb{Z}_N$   $f(x) = a^x \mod N$ .

**Euler's Theorem:** If a and N are coprime then there is a least power 1 < r < N such that  $a^r = 1 \mod N$ . We call r the order a mod N. So f(r) = 1 so f(k+r) = f(k) so f is periodic with period r, with f one-to-one within each period.

2. In order to find r we need a quantum implementation of f. We are always working on finite size registers so restricting  $x \in \mathbb{Z}$  to  $x \in \mathbb{Z}_M$  (for some large enough M):  $f: \mathbb{Z}_M \to \mathbb{Z}_N$ . f will no longer be exactly periodic but this would have negligible effect if M is sufficiently large e.g.  $M = O(N^2)$ 

- 3. Using the classical theory of continued fractions. Define Hilbert spaces  $\mathcal{H}_M \to \{|i\rangle\}_{i\in\mathbb{Z}_M}, \mathcal{H}_N \to \{|i\rangle\}_{i\in\mathbb{Z}_N}.$
- 4.  $|x\rangle \rightarrow |f(x)\rangle$  is not generally a valid quantum operator, so we make it a unitary operation which can be implemented:

$$U_f: \mathcal{H}_M \otimes \mathcal{H}_N \to \mathcal{H}_M \otimes \mathcal{H}_N$$

$$U_f: |i\rangle |k\rangle \rightarrow |i\rangle |k+f(i)\rangle$$

- 5. If  $x \to f(x)$  can be classically computed in O(poly(m)) time  $(m = \log M)$ , then  $U_f$  can be implemented in poly(m) time quantumly too
- 6. We will sometimes view  $U_f$  as a black box/oracle and we will count the number of times the algorithm invokes the oracle.
- 7. Back to factoring to get r we'll use the quantum algorithm for periodicity determination:
- 8. Given an oracle  $U_f$  with the promise that f is periodic of some unknown period  $r \in \mathbb{Z}_N$  so that f(x+r) = f(x) and f is one-to-one in this period (for all  $0 \le x_1 < x_2 < rf(x_1) \ne f(x_2)$ ). To find r in O(poly n) with any prescribed success probability  $1 \epsilon$  we use the following algorithm:
  - Step 1: Create the state

$$\frac{1}{\sqrt{M}} \sum_{i=0}^{M-1} |i\rangle |0\rangle$$

• Step 2: Apply  $U_f$  to get

$$\frac{1}{\sqrt{M}} \sum_{i=0}^{M-1} |i\rangle |f(i)\rangle$$

• Step 3: Measure the 2nd register to get y. By the Born rule the first register collapses to all those i: f(i) = y i.e.  $i = x_0, x_0 + r, x_o + 2r, ..., x_0 + (A-1)r, 0 \le x_0 < r$ .

Discard the second register to get the following state:

$$|per\rangle = \frac{1}{\sqrt{A}} \sum_{j=0}^{A-1} |x_0 + jr\rangle$$

If we measure  $|per\rangle$  in computation basis we will get a value of one of these states  $x_0 + jr$  for uniformly random j. This only gives us a random element of  $\mathbb{Z}_M$  with no information about r.

• Step 4: Apply quantum fourier transform mod M (QFT). Lets recap what QFT does:

$$|x\rangle \to \frac{1}{\sqrt{M}} \sum_{y=0}^{M-1} \omega^{xy} |y\rangle, \forall x \in \mathbb{Z}_M, \omega = e^{2\pi i/M}$$

This can be implement in  $O(m^2)$  time and gives state:

$$QFT |per\rangle = \frac{1}{\sqrt{MA}} \sum_{j=0}^{A-1} \sum_{y=0}^{M-1} \omega^{(x_0 + jr)y} |y\rangle = \frac{1}{\sqrt{MA}} \sum_{y=0}^{M-1} \omega^{x_0 y} \left[ \sum_{j=0}^{A-1} \omega^{jry} |y\rangle \right]$$

The square brackets will be:

$$\begin{cases} A & \text{if } y = KA = k\frac{M}{r}, x = 0, 1, ..., r-1 \\ 0 & \text{otherwise} \end{cases}$$

So gives final state:

$$QFT \left| per \right\rangle = \sqrt{\frac{A}{M}} \sum_{k=0}^{A-1} \omega^{x_0 k \frac{N}{r}} \left| k \frac{M}{r} \right\rangle$$

Now the random shift  $x_0$  only appears in the phase not in the ket labels. So now the measurement probabilities will be independent of  $x_0$ . When we measure this we get some value  $c = \frac{k_0 M}{r}$  with  $k_0$  uniformly random in range  $0 \le k_0 < r$ 

$$\frac{k_0}{r} = \frac{c}{M}$$

As c and M are known, and  $k_0$  is unknown but random in the given range. We want to find r and so we recall several classical facts.

**Co-primality Theorem**: The number of integers less than r that are coprime to r grows with  $O(\frac{r}{\log \log r})$ 

Therefore, the probability of  $k_0$  being coprime to r is  $O(\frac{1}{\log \log r})$ .

**Lemma**: If a single trial has success probability p then if one repeats it  $M^*$  times, for any  $0 < 1 - \epsilon < 1$ . We get probability of at least one success in  $M^*$  trails is greater than  $1 - \epsilon$  if  $M^* = \frac{-\log \epsilon}{p}$ . i.e. roughly O(1/p) trials suffice to achieve probability of success  $> 1 - \epsilon$ 

• After step 4 cancel  $\frac{c}{M}$  down to an irreducible algorithm  $\frac{a}{b}$  there is an efficient algorithm (O(poly n)) for this. This will give us r as denominator b if  $k_0$  is coprime to r with probability  $O(\frac{1}{\log \log r})$ . So check b value by computing f(0) and f(b) and  $b = r \iff f(0) = f(b)$ .

By repeating this process  $M^* = O(\log \log r)$  times this will give us r with any desired probability  $1 - \epsilon$ . Since r < M the whole algorithm takes O(poly m) time!

9. From learning the period r we can use number theory to find a factor of N

#### 1.1.2 Further insights to QFT

Now lets think about the implications of QFT. What does applying quantum fourier transform really achieve?

Lets consider a function:  $f: \mathbb{Z}_M \to \mathbb{Z}_N$  with period  $r \in \mathbb{Z}_M$ ,  $A = \frac{M}{r}$ . Define:

$$R = \{0, r, 2r, 3r, \dots, (A-1)r\} < \mathbb{Z}_M$$
$$|R\rangle = \frac{1}{\sqrt{A}} \sum_{k=0}^{A-1} |kr\rangle$$
$$|per\rangle = |x_0 + R\rangle = \frac{1}{\sqrt{A}} \sum_{k=0}^{A-1} |x_0 + rk\rangle$$

The problem was this random shift  $x_0$  when measuring  $|per\rangle$ . For each  $x_0 < \mathbb{Z}_M$  consider a mapping  $k \to k + x_0$ . "Shift by  $x_0$ ". It is a 1-1 invertible map, and can define a unitary version  $U(x_0)$  on  $\mathcal{H}_M$ :  $U(x_0)|k\rangle = |k + x_0\rangle$ .

$$|x_0 + R\rangle = U(x_0)|R\rangle$$

Since  $(\mathbb{Z}_M, +)$  is an abelian group  $U(x_0)U(x_1) = U(x_0 + x_1) = U(x_1)U(x_0)$ . So all  $U(x_i)$  commute as operators on  $\mathcal{H}_M$ . Therefore they have an orthonomal basis of common eigenvectors  $\{|\chi_k\rangle\}_{k\in\mathbb{Z}_M}$ . These are called shift invariant states as  $U(x_0)|\chi_k\rangle = \omega(x_0,k)|\chi_k\rangle$  for all  $x_0,k\in\mathbb{Z}_M$  with the important caveat that  $|\omega(x_0,k)| = 1$ .

Consider  $|R\rangle$  written in  $\{|\chi_r\rangle\}$  basis:

$$|R\rangle = \sum_{k=0}^{M-1} a_k |\chi_k\rangle$$

 $a_k$  only depend on r not on  $x_0$ . Then:

$$|per\rangle = U(x_0) |R\rangle = \sum_{k=0}^{M-1} a_k \omega(x_0, k) |\chi_k\rangle$$

Here it can be seen that the probability of measuring k is

$$prob(k) = |a_k \omega(x_0, k)|^2 = |a_k|^2$$

So this is all indepedant of  $x_0$  and depends only on r. So measuring in this basis gives us some information about r. So one can think of QFT as the unitary mapping that rotates  $\chi$  basis into the standard computational basis. So can define QFT as:

$$QFT |\chi_k\rangle = |k\rangle$$

How do these mysterious shift invariant states look?

## 1.1.3 Explicit form of shift invariant shapes

$$|\chi_k\rangle = \frac{1}{\sqrt{M}} \sum_{l=0}^{M-1} e^{-2\pi i l \frac{k}{M}} |l\rangle$$

$$U(x_0) |\chi_k\rangle = 1\sqrt{M} \sum_{l=0}^{M-1} e^{-2\pi i l \frac{k}{M}} |l+x_0\rangle = 1\sqrt{M} \sum_{\tilde{l}=0}^{M-1} e^{-2\pi i (\tilde{l}-x_0) \frac{k}{M}} |\tilde{l}\rangle = e^{2\pi i k \frac{x_0}{M}} |\chi_k\rangle$$

giving eigenvalue:  $\omega(x_0, k) = e^{2\pi i k \frac{x_0}{M}}$ . From this we could reconstruct the basis of QFT:

$$[QFT]_{kl} = \frac{1}{\sqrt{M}} e^{2\pi i \frac{kl}{M}}$$

# 2 Lecture 3

# 2.1 Hidden Subgroup Problem

Let G be a finite group of size |G|. We are given an oracle  $f: G \to X$  with X just some set. We are promised there is a subgroup K < G s.t.

f is constant on (left) cosets of K in G

f is distinct on distinct cosets

**Problem**: 'Determine' the 'hidden subgroup' K (e.g. output a set of generators or sample uniformly from elements of K)

We want to solve in time  $O(poly(\log |G|))$  (efficient algorithm) with any consistent probability  $1 - \epsilon$ .

# Examples of problems that can be cast as HSP

Periodicity finding  $f: \mathbb{Z}_M \to X$  periodic, period r 1-1 in period

$$G = \mathbb{Z}_M, K = \{0, r, 2r, ..., (A-1)r\} < G$$

Discrete Logarithm Problem: p - prime number,  $\mathbb{Z}_p^*$  group of integers with multiplication mod  $p, g \in \mathbb{Z}_p^*$  to be a generator (or primitive root mod p). If  $\mathbb{Z}_p^* = \{g^0, g^1, ..., g^{r-2}\}$  and we have  $g^{p-1} = 1 \pmod{p}$ . Fact: These always exist for p is prime. Any  $x \in \mathbb{Z}_p^*$  can be written as  $x = g^y$  for some  $y \in \mathbb{Z}_{p-1}$ ,  $y = \log_q x$  is called the discrete  $\log$  of x to base g. Discrete  $\log$  problem is given

a generator  $g, x \in \mathbb{Z}_p^*$  we want to compute  $y = \log_g x$ . To express this as the HSP:

$$f: \mathbb{Z}_{p-1} \times \mathbb{Z}_{p-1} \to \mathbb{Z}_p^*$$

$$f(a,b) = q^a x^{-b} \mod p = q^{a-yb} \mod p$$

Can check if  $f(a_1, b_1) = f(a_2, b_2) \iff (a_1, b_1) = (a_2, b_2) + \lambda(y, 1) \in \mathbb{Z}_{p-1}$ :

$$G = \mathbb{Z}_{p-1} \times \mathbb{Z}_{p-1}$$

$$K = \{\lambda(y, 1) : \lambda \in \mathbb{Z}_{p-1}\} < G$$

Then f is constant and distinct on cosets of K and generator (y, 1) of K gives  $y = \log_q x$ 

#### Graph Problems:

So we can solve problems like those above were G is abelian, but we can also solve graph problems.

Consider graph A = (V, E), |V| = n lets say that the graph is undirected and there is at most one edge between any two vertices. Vertices here are labelled by numbers from 1 to n.

Lets define an adjacency matrix  $M_A$ :  $[M_a]_{ij} = \begin{cases} 1 & \iff (i,j) \\ 0 & \text{otherwise} \end{cases}$ . The per-

mutation group of [n],  $|P_n| = n!$ ,  $\log |P_n| \sim O(n \log n)$ . Define a group of automorphisms of group A which is a set of permutations with the following property:  $\pi \in P_n$  s.t.  $\forall i, j(i,j)$  is an edge in  $A \iff (\pi(i), \pi(j))$  is also an edge in A.

An associated HSP (the case of non-abelian G):

$$G = P_n, X = \text{set of all labelled graphs on } n \text{vertices}$$

For any  $A \in X$ , define  $f_A : G \to X$ ,  $f_A(\pi) =$ "A with vertex labels permuted by  $\pi$ "

$$K = Aut(A)$$

(Check f(K) is constant and discrint on cosets of Aut(A))

#### **Applications**:

If we can sample uniformly from K, then we can solve Graph Isomorphism problem (GI). This has a number of different applications in areas of computer science. Two labelled graphs A and B with n vertices are isomorphic if there is a 1-1 map (i.e. permutation)  $\pi[n] \to [n]$  s.t.  $\forall i, j \in [n](i,j)$  is an edge in  $A \Leftrightarrow (\pi(i), \pi(j))$  is an edge in B. The GI problem is given to graphs A and B and deciding if they are isomorphic. This can be represented as a non-abealian HSP. There is no known poly(m) time classical algorithm to solve this problem,

so GI is clearly in NP but not believed to be NP-complete (a class of problems such that every problem in NP can be reduced to an NP-complete problem these are the hardest NP problems). In 2017, L Babai presented a quasi-polynomial algorithm for GI runtime  $n^{O((\log n)^2)}$ . This ranks in between polynomial runtime and exponential algorithms.

# 3 Lecture 4

Quantum algorithm for finite abelian  $\operatorname{HSPs}$  - Generalisation of period-finding algorithm

Write our abelian group (G, +) additively

Construction of shift-invariant states and Fourier transform for G.

Representations of abelian G:

Consider the mapping  $\chi: G \to \mathbb{C}^* = \mathbb{C} - \{0\}$  with multiplication that satisfies:

$$\chi(g_1 + g_2) = \chi(g_1)\chi(g_2), \forall g_1, g_2 \in G$$

 $\chi$  is a group homomorphism from G to  $\mathbb{C}^*$ . Such  $\chi$ 's are called irreducible representations of G. They have the following properties: **Theorem 1**:

- 1) any value  $\chi(g)$  is a |G|-th root of unity ( $\chi \in \to S^1$  the unit circle)
- 2) Schur's lemma (orthogonality): If  $\chi_i, \chi_j$  satisfy (HOM) then

$$\frac{1}{|G|} \sum_{g \in G} \chi_i(g) \bar{\chi_j}(g) = \delta_{ij}$$

3) There are always exactly |G| different functions  $\chi$  satisfying (HOM).

**Examples**:  $\chi(g) = 1, \forall g \in G$  is an irrep/ called a trivial irrep Label the trivial irrep as  $\chi_0$ ,  $0 \in G$ . Then for any other irrep  $\chi \neq \chi_0$  orthogonality to  $\chi_0$  gives:

$$\sum_{g \in G} \chi(g) = 0 \text{ if } \chi \neq \chi_0$$

Going back to constructing shift-invariant states

#### 3.0.1 Shift-invariant states

Consider a state space  $\mathcal{H}_G$ ,  $dim\mathcal{H}_G = |G|$  with basis  $\{|g\rangle\}_{g\in G}$ . Now introduce shift operators U(k) for  $k \in G$  defined as follows:

$$U(k):\left|g\right> 
ightarrow \left|g+k\right>,g,k\in G$$

All shift operators commute so there exists a simultaneous eigenbasis. For each  $\chi_k, k \in G$ :

$$|\chi_k\rangle = \frac{1}{\sqrt{|G|}} \sum_{g \in G} \bar{\chi_k}(g) |g\rangle$$

By theorem 1  $\{\chi_k\}$  form an orthonormal basis.

$$U(g) |\chi_k\rangle = \chi_k(g) |\chi_k\rangle$$

**Proof**:

$$U(g) |\chi_k\rangle = \frac{1}{\sqrt{|G|}} \sum_{h \in G} \bar{\chi}_k(h) |h + g\rangle$$

$$\{h' = h + g\} = \frac{1}{\sqrt{|G|}} \sum_{h' \in G} \bar{\chi_k} (h' - g) |h'\rangle$$

using HOM  $\chi_k(-g) = \chi_k(g)^{-1} = \bar{\chi}_k(g) \implies \overline{\chi_k(h'-g)} = \bar{\chi}_k(h')\bar{\chi}_k(-g) = \bar{\chi}_k(g)$  $\bar{\chi}_k(h')\chi_k(g)$ . Therefore,

$$U(g) |\chi_k\rangle = \frac{1}{\sqrt{|G|}} \sum_{h' \in G} \chi_k(g) \bar{\chi}_k(h') |h'\rangle = \chi_k(g) |\chi_k\rangle$$

So  $|\chi_k\rangle$ 's form a common eigenbasis

Introduce Fourier transform QFT for a group G

- consider a unitary mapping on  $\mathcal{H}_G$  mapping  $|\chi_k\rangle$  basis to  $|g\rangle$  basis

$$QFT |\chi_q\rangle = |q\rangle, \forall q \in G$$

$$QFT^{-1}|g\rangle = |\chi_g\rangle$$

k-th column of  $QFT^{-1}$  in  $|g\rangle$  basis is mode of components of  $|\chi_k\rangle$ :

$$[QFT^{-1}]_{gk} = \frac{1}{\sqrt{|G|}}\bar{\chi_k}(g)$$

**Example**:  $G = \mathbb{Z}_M L$ Check  $\chi_a(b) = e^{\frac{2\pi i a b}{M}}, a, b \in \mathbb{Z}_M$  satisfies HOM and has its irreps labelled by  $a \in \mathbb{Z}_M$  with  $\chi_0(b) = 1 \forall b \in \mathbb{Z}_m$ .

$$G = \mathbb{Z}_{M_1} \times ... \times \mathbb{Z}_{M_r}$$

$$(a_1, ..., a_r) = g_1, (b_1, ..., b_r) = g_2$$

$$\chi_{q_1}(q_2) = e^{2\pi i(\frac{a_1b_1}{M_1} + \dots + \frac{a_rb_r}{M_r})}$$

This satisfies HOM and our  $QFT_G = QFT_{M_1} \otimes ... \otimes QFT_{M_r}$  on  $\mathcal{H}_G = \mathcal{H}_{M_1} \otimes ... \otimes QFT_{M_r}$  $...\otimes \mathcal{H}_{M_r}.$ 

This second example is exhaustive since we have a classification theorem:

Classification theorem: Any finite abelian group G is isomorphic to a direct product of the form  $G = \mathbb{Z}_{M_1} \otimes ... \otimes \mathbb{Z}_{M_r}$ . So  $M_1$  can be taken in a form  $p_1^{s_1}, ... p_r^{s_r}$ .

#### 3.0.2 Quantum algorithm

$$f:G\to X$$

with hidden subgroup K and cosets  $k=0+k, g_2+k, ... g_m+k, m=\frac{|G|}{|K|}$ . we will work on  $\mathcal{H}_{|G|}\otimes\mathcal{H}_{|X|},\{|g\rangle|x\rangle\}_{g\in G,x\in X}$ .

Create a state  $\frac{1}{\sqrt{|G|}} \sum_{g \in G} \left| g \right\rangle \left| 0 \right\rangle$ 

Apply  $U_f$  and  $\frac{1}{\sqrt{|G|}} \sum_{g \in G} |g\rangle |f(g)\rangle$ 

 $Measure the second register to get f(g_0)$ . The first register will not give the coset state:

$$|g_0 + k\rangle = \frac{1}{\sqrt{|k|}} \sum_{k \in K} |g_0 + k\rangle = U(g_0) |K\rangle$$

apply QFT and measure to get a result  $g \in G$ 

# 4 Lecture 6

We can write  $|K\rangle$  in the shift-invariant basis  $\{\chi_q\}_{q\in G}$ 

$$|K\rangle = \sum_{g} a_g |\chi_g\rangle$$

$$|g_0 + K\rangle = U(g_0) |K\rangle = \sum_g a_g \chi_g(g_0) |\chi_g\rangle$$

as  $QFT |\chi_g\rangle = |g\rangle$  so after we apply QFT

$$prob(g) = |a_g \chi_g(g_0)|^2 = |a_g|^2, |\chi_g(g_0)| = 1$$

$$QFT |K\rangle = \frac{1}{\sqrt{|G|}} \frac{1}{\sqrt{|K|}} \sum_{l \in G} (\sum_{k \in K} \chi_l(k) |l\rangle)$$

 $\sum_{k \in K} \chi_l(k) |l\rangle$  involves irreps  $\chi_l$  of G restricted to subgroup K < G, and each such object is itself an irrep in K. Hence we have the following relation:

$$\sum_{k \in K} \chi_l(k) = \begin{cases} |k| & if \chi_l \text{ restricts to the trival irrep of K} \\ 0 & \text{otherwise} \end{cases}$$

$$QFT|K\rangle = \sqrt{\frac{|K|}{|G|}} \sum_{l \in G} |l\rangle$$

Then a measurement gives a uniformly random choice of l s.t.  $\chi_l(k) = 1$ . If k has generators  $k_1, ..., k_n$  where  $M = O(\log(K)) = O(\log|G|)$ . Then the output of a measurement gives us  $\chi_l(k) = 1 \forall i$ . It can be shown that if  $O(\log(|G|))$  values of l chosen uniformly at random then with probability  $> \frac{2}{3}$  they will suffice to determine a generating set for kvia the equations  $\chi_l(k) = 1$ .

Example:  $G = \mathbb{Z}_{M_1} \times ... \times \mathbb{Z}_{M_l}$ 

 $l=(l_1,...,l_q)\in G,\,g=(b_1,...,b_q)\in G \text{ gives } \chi_l(g)=e^{2\pi i(\frac{l_1b_1}{M_1}+...+\frac{l_qb_q}{M_q}}$  For  $k=(k_1,...,k_q)\in K$  with  $\chi_l(k)=1 \implies \frac{l_1k_1}{M_1}+...+\frac{l_qk_q}{M_q}=0 mod 1.$  This is a homogenous linear equation on k and  $O(\log(k))$  such equations determine kas null space.

#### 4.0.1 Remarks on HSP for non-abelian groups G

Now we will consider multiplicative shifts. As before we can generate a bunch of coset states but it is curious to investigate what breaks down.

$$|g_0K\rangle = \frac{1}{\sqrt{|K|}} \sum_{k \in K} |g_0k\rangle, g_0 \in G$$
 is chosen randomly

The real problem with QFT construction is that there is no good basis of shift invariant states. This is because  $U(q_0)$  don't commute.

#### Construction of non-abelian QFT

Consider a d-dimensional representation of G and a group homomorphism  $\chi$ :

 $\chi$  is irreducible if no subset of  $\mathbb{C}^d$  is left invariant by all matrices  $\chi(g), g \in G$ . (i.e. we cannot simulatenously block-diagonalize all of  $\chi(q)$ 's by a simple basis change)

Let's define a complete set of irreps. It is a set  $\chi_1,...,\chi_m$  s.t. that any irrep is unitarily equivalent to one of them. e.g.  $\chi \sim \chi' = V\chi CV^{-1}, V \in U(d)$ 

**Example**: G is abelian, all irreps have d = 1, since all  $\chi(g)$  commute. Theorem (non-abelian analogue of Theorem 1) (consult Fulton and Hardes "Representation Theory" fo more information)

If  $d_1, ..., d_m$  are the dimensions of a complete set of irreps  $\chi_1, ..., \chi_M$  then:

- 1)  $d_1^2 + \dots + d_m^2 = |G|$
- 2)  $\chi_{i,jk}(g)$  is (j,k) th matrix entry of  $\chi_i(g)$  then by Schur orthogonality:

$$\sum_{g} \chi_{i,jk}(g) \bar{\chi}_{i',j'k'}(g) = |G| \delta i i' \delta_{jj'} \delta k k'$$

Now if we look at the states that correspond to these irreps  $\chi_{i,jk} = \sum_{a} \in$  $G\bar{\chi}_{i,jk}(g)|g\rangle$  they form an orthonormal basis.

QFT on G is defined to be a unitary rotation between two basis of  $\{\chi_{i,jk}\}$  basis  $\rightarrow \{|g\rangle\}_{g\in G}.$ 

These takes  $|\chi_{i,jk}\rangle$  are not shift-invariant for all  $U(g_0)$  so this implies that measuring coset state  $|g_0k\rangle$  in the  $\{|\chi\rangle\}$  basis results in an output distribution that is not independent of  $g_0$ .

A "partial" shift-invariance survives. Consider a measurement  $M_{rep}$  on  $|g_0k\rangle$  this measurement will only distinguish the irreps (i values) and not all (i, j, k)'s. the outcome i will be associated with  $d_i^2$  dimensional orthogonal subspaces that are spanned by  $\{\chi_{i,jk}\}_{j,k=1}^{d_i}$ . Then  $\chi_i(g_1g_2) = \chi_i(g_1)\chi_i(g_2) \implies$  the output distribution of i values is indeed

independent of  $g_0$ .

So this gives us direct (but incomplete) information about k. For instance, conjugate subgroups k and  $L = g_0 K g'_0, g_0 \in G$ . This measurement will give us the

 $M_{rep}$  will result in the same output statistics

Not everything is lost there are some cases when this information is enough.

The reason HSP is good in the abelian case is we have an efficient QFT transform. In other words QFT can be implemented in  $poly(\log(|G|))$  times. This is true for abelian groups and some non-abelian groups (e.g.  $P_n$ ). Some partial results:

For normal subgroups  $gk = kgg \in G$  we have a theorem proven by Hall green, russel Ta shma in 2003 SIAM J Comp 32, p 916-934:

Suppose G has QFT that is efficiently implementable. Then if a hideen subgroup k is a normal subgroup, then there is an efficient quantum HSP

Theorem (Edinsguin, Hoyer, Knill, 2004)

For general non-abelian HSP,  $M = O(poly \log(|G|))$  then random coset states  $|g_1k\rangle,...,|g_nk\rangle$  suffice to determine k. But it is not known how to efficiently determine k from the M coset states.

#### 5 Lecture 6

#### Phase estimation Algorithm 5.1

- Unifying principle for quantum algorithms based on QFT
- also gives an alternative way of factoring (originally discovered by Kitaev)

The fact the phase estimation algorithm became so wide spread and you can cast every algorithm that is tangenetially ralted to QFt in QPE

Given a unitary operator U, and eigenstates  $|v_{\phi}\rangle = U|v_{\phi}\rangle = e^{2\pi i \phi}|v_{\phi}\rangle$ 

Want to estiamte the phase  $\phi$  0 <  $\phi$  < 1 (up to n bits of precision  $\phi$  =  $0.i_1i_2,...i_n = i_1/2 + i_2/4 + ...$  for any given n)

We will have to implement controlled unitary opearotrs and in particular we will need Controlled- $U^k$  for integers k

$$C - U^k |0\rangle |\xi\rangle = |0\rangle |\xi\rangle, C - U^k |1\rangle |\xi\rangle = |1\rangle U^k |\xi\rangle$$

 $|\xi\rangle$  has a general dimension d:

$$U^k |v_{\phi}\rangle e^{2\pi i k \phi} |v\phi\rangle$$
,  $C - U^k = (C - U)^k$ 

If we are given U as a circuit description, we can easily implement C-U by controlling each gate in U's circuit. However, if U is given as black box (e.g. a physical operation in the lab) we need further information as there is an inherent ambiguity as we have to account for local phase  $e^{i\theta U}$  as it has no effect normally unless you use a controlled operation.

If the unitary is specified in this ambiguous way we need to figure out what to do. It suffices to have an eigenstate  $|l\rangle$  with a known eigenvalue  $U|l\rangle = e^{i\alpha}|l\rangle$ then  $e^{i\theta}U$  will map  $\alpha \to \alpha + \theta$ . Consider the following circuit:

with  $P(-l) = \begin{pmatrix} 1 & 0 \\ 0 & e^{-il} \end{pmatrix}$ . This correctly gives C - U. We'll want a "generalised controlled-U" that gives:

$$|x\rangle |\xi\rangle \to |x\rangle U^x |\xi\rangle x \in \mathbb{Z}_{2^n}$$

For  $x = x_{n-1}...x_1x_0 = 2^0x_0 + 2^1x_1 + 2^2x_2... + 2^{n-1}x_{n-1}$ :

$$\begin{vmatrix} x_{n-1} \rangle & & & & & \\ |x_1\rangle & & & & & \\ |x_0\rangle & & & & & \\ |\xi\rangle & - U^{2^0} - U^{2^1} & - U^{2^{n-1}} \end{vmatrix}$$

circuit, take  $|\xi\rangle = |v_{\phi}\rangle$ :

This gives output  $|A\rangle=\frac{1}{\sqrt{2^n}}\sum_x e^{2\pi i\phi x}\,|x\rangle$ . Applying  $QFT^{-1}$  to  $|A\rangle$  and measure. We get  $y_0y_1...y_{n-1}$ . Then output the number  $0.y_1y_2...y_{n-1}=\frac{y_0}{2}+\frac{y_1}{4}+...+\frac{y_{n-1}}{2^n}$  as an estimate of  $\phi$ .

Now lets assume an idealised situation where  $\phi$  has only n binary digits:

$$\phi = 0.z_1...z_{n-1}$$

Then  $\phi = \frac{z}{2^n}$  where z is an n-bit integer in  $\mathbb{Z}_{2^n}$ :

$$|A\rangle = \frac{1}{\sqrt{2^n}} \sum_{n} e^{2\pi i 2^n z/2^n} |x\rangle$$

is a QFT of  $|z\rangle$ . Applying  $QFT^{-1}|A\rangle=|z\rangle$  and we get  $\phi$  exactly with certainity.

Note the algorithm up to the final measurements is a unitary operation mapping:

$$|0\rangle |0\rangle \dots |0\rangle |v_{\phi}\rangle \rightarrow |z_{0}\rangle |z_{1}\rangle \dots |z_{n-1}\rangle |v_{\phi}\rangle$$

If  $\phi$  has more than n bits  $\phi = 0.z_0z_1...z_{n-1}|z_nz_{n+1}$ .

**Theorem (PE)**: If measurements in the algorithm give  $y_0y_1...y_n$  and the aoutput  $\Theta = 0.y_0y_1...y_{n-1}$  then :

a) Prob (  $\Theta$  is closeset *n*-binary digit approx to  $\phi$ )  $\geq \frac{4}{\pi^2} = 0.4$ 

b) Prob( $|\Theta - \phi| \ge \epsilon$ )  $\le O(\frac{1}{2^n \epsilon})$ 

We will show that Prob(  $|\Theta - \phi| \ge \epsilon$ )  $\le \frac{1}{2^{n+1}\epsilon}$ )

# 6 Lecture 8

Today we will prove the Phase estimation theorem. We need to change the defintion of distance as we need a distance on a circle.

Define  $d(\theta, \phi) = \min\{|\theta - \phi|, |1 + \phi - \theta|, |1 + \theta - \phi|\}$  which is the distance on the circle. Lets consider the normal binary expansion 0.999999 the closest string should be 1.

**Theorem (Phase Estimation)**: If the output of PE algorithm with n lines (initited in as zeros) is  $\theta = 0.y_0y_1...y_{m-1}$ , then:

a) Prob ( $\theta$  is closeset *n*-binary digit approx to  $\phi$ ,  $d(\theta,\phi) \leq \frac{1}{2^{n+1}}$ )  $\geq \frac{4}{\pi^2} = 0.4$ 

b) Prob( $|\theta - \phi| \ge \epsilon$ )  $\le O(\frac{1}{2^{n_{\epsilon}}})$  for  $\epsilon$  fixed

Recall: the output is obtained by measuring an n-qubit state  $QFT^{-1}|A\rangle$  where  $|A\rangle=\frac{1}{\sqrt{2^n}}\sum_{x=0}^{2^n-1}e^{2\pi ilx}|x\rangle$ 

$$QFT^{-1}|x\rangle = \frac{1}{\sqrt{2^n}} \sum_{y=0}^{2^n - 1} e^{-2\pi i \frac{yx}{2^n}|y\rangle}$$

Soon we will change the notation to make sure it is not overloaded with these powers of  $2^n$ 

$$QFT^{-1}|A\rangle = \frac{1}{2^n} \sum_{y=0}^{\infty} \sum_{x=0}^{\infty} e^{2\pi i(\psi - \frac{y}{2^n})x} |y\rangle$$

Let  $\{\phi = 2^n \psi\}$ :

$$QFT^{-1}|A\rangle = \frac{1}{2^n} \sum_{y=0} \sum_{x=0} e^{2\pi i \frac{\phi - y}{2^n} |x|} |y\rangle = \frac{1}{2^n} \sum_{y} \frac{1 - e^{2\pi i (\phi - y)}}{1 - e^{2\pi i \frac{\phi - y}{2^n}}} |y\rangle$$

In the case  $\phi - y \neq 0$ :

$$Prob(seey) = \frac{1}{2^n} |\sum_{x=0}^{2^n-1} e^{2\pi i \frac{(\phi-y)}{2^n} x}|^2 = \frac{1}{2^{2n}} \frac{|1 - e^{2\pi i (\phi-y)}|^2}{|1 - e^{2\pi i \frac{\phi-y}{2^n}}|^2} = \frac{1}{2^{2n}} \frac{|1 - e^{2\pi i (\psi - \frac{y}{2^n}} 2^n|^2)}{|1 - e^{2\pi i (\psi - \frac{y}{2^n}}|^2)}$$

As  $\theta_y = \frac{y}{2^n}$  and observe that  $|1 - e^{2\pi i (\psi - \theta_y)}|^2 = |1 - e^{2\pi i d(\psi, \theta_y)}|^2$  therefore:

$$Prob(see y) = \frac{1}{2^{2n}} \frac{|1 - e^{2\pi i 2^n d(\psi, \theta_y)}|^2}{|1 - e^{2\pi i d(\psi, \theta_y)}|^2}$$

As  $0 < d(\psi, \theta_y) \le \frac{1}{2}$  we will ise the following bounds: i)  $|1-e^{i\alpha}| \le 2|$ 

- ii)  $|1 e^{i\alpha}| \le |alpha|$

iii) For 
$$|alpha| \le \phi |1 - e^{i\alpha}| = 2|sin(\frac{\alpha}{2})| \ge \frac{2|alpha|}{\pi}$$

The last of the comes from the fact that for positive  $\alpha$  we have  $\sin(\alpha/2) \geq \frac{\alpha/\pi}{2}$ 

When  $d(\psi,\theta) \leq \frac{1}{2^{n+1}}$  implies that y is the best approximation for  $\psi$  and  $2\pi d(\psi,\theta_y)2^n \leq \frac{2^{n+1}}{2^{n+1}}\pi$  so:

$$Prob(yisbestapproximation) \ge \frac{1}{2^{2n}} |\frac{2}{\pi} \frac{2\pi d(\psi, \theta_y)}{2\pi d(\psi, \theta)}|^2 = \frac{4}{\pi^2}$$

The calculations for the above will be on the moodle

#### Further remarks:

If  $C-U^{2^n}$  is implement as  $(C-U)^{2^{\alpha}}$ , then PE algorithm needs exponential time  $(1+2+...+2^{n-1}=2^n-1)$ . But for some some U implementing  $C-U^{2^k}$ requries only polynomial time so we get a poly-time PE algorithm. Harks back to the algorithm for finding powers by repeated squaring, expressing the exponent in binary and then doing repeated squaring. The number of applications of controlled unitaries does not depend on d the dimension of the space. This can be used to rpovide an alternative factoring algorithm (due to A.Kitaev) ( see example sheet).

In many applications we feed an arbitary state to the last register rather than an eigenstate. If instead of  $|v_{\phi}\rangle$  we input general state  $|\xi\rangle$ , expand in eigenbasis of U:

$$|\xi\rangle = \sum_{j} c_{j} |v_{\phi_{j}}\rangle, U |v_{\phi_{j}}\rangle = e^{2\pi i \phi_{j}} |v_{\phi_{j}}\rangle$$

Then we get (before the final measurement) a unitary process  $U_{PE}$ :

$$|00..00\rangle |\xi\rangle \rightarrow^{U_{PE}} \sum_{i} c_{i} |\psi_{i}\rangle |v_{\phi_{i}}\rangle$$

The Born rule implies that the final measurement will give a choice of  $\phi_i$ 's (or an approximation) it can be chosen with probability  $|c_i|^2$ . This is not some average of the  $\phi_j$  values.

Will be elaborated more in the notes on moodle on the following:

If you want to have n-qubits and want to get m-bits correctly probability of success  $1 - \eta$ , then must have :

$$n \ge m + \log \frac{1}{\eta}$$

# 6.1 Amplitude amplification

Much like when we multiple up to HSP we revisited shors alogrithm, in this case we revisit grovers algorithm. This is an apotheosis of technique in Grover's algorithm.

#### 6.1.1 Background

**Reflection Operators**: State  $|\alpha\rangle$  in  $\mathcal{H}_d$ , n-dim subspace  $L_\alpha$  with (d-1) dim orthogonal  $l_\alpha^\alpha$ 

$$\begin{split} I_{|\alpha\rangle} &= I - 2 \, |\alpha\rangle \, \langle \alpha| \\ I_{|\alpha\rangle} &= - \, |\alpha\rangle \\ I_{|\alpha\rangle} \, |\beta\rangle &= |\beta\rangle \end{split}$$

for any  $|\beta\rangle |\alpha\rangle$ 

For any unitary  $U: UI_{|\alpha\rangle}U^{\dagger} = I_{U|\alpha\rangle}, U|\alpha\rangle\langle\alpha|U^{\dagger} = |\beta\rangle\langle\beta| \text{ for } |\beta\rangle = U|\alpha\rangle.$ 

Consider a k-dimensional subspace  $A < \mathcal{H}_d$ , any orthonormal basis  $|a_1\rangle, ..., |a_k\rangle$ . Lets consider a projection operator on to this subspace:

$$P_A = \sum_{i=1}^k |a_i\rangle \langle a_i|$$

Define a generalised projection operator:

$$I_A = I - 2P_A$$

$$\{|\xi\rangle \quad |\xi\rangle \in A$$

$$I_A |\xi\rangle = \begin{cases} |\xi\rangle & |\xi\rangle \in A^{\dagger} \\ -|\xi\rangle & |\xi\rangle \in A \end{cases}$$

Now lets recall what grover does very briefly (have a look at Part II course):

Search for a unique "goal" item in unstructured database of  $N=2^n$  items.

Write  $B_n = \text{set}$  of all n-bit strings. Given an oracle for  $f: B_n \to B_n$  with the promise that there is a unique element  $x_0 \in B_1$  with f(0) = 1. Problem is to find  $x_0$ .

Closely related to class NP and Boolean satisfiability problem. This is explained in part II lecture notes.

# 7 Lecture 8

**Recap of grovers algorithm** We are searching for a unique "good" element in an unstructured database,  $N=2^n$  items

We are given an oracle f that maps from  $B_n \to B_1$ Promise: There is a unique  $x_0 \in B_n$  with  $f(x_0) = 1$ Problem: Find  $x_0$ 

Consider Grover iteration operator on n qubits

$$Q = -H_n I_{|0\rangle} H_n I_{|x_0\rangle} = -I_{|\psi_0\rangle} I_{|x_0\rangle}$$

where 
$$H_n = H \otimes ... \otimes H$$
,  $|\psi_0\rangle = H_n |00.000\rangle = \frac{1}{\sqrt{2n}} \sum_{x \in B_n} |x\rangle$ 

One application of Q uses 1 query of  $U_f$ 

**Thereom (Grover '96)**: In 2-dim span of  $|\psi_0\rangle$  and (unknown)  $|x_0\rangle$  the action of Q is rotation by angle  $2\alpha$  where  $\sin\alpha = \frac{1}{\sqrt{N}}$ . Hence grover's algorithm to find  $x_0$  given  $U_f$  is:

Make  $|\psi_0\rangle$ 

Apply Q m times where  $m = \frac{\arccos \frac{1}{\sqrt{N}}}{2} \arcsin \frac{1}{\sqrt{N}}$  to rotate  $\psi_0$  very close to  $x_0$  (within the angle  $\pm \alpha$ 

Measure to see  $x_0$  with high probability  $1 - \frac{1}{N}$  For large N arccos  $\frac{1}{\sqrt{N}} = \frac{\pi}{2}$  and arcsin  $\frac{1}{\sqrt{N}} = \frac{1}{\sqrt{N}}$  so  $m = \frac{\pi}{4}\sqrt{N}$  iteractions or queries to  $U_f$  needed.

Classically we need O(N) quereies to find  $x_0$  with any constant probability that does not depend on N, so this achieves a quadratic speed-up.

## 7.1 Amplitude Amplification

Let G be any subspace ('good subspace') of state space  $\mathcal{H}$  G be its orthogonal complement ('bad subspace')  $\mathcal{H} = G \oplus G$ 

Given any  $|\psi\rangle \in \mathcal{H}$ , we have unique decomposition with real positive coefficient  $|\psi\rangle = \sin\theta |g\rangle + \cos\theta |b\rangle, |g\rangle \in G, |b\rangle \in G$ Introduce reflection operators that flip  $|\psi\rangle$  and good vectors:

$$I_{|\psi\rangle} = I - 2 |\psi\rangle \langle \psi|, I_G = I - 2P_G$$

 $\sin\theta = ||P_G\left|\psi\right\rangle||$  =length of good projection of  $|\psi\rangle$  Introduce  $Q=-I_{|\psi\rangle}I_G$ 

#### 7.1.1 Amplitude Amplification Theorem

In the 2-dim space spanned by  $|g\rangle$  and  $|\psi\rangle$  Q is rotaion by  $2\theta$  where  $\sin\theta$  =length of good projection of  $|\psi\rangle$ 

**Proof**: We have  $I_G |g\rangle = -|g\rangle$ ,  $I_G |b\rangle = |b\rangle$ :

$$Q |g\rangle = I_{|\psi\rangle} |g\rangle, Q |b\rangle = -I_{|\psi\rangle} |b\rangle$$

$$I_{|\psi\rangle} = I - 2(\sin\theta |g\rangle + \cos\theta |b\rangle)(\sin\theta \langle g| + \cos\theta \langle b|)$$

using the fact that  $\langle b | | g \rangle = 0, \langle g | | g \rangle = \langle b | | b \rangle = 1$ :

$$Q|b\rangle = \cos 2\theta |b\rangle + \sin 2\theta |g\rangle$$

$$Q|G\rangle = -\sin 2\theta |b\rangle + \cos 2\theta |g\rangle$$

So

$$Q\left|g\right\rangle = I\left|g\right\rangle - 2\sin^{2}\theta\left|g\right\rangle - 2\sin\theta\cos\theta\left|b\right\rangle = \cos2\theta\left|g\right\rangle - \sin2\theta\left|b\right\rangle$$

In the  $\{|b\rangle, |g\rangle\}$  basis, the Q matrix is a rotation matrix by  $2\theta$ :

$$Q = \begin{pmatrix} \cos 2\theta & \sin 2\theta \\ -\sin 2\theta & \cos 2\theta \end{pmatrix}$$

As we apply Q n  $timesweget Q^n | \psi \rangle = \sin(2n+1)\theta | g \rangle + \cos(2n+1)\theta | b \rangle$ . If  $timesweget Q^n | \psi \rangle = \sin(2n+1)\theta | g \rangle + \cos(2n+1)\theta | b \rangle$ . If  $timesweget Q^n | \psi \rangle = \sin(2n+1)\theta | g \rangle + \cos(2n+1)\theta | g \rangle$  so this is maximised when  $(2n+1)\theta = \frac{\pi}{2}$ . So forthenear estint egert  $times = \pi \frac{1}{4\theta - \frac{1}{2}}$  we will be within  $\theta$  of the element.

**Example**: If we have  $\theta = \frac{\pi}{6}$ , n = 1 we can see that  $Q^1$  rotates  $|\psi\rangle$  exactly onto  $|g\rangle$ .

Generally, for a given  $\theta$  n is not an integer, so we use n = nearest integer to  $\frac{4\pi}{\theta} - 1 \approx \frac{4\pi}{\theta} = O(\frac{1}{\theta}) = O(\frac{1}{\sin \theta}) = O(\frac{1}{\operatorname{length of good projection of } |\psi\rangle}$  and then  $Q |\psi\rangle$  will be witin angle  $\theta$  of  $|g\rangle$  so the probability of seeing a good value is:  $P \geq \cos \theta = 1 - O(\theta^2)$ .

All this can implement if  $I_{|\psi\rangle}I_G$  can be impreneted efficiently. see ES2.

For  $I_G$  is suffices for G to be spanned by computational basis states and have an indicator function f.

$$f(x) = 1$$
, xifx is good,  $f(x) = 0$  ifx is bad

For  $I_{|\psi\rangle}$  we usually have  $|\psi\rangle = H_n |00..000\rangle$ , then  $|\psi\rangle$  can be implemented in O(n) time where n is the number of qubits.

In the amplitude amplification algorithm the relative amplitudes of good elements remain the same as they were in  $|\psi\rangle = \sin\theta \, |g\rangle + \cos\theta \, |b\rangle$  so  $|g\rangle$  remains the same just the amplitude varies. So AA amplifies overall  $|g\rangle$  amplitude at the expense of reducing the amplitude of  $|b\rangle$ .

Second remark: Final state is generally not exactly  $|g\rangle$ , however, if  $\sin\theta$  is known then there is a modification of this algorithm that uses a modest amount of resources t omake it exact.

The rotuine is useful for state prepation e.g.

$$\sum_{x < N, \mathbf{x} \text{ is coprime to N}} |x\rangle$$

# 8 Examples Class 1

Try to prove the same thing in 1(ii) in shors algorithm and it will be slightly nicer language than group theory.

For  $\mathbb{Z}_2$  the irreps are  $\chi_a(x)=(-1)^{ax}$  for  $a,x\in\mathbb{Z}_2$  therefore  $\mathbb{Z}_2^n$  has irreps  $\chi_a(x)=(-1)^{a_1x_1+a_2x_2+...+a_nx_n}$  so

$$|\chi\rangle = \frac{1}{\sqrt{|G|}} \sum_{g} \bar{\chi}(g) |g\rangle = \frac{1}{\sqrt{|G|}} \sum_{b \in \mathbb{Z}_2^n} (-1)^{ab} |b\rangle$$

and

$$[QFT]_{ab} = \frac{1}{\sqrt{|G|}}(-1)^{ab} = \frac{1}{\sqrt{2^n}}(-1)^{a_1b_1}(-1)^{a_2b_2}...(-1)^{a_nb_n}$$

could be written using Hadmards with  $QFT = H \otimes H \otimes ... \otimes H$  as  $H = \frac{1}{\sqrt{2}} \begin{pmatrix} 1 & 1 \\ 1 & -1 \end{pmatrix}$  with the columns being a and the rows being b so it is only -1 if they are both 1 as otherwise it will have a zero in the exponent.

First part of HSP problems is generate coset states and the second part is repeation (how many times you need to repeat the process). In standard HSP we make one query to f to make the following state:

$$|y \oplus k\rangle = \frac{1}{2^k} \sum_{x \in K} |y \oplus x\rangle, y \in \mathbb{Z}_2^n$$

Apply  $QFT = H^{\otimes n}$  and measure. This gives an uniformly random output  $c \in \mathbb{Z}_2^n$  which is such that the irrep  $\chi_c$  of G that is restricted to K is the trivial irrep of K. In other words  $\chi_c(a) = 1$  for all  $a \in K$ . In other word,  $(-1)^{ac} = 1 \implies ca = 0 \mod 2$ . We know that this k viewed as a subspace of  $\mathbb{Z}_2^n$  has dimension k. So we need (n-k) linearly independent  $c_i$  with  $c_i a = 0$  to determine K.

In order to succeed with probability  $1 - \epsilon$  when given a constant probability p we run the process some M times. Then you calculate the probability of M runs failing to determine k and show that for a constant overhead we can get any accuracy.  $1 - (1 - p)^M > 1 - \epsilon$ .

Phase gates act with  $P(\alpha):|0\rangle \rightarrow |0\rangle, |1\rangle \rightarrow e^{i\alpha}|1\rangle$ . Therefore can apply a fractional phase by first preparing a register with  $i_1,...,i_n$  s.t.  $y=0.i_1...i_n$  and then by applying phase gates  $P(\frac{1}{2}...P(\frac{1}{2^n}\text{to get the state }e^{iy}|y\rangle.$ 

When inventing algorithms whilst it is useful to think about eigenstates to start with make sure to run it through with a general state as you might need to uncompute at the end. e.g. need to apply an inverse controlled unitary in question 5 after computing the correct eigenvalue.

We can't just discard additional registers like  $\sum_j \lambda_j |u_j\rangle |c_j\rangle$  to  $\sum_j \lambda_j |u_j\rangle$  we need to uncompute to remove these extra stuff.

# 9 Lecture 9

# 9.1 Amplitude amplification

We have a Hilbert space we can partition into a good part and a bad part:  $\mathcal{H} = G_{good} \oplus G_{good}^{\dagger}$ , for every  $|\psi\rangle$  we have  $|\psi\rangle = \sin\theta |g\rangle + \cos\theta |b\rangle$ . We have:

$$I_{|\psi\rangle} = I - 2 |\psi\rangle \langle \psi|, I_G = I - 2P_G, Q = -I_{|\psi\rangle}I_G$$

We proved AA Thereoem: That in the plane spanned by  $|b\rangle$  and  $|g\rangle$  Q is a rotation by  $2\theta$ .

Typically we are given some  $|\psi\rangle$  and we use this property to rotate it to  $|g\rangle$ .

# 9.2 Applications of Amplitude Amplification

# Grover search with one or more (k) good items in N:

Re look over Part II lecture notes and reprove for k good elemetrs. Maybe try problem sheet again.

$$|\psi\rangle = \frac{1}{\sqrt{2^n}} \sum_{x \in B_n} |x\rangle = \sqrt{\frac{k}{N}} (\frac{1}{\sqrt{k}} \sum_{goodelem} |x\rangle) + \sqrt{\frac{N-k}{N}} (\frac{1}{\sqrt{N-k}} \sum_{bad} |x\rangle)$$

G is spanned by k good  $|x\rangle$ 's,  $\sin\theta = \sqrt{\frac{k}{N}}$  so Q is rotation by  $22\theta$ ,  $\theta = \arcsin\sqrt{\frac{k}{N}} \implies O(\sqrt{\frac{N}{k}})$  queries.

Note that for 2-bit case N=4 and k=1 good element we get  $\theta=\arcsin\frac{1}{2}=\pi 6$  and one application of Q rotates  $|\psi_0\rangle$  exactly onto  $|g\rangle$ 

### Square-root speedy of general quantum algorithms:

Let A be a quantum algorithm/circuit(sequence of 'basic' unitary gates) on input states  $|00.00\rangle$ . The final state is  $A |00...000\rangle$ .

Good labels = desired computation outcomes.

 $A|0...0\rangle = \alpha |a\rangle + \beta |b\rangle$ ,  $\alpha = \sin \theta$  with  $|a\rangle$  normalised but generally unequa superposition  $\sum_{qoodx} c_x |x\rangle$ .

So probability of success in 1 run is  $|\alpha|^2$  so need to do  $O(\frac{1}{|\alpha|^2})$  repetitions of A to succeed with any given const probability  $1 - \epsilon$ .

Now lets try amplitude amplification instead of repeated measurement. We need to check we satissify some assumptions.

Assume we can check if the answer is good or bad.

So that we can implement  $I_G|x\rangle \to \begin{cases} -|x\rangle & \text{x is good} \\ |x\rangle & \text{x is bad} \end{cases}$  Consider  $|\psi\rangle = A|0.00\rangle$ 

$$Q = -I_{A|00.000\rangle}I_G = -AI_{|0..00\rangle}A^{\dagger}I_G$$

so all parts are implementable

By the amplitude amplifacation theorem Q is a rotation by  $2\theta$  with  $\sin\theta = |\alpha|$ . So after  $n = \frac{\pi}{4\theta} = O(\frac{1}{|\theta|}) = O(\frac{1}{\sin\theta}) = O(\frac{1}{|\alpha|}$  repetitions  $A|0..0\rangle$  will be rotated very near  $|g\rangle$  and the final measurement will succeed with high probability. So we get a square root time speed up over a discrete repetition method.

Each application of Q need one A and one  $A^{\dagger}$ . and can think of  $A^{\dagger}$  as inverse gates in reverse order so it has a similar complexity to A.

If success probability of A is known then we can do even better and we can cook up the rotation so exactly get to  $|g\rangle$  (this exact method will be covered in ES2). This will convert the probabilistic algorithm A into a deterministic one.

#### Quantum counting:

Here we apply amplitude amplification and phase estimation.

Given  $f: B_n \to B_1$  with k good x's we want to estiamte k with (instead of finding a good x).

Let's recall the grover operator:  $Q_G$  for f is a rotation by  $2\theta$  in this 2D plane spanned by  $|\psi_0\rangle = \frac{1}{\sqrt{2^n}} \sum_{x \in B_n}$  and its good projection  $|g\rangle = \frac{1}{\sqrt{k}} \sum_{goodx} |x\rangle$  with  $\sin \theta = \sqrt{\frac{k}{N}} \approx \theta$ 

Recall that rotations in  $\{|b\rangle, |g\rangle\}$  palne have eigenvaleus and eigenvectors:

$$|e_{\pm}\rangle = \frac{1}{\sqrt{2}}(|b\rangle \pm i |g\rangle)$$

$$\lambda_{\pm} = e^{\pm 2i\theta}$$

Then check that  $|\psi_0\rangle = \sin\theta |g\rangle + \cos\theta |b\rangle = \frac{1}{\sqrt{2}}(e^{i\theta} |e_+\rangle + e^{-i\theta} |e_-\rangle)$ Writing  $e^{\pm 2i\theta}$  as  $e^{2\pi i\phi_{\pm}}$  with  $\phi_{\pm} \in (0,1)$ :

$$\phi_+ = \frac{2\theta}{2\pi} = \frac{\theta}{\pi}$$

$$\phi_- = \frac{-2\theta + 2\pi}{2\pi} = 1 - \frac{\theta}{\pi}$$

$$\frac{\theta}{\pi} = \frac{1}{\pi} \sqrt{\frac{k}{N}}$$

Run QPE algorithm with U= "Grover Q" and estimate register set to  $|\psi_0\rangle$ : will output  $\frac{\theta}{\pi}$  or  $1-\frac{\theta}{\pi}$  with probability  $\frac{1}{2}$ 

# 10 Lecture 10

$$f: B_n \to B_1$$

which has k good x's and we want to estimate k. This is harder than just finding a solution.

$$|\psi_0\rangle = \frac{1}{\sqrt{2^n}} \sum_{x \in B_n} |x\rangle, |g\rangle = \frac{1}{\sqrt{k}} \sum_{goodx} |x\rangle$$
  
$$\sin \theta = \sqrt{\frac{k}{N}}$$

for  $k \ll N$  so

$$|\psi_0\rangle = \sin\theta |g\rangle + \cos\theta |b\rangle = \frac{1}{\sqrt{2}} (e^{i\theta} |e_+\rangle + e^{-i\theta} |e_-\rangle)$$

Run QPE -;  $\frac{\theta}{\pi}$  or  $1 - \frac{\theta}{\pi}$  with probability  $\frac{1}{2}$ . We can get approximately  $\theta$  in either case. For any m, running QPE with m qubit lines will give an m-bit approximation to  $\sqrt{\frac{k}{N}}$  this uses  $2^n$  C-Q gates.

By Theroem PE we learn  $\sqrt{\frac{k}{N}}$  to an additive error with constant probability  $\frac{4}{\pi^2}$  using  $O(2^m)$  queries. Write  $\frac{1}{2^n}$  as  $\frac{\delta}{\sqrt{N}}$   $\delta > 0$ . Thus we can learn  $\sqrt{k}$  to additive error  $O(\delta)$  using  $O(2^n) = O(\frac{\sqrt{N}\delta}{2})$  queries.

Classically the same approximation (obtained with constant probability) requires  $O(\frac{N}{\delta^2})$  so requires quadratically more effort.

#### 10.1 Hamiltonian simulation

One of the most promising applications of qunatum computers is simulating quantum systems. This is difficult for classical computers as these physical systems have so many degrees of freedom. We won't do the nitty gritty but cover some of the core fundamental ideas and talk about the complexity of it.

We want to use a quantum computer to simulate the evolution/dynamics of a quantum system given its Hamiltonian H.

For n qubits generally requires  $O(2^n)$  time on a classical computer. We will do it in poly(n) time for a class of Hamiltonian. A lot of good hamiltonians that corrospond to real problems do admit a good representations in quantum computers. First we will focus on evolution and dynamics and then we will move on to studying the ground state properties of these hamiltonians.

# 10.1.1 Hamiltonain and quantum evolution

Consider a physical system in state  $|\psi\rangle$ , with hamiltonain H: selfadjoint (Hermitian) operator/matrix - the quantum energy observable. Time evolution given

by Schrodinger equation  $(\hbar = 1)$ 

$$\frac{d}{dt} |\psi(t)\rangle = -iH |\psi(t)\rangle$$

We will consider time-independent Hamiltonians H(t) = H

$$|\psi(t)\rangle = e^{-iHt} |\psi(0)\rangle$$

$$e^A = I + A + \frac{A^2}{2!} + \dots$$

Thus given H and time t we want to simulate the action of a unitary  $U(t) = e^{-iHt}$  to a suitably good approximation.

Approximation (closeness) of unitary operators Operator norm (spectral norm)  $||A|| = \max_{||\psi||=1} ||A|\psi\rangle|| = |\max|\min|\min|\exp|-(ifAisdiagonalisable)$ .

$$||A + B|| \le ||A|| + ||B||$$

We say that U approximates  $\tilde{U}$  to within error  $\epsilon$  if  $||U - \tilde{U}|| \le \epsilon$  (for any  $|\psi\rangle$  the action of U  $\tilde{U}$  are at most  $\epsilon$  apart.

In general this hamiltonain will be a huge  $2^n \times 2^n$  matrix for an nqubit system so hard to write down so we will consider a special class of k-local hamiltonains. We will want to simulate  $U = e^{-itH}$  with a circuit of poly(n,t) "basic" unitary gates (this is the defintino of efficient simulation on a quantum computer).

Not all H's can be efficiently simulated. One such class of efficient hamiltonains is k-local hamiltonains: k-local Hamiltonain: H is k-local on n qubits if  $H = \sum_{j=1}^m H_j$  where  $H_j$  is a hermitain matrix acting on at most k qubits. This action does not have to be continuous a.k.a they don't need to be qubits 1-k it can be any subset of k qubits.

$$H_i = \tilde{H}_i \otimes I$$

with  $\tilde{H}_j$  encoding the hamiltonain on the k qubits and the I is on the rest of the system.

 $m \leq \binom{n}{k} = O(n^k) = poly(n)$  terms in H. We will see that many important classes of hamiltonains fall into this category e.g.:

1. (3 qubits, 2-local)

$$H = X \otimes I \otimes I - 5Z \otimes I \otimes Y$$

2. Write  $M_{(k)}$  to denotes an operator M acting on k qubits and I on the rest. This can correspond to the Ising model on  $n \times n$  square lattice of qubits.

$$H = J \sum_{i,j=1}^{n-1} Z_{(i,j)} Z_{(i,j+1)} + Z_{(i,j)} Z_{(i+1,j)}$$

Heisenberg model on a line:

$$H = \sum_{i=1}^{n-1} J_x X_{(i)} X_{(i+1)} + J_y Y_{(i)} Y_{(i+1)} + J_z Z_{(i)} Z_{(i+1)}$$

with  $J_x, J_y, J_z$  are real coeff.

Remember that  $e^{-\sum_i H_j t} \neq \prod_j e^{-H_j t}$  if  $H_j$  don't commute Also remember that  $e^{-H_j t}$  are local unitaries acting on k qubits.

# 11 Lecture 11

If we want to use some standard universal gate set, then we will invoke the theorem:

**Solovay-Kitaev Theorem**: Let U be a unitary operator on k (const) qubits, and S be any universal set of quantum gates. Then U can be approximated to within accuracy  $\epsilon$  using a logarithm sequence of gate  $O(\log^c(\frac{1}{\epsilon}))$  gates from S with c < 4.

We don't prove this result but it is in Nielsen and Chuang.

**Lemma A** about accumulation of errors: We need to prove this on ES2. Let  $\{u_i\}$ ,  $\{v_i\}$  be sets of m unitary operators such that  $||U_i - V_i|| \le \epsilon$ . So if we want to approximate a whole sequence of  $U_i$ s:  $||U_m...U_1 - V_m...V_1|| \le m\epsilon$  so errors will accumulate linearly. This is kind of the worst case error as it is the maximium error over all  $|\psi\rangle$ .

Proof: using induction on m.

First lets consider a little warm up with the (easy) commuting case.

Consider the case:

$$H = \sum_{j=1}^{m} H_j$$
, any k-local Hamiltonian with commuting  $H_j$ 

Then for any power  $t,\ e^{-iHt}$  can be approximated to within  $\epsilon$  by a circuit  $O(mpoly(\log(\frac{m}{\epsilon})))$  from any given universal set. (Note that  $m=\binom{n}{k}=O(n^k)$ , this is  $\operatorname{poly}(n_1\log\frac{1}{\epsilon})$  too and  $\log\frac{1}{\epsilon}$  is the number of digits in the precision of the approximation

Proof: Using SK theorem each  $e^{-iH_jt}$  can be approximated to within  $\frac{\epsilon}{m}$  with  $O(poly(\log \frac{m}{\epsilon}))$  gates.

Lemma A implies the full product  $\prod_{i=1}^{m} e^{-iH_{j}t}$  is then approximated to within  $m\frac{\epsilon}{m} = \epsilon$  using a total of  $O(mpoly(\log \frac{m}{\epsilon}))$  gates.

#### The full non-commuting case

For any matrix X we write  $X + O(\epsilon)$  for  $X + \epsilon$ , where  $||E|| = O(\epsilon)$ 

**Lemmna B**(Lee-Trotter formula):

Let A, B be the following matrices with  $||A|| \le K$  and  $||B|| \le K$  with K < 1 (small). Then:

 $e^{-iA}e^{-iB} = e^{-i(A+B)} + O(K)^2$ 

Proof:  $e^{-iA} = I - iA + \sum_{k=2} \frac{(-iA)^k}{k!} = I - iA + (iA)^2 \sum_{k=0} \frac{(-iA)^k}{(k+2)!}$  we want to show that the sum in the last term has norm  $< e^{-K} < 1$ . Therefore,  $e^{-iA} = I - iA + O(K^2)$  so therefore:

$$e^{-iA}e^{-iB} = (I - iA + O(K^2))(I - iB + O(K^2)) = I - i(A + B) + O(K^2) = e^{-i(A + B)} + O(K^2)$$

Now, apply Lee-Trotter formula repeatedly to accumulate sums of  $H_1...H_m$  in the exponent

Note that if each term  $||H_j|| < K$ , then  $||\sum_{i=1}^l H_i|| < lK$  and we'll want this expression here to be < 1 for all  $l \le m$ . For now, we'll assume that  $K < \frac{1}{m}$ . For now we will also take t = 1 (deal with general t later).

$$e^{-iH_1}e^{-iH_2}...e^{-iH_m} = [e^{-i(H_1+H_2)} + O(K^2)]e^{-iH_3}...e^{-iH_m}$$

So as  $||H_1 + H_2|| < 2K$  and the error that is error that is generated at each stage stays the same for subsequent unitaries  $U_i$  as ||AU|| = ||A|| therefore:

$$e^{-iH_1}e^{-iH_2}...e^{-iH_m} = e^{-i(H_1+H_2)}e^{-iH_3}...e^{-iH_m} + O(K^2)$$

Then by Lee-Trotter for  $(H_1 + H_2)$  and  $H_3$ 

$$e^{-iH_1}e^{-iH_2}...e^{-iH_m} = [e^{-i(H_1+H_2+H_3)} + O((2K)^2)]e^{-iH_4}...e^{-iH_m} + O(K^2)$$

as  $||H_1 + H_2|| < 2K$ . We can continue in the same fashion until we get **error estimate 1**:

$$e^{-iH_1}e^{-iH_2}...e^{-iH_m} = e^{-i(H_1 + ... + H_m)} + O(k^2) + O((2k)^2) + ... + O(((m-1)k)^2) = e^{-i(H_1 + ... + H_m)} + O(m^3k^2)$$

For general finite  $||H_j||$ 's and t values  $||H_jt|| < Kt$  (this value can be large). So we introduce a way of breaking down the sequence by introducing N (that we will fix later) s.t.  $\frac{H_jt}{N}$  gives  $\tilde{K} = ||\frac{H_jt}{N}|| < \frac{Kt}{N}$ . We can think about this as dividing the time T up into small  $\frac{1}{N}$  intervals and then our unitary has the value:

$$U = e^{i(H_1 + \dots + H_m)t} = \left(e^{i(\frac{H_1t}{N} + \dots + \frac{H_2t}{N})}\right)^N$$

We want final error for U to be less than  $\epsilon$  so by lemma A we want the error of each step to be less than  $\frac{\epsilon}{N}$ . So using error estimate 1 we get

$$Cm^3\tilde{K}^2 < \frac{\epsilon}{N} \implies N > \frac{Cm^3k^2t^2}{\epsilon}$$

Then —  $e^{-i\frac{H_1t}{N}}...e^{-i\frac{H_nt}{N}} - e^{-i(H_1t+...H_nt)}|| < \frac{\epsilon}{N}$ . BY Lemma A:

$$||(e^{-i\frac{H_1t}{N}}...e^{-i\frac{H_nt}{N}})^N - e^{-i(H_1t+...H_n)}|| < \epsilon$$

So the total circuit size is  $O(m^4 \frac{(Kt)^2}{\epsilon})$  we are dealing with Hamiltonians that are k-local. So for n qubits and k-local terms  $m = O(n^k)$ . Therefore, the total circuit size is  $O(n^{4k} \frac{(Kt)^2}{\epsilon})$ . You can actually refine this method to get order  $t^{1+\delta}$  but they are much more technical and would need many more lectures to explain.

We have a circuit of size  $|C| = O(\frac{m^4(kT)^2}{\epsilon})$  and if we want to use gates from the universal set Lecture 12

# 11.1 The Local Hamiltonaian Problem and QMA

Recap: Defintion of NP. There is no known way to solve it in polynomial time but it is easy to verify the solution in polynomial time.

We will adopt the language of the part II lecture course. A language is in NP if it has an efficient (poly-time) verifier V. A verifier V for a language L is a computation with two inputs w, c s.t. if  $w \in L$  then for some c V(w,c) halts with "accept". such c is called a certificate/proof (of membership) for w. If  $w \not\in L$ , then for all C, V(w,c) halts outputing "reject". V is a poly-time verifier for all pairs (w,c). V runs in poly(n) time n=|w|.

#### 11.1.1 The satisficatibility problem (SAT)

Boolean forumlas  $\phi(x_1...x_n): \{0,1\}^n \to \{0,1\}$  and every  $(b_1...b_n), b_i \in \{0,1\}$  s.t.  $\phi(b_1...b_n) = 1$  is called a satisfying assignment.  $V(\phi,c)$  evaluates  $\phi(c)$ . SAT is not known to be in P.

Theory of NP-completeness shows that many different problems that look very different (SAT, Travelling Salesman Problem, integer linear programming) are essentially the same problem. We can translate instances of one to another in deterministic poly time.

Now lets try to relax some of the requirements on NP. Consider a setting where the prover and verifier may use randomness and sometimes make errors ( allowing for some probability of error say  $\frac{1}{3}$ ). In other words when  $w \in L$  a prover shoul dbe able to prepare a certificate/proof s.t.  $\operatorname{Prob}(V(w,c))$  accepts)  $\geq \frac{2}{3}$  and  $\operatorname{Prob}(V(w,c))$  rejects)  $\leq \frac{1}{3}$ . This defines complexity class MA (Merlin-Arthur). Merlin is regarded as the omniscient prover and Arthur is a randomised poly time verifier.

### Quantum Merlin Arthur class

The direct analgoue of NP.

It is a class of promise problems. A promise problem L partitions  $\{0,1\}^k$  of all binary strings into  $L_0, L_1, L_*$ .

An algorithm is promised that it never recieves inputs from  $L_*$ . However, if the input is from  $L_0/L_1$  then it has to determine 0/1.

**QMA**: A primise problem  $L = (L_0, L_1, L_*)$  is in the class QMA if there exists a uniform family of circuits  $\{C_n\}$  with two input registers and one output qubit

and a prolynomial p() s.t.  $\forall w \in \{0,1\}^k$ 

Completeness: If  $w \in L_1 \cap \{0,1\}^n$  then there exists a p(n)- qubit state  $|\psi\rangle$  (a proof/witness state) such that  $C_n$  outputs 1 with probability  $\geq \frac{2}{3}$  when run on w and  $|\psi\rangle$ 

Soundness: If  $w \in L_0 \cap \{0,1\}^n$  then for every p(n) qubit state  $|\psi\rangle$  the circuit  $C_n$  outputs 1 with probability  $\leq \frac{1}{3}$  when run on w and  $|\psi\rangle$ 

**QMA**:  $Aprmiseproblem L = (L_0, L_1, L_*)$  is in the class QMA if there exists a uniform family of circuits  $\{C_n\}$  with two input registers and one output qubit and a prolynomial p() s.t.  $\forall w \in \{0,1\}^k$ 

Completeness: If  $w \in L_1 \cap \{0,1\}^n$  then there exists a p(n)- qubit state  $|\psi\rangle$  (a proof/witness state) such that  $C_n$  outputs 1 with probability  $\geq \frac{2}{3}$  when run on w and  $|\psi\rangle$ 

Soundness: If  $w \in L_0 \cap \{0,1\}^n$  then for every p(n) qubit state  $|\psi\rangle$  the circuit  $C_n$  outputs 1 with probability  $\leq \frac{1}{3}$  when run on w and  $|\psi\rangle$ 

#### Remarks:

- 1) If we replace  $|\psi\rangle$  with a classical bitstring we get the class QCMA
- 2) If we additionally to (1) force the verifier to be classical MA
- 3) If we additionally to (1),(2) repalce success probability  $(\frac{2}{3})$  by 1 -; NP

$$NP \le MA \le QCMA \le QMA$$

SAT is NP-complete (Cook-Levin theorem). Consider a special case: k-SAT  $\phi$  is the conjunction of clauses each of which is a disjunction of k-literals. e.g. for k=3:

$$(x_1 \cup \bar{x}_2 \cup x_3) \cap (\bar{x}_1 \cup x_2 \cup x_4) \cap (x_1 \cup \bar{x}_4 \cup x_5)$$

k-SAT is NP-complete for  $x \geq 3$  and k-SAT is in P when k = 2.

We will reformulate k - SAT and relate it to a minimal eigenvalue of a certain hamiltonian which is diagonal in computational basis.

Fix k=3. Consider a claus  $C=x_1\cup \bar{x}_2\cup x_3$  (has one non-satsifiy assignment 010).

Lets associate a particular diagonal hamiltonian with 0 everywhere but a 1 at localation indexed by the bit string above.

 $H_c$  gives "penalty" of 1 to the bitstring of  $x_1...x_n$  if x does not satisfy clause C. We will regard  $H_c$  as part of n-qubit Hamiltonian as  $H_c = H_c \otimes I$ .

When we evaluate this term:  $\langle x|H|x\rangle = 0$  if clause C is satsified or 1 otherwise.

Suppose we have a 3-SAT forumla  $\phi = C_1 \cap ... \cap C_n \implies H_{\phi} = \sum_{j=1}^m H_{c_j}$ . The eigenvalues of  $H_{\phi}$  lie in the interval [0, m] and  $H_{\phi}$  is 3-local.

Each assignment  $x \in \{0,1\}^n$  generates the "energy"

 $\langle x|H_{\phi}|x\rangle=\sum_{j=1}^{n}\langle x|H_{C_{j}}|x\rangle$  this counts the number of unsatisfied clauses under x.

We want to find the minimal energy of this hamiltonain which will correspond to the lowest eigenvalue  $\lambda_n$  ( the smallest number of unsatsified clauses).

# 12 Lecture 12

The k-local Hamiltonian problem (LH): Given a classical description of a Hamiltonian of an n qubit Hamiltonian:

$$H = \sum_{j=1}^{m} H_j$$

where each  $H_j$  is k-local and is positive semi-definite. Also given two parameters  $a,b \in [0,m]$  with  $b-a \geq \frac{1}{poly(n)}$  with a promise that  $\lambda_{min}$ - minimal eigenvalue of H is either  $\leq a$  or  $\geq b$ . Decide which is the case.  $\lambda_{min}$  is often called the ground state energy. Determining  $\lambda_{min}$  ( $\leq a$  or  $\geq b$ ) is tantemount to approximating the ground state energy  $\lambda_{min}$  up to an additive error O(b-a). Further studied on ES.

We would like to prove that LH is QMA-complete

We first need to obtain that LH is in QMA: Given the corresponding witness states  $|\psi\rangle$  for  $x \in L_1$ , we can efficiently approximate the energy to check if it is  $\leq a$  or  $\geq b$ .

Note: LH is a promise problem, which means that H with  $\lambda \in (a,b)$  won't satisfy the condition.

Next step, is we want to show that LH is QMA-hard (any other problem in QMA can be reduced to it)

Take any  $L = (L_1, L_0, L_*)$  in QMA and fix  $x \in L_1 \cup L_0$  which will be our n bit string.

Plan: Consider a circuit  $C_n$  and convert it to Hamiltonian H s.t.  $\lambda_{min}$  is small iff  $C_n$  has high acceptance probability on some state  $|\psi\rangle$ . We have  $C_n = U_T...U_1$ , each  $U_i$  is a 1 or 2 qubit gate. We'll take error probability to be  $\frac{1}{4T}$ .

The circuit  $C_n$  take n+s+p(n) qubits as inputs. n qubits encode the classical input x. s qubits is the workspace of clean qubits in state 0. p(n) is a witness state  $|\psi\rangle$ .

Given 
$$|\psi\rangle$$
, define  $|\psi_0\rangle = |0\rangle^{\otimes S} |\psi\rangle$ ,  $|\psi_t\rangle = U_t |\psi_{t-1}\rangle$  for  $i = 1...T$ .

We will define the following Hamiltonian which has three parts.

$$H_{init} = \sum_{i=1}^{s} |1\rangle \langle 1|_{i} \otimes |0\rangle \langle 0|_{C}$$

with i = 1, ..., s + p(n)

$$H_{t} = \frac{1}{2} [I \otimes (|t-1\rangle \langle t-1|_{C} + |t\rangle \langle t|_{C}) - U_{t} \otimes |t\rangle \langle t-1|_{C} - U_{t}^{\dagger} \otimes |t-1\rangle \langle t|_{C}]$$

$$H_{final} = |0\rangle \langle 0| \otimes |T\rangle \langle T|_{C}$$

$$H = H_{init} + \sum_{i=1}^{T} H_i + H_{final}$$

We should note that H "follows" the state  $|\psi\rangle$  and percribes penalties when deviating form  $|\psi_0\rangle$ ,  $|\psi_1\rangle$ , ...,  $|\psi_T\rangle$ 

H also penalizes the zero output in the final measurement

C-is an extra "clock" register of log(T+1) qubits

Check locality of H is  $k = \log(T+1) + 2 = O(\log(n))$ . We will further reduce k to a constant k=5.

We need to check that we can distinguish  $x \in L_1$  and  $x \in L_0$  by looking at  $\lambda_{min}$ .

Consider first case when  $x \in L_1$  this means that there exists a p(n) qubit witness state  $|\psi\rangle$  that makes  $C_n$  accept with probability  $1-\frac{1}{4T}$ . Lets consider another state  $|\psi'\rangle = \frac{1}{\sqrt{T+1}} \sum_{t=0}^{T} |\psi_t\rangle |t\rangle$  this is a "History state". When we look at this state we want to evaluate the energy. State  $|\psi'\rangle$  gets no contribution to penalty from  $H_{int}$  and  $H_t$  terms. Since the probability of measuring the (incorrect) outcome is  $\frac{1}{4T}$  we have  $\lambda_{min} = \langle \psi' | H | \psi' \rangle = \frac{1}{T+1} \langle \psi_T | \langle T | H_{final} | T \rangle | \psi_T \rangle \leq \frac{1}{T+1} \langle \psi_T | \psi_T \rangle$  $\frac{1}{T+1}\frac{1}{4T} = a$ . This will show us completeness.

Now take  $x \in L_0$ . We want to show that  $\lambda_{min}$  is at least b = 2a. Consider any witness state  $|\psi'\rangle = \sum_{t=0}^{T} \alpha_t |\phi_t\rangle |t\rangle$  with  $\alpha_t \ge 0$  and  $\{|\phi_t\rangle\}$  is normalised. Then evaluate

$$\langle \psi' | H | \psi' \rangle = \frac{1}{2} (\alpha_{t-1}^2 + \alpha_t^2 - \alpha_{t-1} \alpha_t \langle \phi_t | U_t | \phi_{t-1} \rangle - \alpha_t \alpha_{t-1} \langle \phi_{t-1} | H | \phi_t \rangle)$$
$$\langle \psi' | H | \psi' \rangle = \frac{1}{2} ||\alpha_t | \phi_t \rangle - \alpha_{t-1} U_t |\phi_{t-1} \rangle ||^2$$

We will be making several simplifying assumptions:

- 1) All  $\alpha_t$ s are  $\frac{1}{\sqrt{T-1}}$
- 2)  $|\phi\rangle = |0\rangle^{\otimes s} |\psi\rangle$  for some p(n) qubit state  $|\psi\rangle$
- 3)  $|\phi_T\rangle$  has accept probability close to 1.

Using these and the fact that  $x \in L_0$ ,  $C_n |\phi_0\rangle$  must have accept probability

From 3 we have  $|\phi_T\rangle$  and  $C_n |\phi_0\rangle$  must be nearly orthogonal:

$$1 \leq || |\phi_{T}\rangle - C_{n} |\phi_{0}\rangle || = ||sum_{t=1}^{T} U_{T} - U_{t-1} |\phi_{t}\rangle - U_{T} ... U_{t} |\phi_{t-1}\rangle || \leq \sum_{t=1}^{T} ||U_{T} ... U_{t} |\phi_{t}\rangle - U_{T} ... U_{t} |\phi_{t-1}\rangle || = \sum_{t=1}^{T} ||U_{T} ... U_{t} |\phi_{t}\rangle - U_{T} ... U_{t} |\phi_{t-1}\rangle || \leq \sum_{t=1}^{T} ||U_{T} ... U_{t} |\phi_{t}\rangle - U_{T} ... U_{t} |\phi_{t-1}\rangle || = \sum_{t=1}^{T} ||U_{T} ... U_{t} |\phi_{t}\rangle - U_{T} ... U_{t} |\phi$$

Using the evaluation of  $\langle \psi' | H | \psi \rangle$  assumptions and setting  $\alpha_t = \alpha_{t-1} = \frac{1}{\sqrt{T+1}}$ :

$$\left\langle \psi' | \, H \, | \psi' \right\rangle = \sum_{t=1}^{T} \left\langle \psi' | \, H_{t} \, | \psi' \right\rangle = \frac{1}{2} \frac{1}{T+1} \sum_{t=1}^{T} || \, | \phi_{t} \rangle - U_{t} \, | \phi_{t-1} \rangle \, ||^{2} \\ \geq \frac{1}{2T(T+1)} (\sum_{t=1}^{T} || \, | \phi_{t} \rangle - U_{t} \, | \phi_{t-1} \rangle \, ||^{2}) \\ \geq \frac{1}{2T(T+1)} (\sum_{t=1}^{T} || \, | \phi_{t} \rangle - U_{t} \, | \phi_{t-1} \rangle \, ||^{2}) \\ \geq \frac{1}{2T(T+1)} (\sum_{t=1}^{T} || \, | \phi_{t} \rangle - U_{t} \, | \phi_{t-1} \rangle \, ||^{2}) \\ \geq \frac{1}{2T(T+1)} (\sum_{t=1}^{T} || \, | \phi_{t} \rangle - U_{t} \, | \phi_{t-1} \rangle \, ||^{2}) \\ \geq \frac{1}{2T(T+1)} (\sum_{t=1}^{T} || \, | \phi_{t} \rangle - U_{t} \, ||^{2}) \\ \geq \frac{1}{2T(T+1)} (\sum_{t=1}^{T} || \, | \phi_{t} \rangle - U_{t} \, ||^{2}) \\ \geq \frac{1}{2T(T+1)} (\sum_{t=1}^{T} || \, | \phi_{t} \rangle - U_{t} \, ||^{2}) \\ \geq \frac{1}{2T(T+1)} (\sum_{t=1}^{T} || \, | \phi_{t} \rangle - U_{t} \, ||^{2}) \\ \geq \frac{1}{2T(T+1)} (\sum_{t=1}^{T} || \, | \phi_{t} \rangle - U_{t} \, ||^{2}) \\ \geq \frac{1}{2T(T+1)} (\sum_{t=1}^{T} || \, | \phi_{t} \rangle - U_{t} \, ||^{2}) \\ \geq \frac{1}{2T(T+1)} (\sum_{t=1}^{T} || \, | \phi_{t} \rangle - U_{t} \, ||^{2}) \\ \geq \frac{1}{2T(T+1)} (\sum_{t=1}^{T} || \, | \phi_{t} \rangle - U_{t} \, ||^{2}) \\ \geq \frac{1}{2T(T+1)} (\sum_{t=1}^{T} || \, | \phi_{t} \rangle - U_{t} \, ||^{2}) \\ \geq \frac{1}{2T(T+1)} (\sum_{t=1}^{T} || \, | \phi_{t} \rangle - U_{t} \, ||^{2}) \\ \geq \frac{1}{2T(T+1)} (\sum_{t=1}^{T} || \, | \phi_{t} \rangle - U_{t} \, ||^{2}) \\ \geq \frac{1}{2T(T+1)} (\sum_{t=1}^{T} || \, || \phi_{t} \rangle - U_{t} \, ||^{2}) \\ \geq \frac{1}{2T(T+1)} (\sum_{t=1}^{T} || \, || \phi_{t} \rangle - U_{t} \, ||^{2})$$

# 13 Example Sheet 2

For Grover's search algorithm, if the indicator function is classically efficiently computation then  $I_G$  is implementable.

Look at Part II course to recap coprimality conditions. They also show that to factor N with constant probability it suffices to find an approximations  $\xi$  to  $\frac{s}{r}$  for random s in range [1,r) to 2m+1 binary digits of accuracy with  $m=O(\log N)$  s.t.  $|\xi-\frac{s}{r}|<\frac{1}{2N^2}$ 

# 14 Lecture 13

The Hamiltonian is not k-local for k constant it is of  $\log N$  we will fix this today.

### 14.0.1 Making Hamiltonian k-local

What is responsible for this non-locality effect is the clock-register. So we can reduce locality from  $O(\log n)$  to 5 by making the clock unary. For t=0 we will write  $|0\rangle^{\otimes T}$  for t=1 we write  $|1\rangle |0\rangle^{\otimes T-1}$  for t=3 we write  $|11\rangle |0\rangle^{\otimes T}$ . When we had a binary format we used  $\log(T)$  qubits in the unary format we will use T qubits. This construction will require an extra bit that we will see shortly. Let's crite  $C_t$  to denote the t-th qubit of the clock register. Now our hamiltonian H has the following terms:

$$H_{init} = \sum_{i=1}^{s} |1\rangle \langle 1|_{i} \otimes |0\rangle \langle 0|_{C_{i}}$$

Need an extra term to penalise clock registers that do not confrom to the unary format:

$$H_{clock} = \sum_{t=1}^{T-1} |01\rangle \langle 01|_{C_t C_{t+1}}$$

$$H = H_{init} + \sum_{t=1}^{T} H_t + H_{final} + H_{clock}$$

So each term accesses at most 5 qubits, so this is a 5 local hamiltonian. What is the smallest k and still retain QMA completeness. k-local Hamiltonian is QMA-complete for k=2.

#### 14.0.2 QMA-complete problems

Fairly exhaustive list in a paper by A Bookatz arxiv 1292.6312. We won't really look through many of them but we will think about a few interesting problems.

**Non-identity check**: Given a poly(n)-sized quantum circuit U on n qubits. We want to determine whether U is not close to the identity up to some global phase:

- 1) accept if  $\forall \phi \in [0, 2\pi] ||U e^{i\phi}I|| \ge b$  or
- 2) reject if there exists  $\phi in[0, 2\pi]$  s.t.  $||U e^{i\phi}I|| \le a$ .

This comes with a promise that one of the above options holds and  $b-a \geq \frac{1}{poly(n)}$ 

**Excited k-local hamiltonian**: For some fixed constant c and given a k-local Hamiltonian H on n qubits. We want to determine whether:

- 1) ACCEPT: The c-th eigenvalue of H is  $\leq a$
- 2) REJECT: the c-th eigenvalue of H is  $\geq b$

Have promise that one of these cases hold and  $b-a \geq \frac{1}{(n)}$ .

Many Hamiltonians that correspond to real physical systems are QMA-complete so studying the behavoir of these systems is also hard on a quantum computer. An example of this is a 2-local Ising model  $H_{zzxx} = \sum_{i=1}^n h_i Z_i + \sum_{i=1}^n \Delta_i X_i + \sum_{i,j=1}^n J_{ij} Z_i Z_j + \sum_{i,j=1}^n K_{ij} X_i X_j$  for real  $h_i, \Delta_i, J, K$ . Also true for 2D Heisenberg model or the Rose-Hubbard model.

# Group non-membership: Problem that is in QMA, but not known to be QMA-complete:

Given a finite group G, subgroup H < G, an element  $g \in G$  and determine whether: 1)  $g \not\in H$  2)  $g \in H$ 

### 14.1 Harrow-Hassidian-Lloyd (HHL)

Quantum algorithm for solving linear systems of equations. We want to solve the linear system of equations: Ax = b,  $b, x \in \mathbb{C}^N$  where dimensions N is potentially very large  $N = 2^n$  or let  $2^n$  be the least power of 2 that is greater than N. Rather than outputting the full solution x itself which would take at least O(N) time we want to compute suitable approximations to the value of properties of the solution x such as quadractic expression  $x^T Mx$  (e.g. the total weight of some subset of components).

Very large systems becoming increasingly important in applications:

- data mining/machine learning on datasets that are very large (petabytes...) to discover pattern properties within the data
- numerical solutions to PDEs. Use discretisation techniques (finite element methods). These techniques lead to systems that have the size for larger than

the original problem description.

The best known classical techniques take poly(N) times to solve such problems - in general even to compute the properties of the solutions, no better method is known than computing solutions itself.

Important parameters/for both classical and quantum algorithms

- the system size N
- the desire approximation tolerance  $\epsilon$
- the condition number  $\kappa$  of the matrix A is defined as  $\kappa = \left| \frac{max}{\lambda_{min}} \right|$  (this tells us how close A is to being non-invertible)

If we renormalise A to have  $\lambda_{max} = 1$  then  $_{min} = \frac{1}{\kappa}$ . This is important as the numerical comptation of  $A^{-1}$  becomes less stable with increasing  $\kappa$ . This means we need more significant digits of computation.

# 15 Lecture 15

This algorithm has been improved a lot of times, including as recently as last year. So the exponential improvement in error scaling from  $\frac{1}{\epsilon}$  to  $\log \frac{1}{\epsilon}$  was shown in Childs, Kothan, Smomann in 2015.

We have an N-dimensional space of  $n = \log N$  qubits. Let  $\{|i\rangle : i = 0,...,N-1\}$  be the computational basis. Let eigenvectors and corresponding eigenvalues of A be:  $|u_j\rangle$  and  $\lambda_j$  for j=0,...,N-1. We being by implementing RHS  $|b\rangle$  and consider it in the eigenbasis of A,  $|b\rangle = \sum_{i=0}^{N-1} b_i |i\rangle = \sum_{i=0}^{N-1} \beta_j |u_j\rangle$ . Then  $|x\rangle = A^{-1} |b\rangle = \sum_{j=0}^{N-1} \beta_j \frac{1}{j} |u_i\rangle$ . This transformation of  $|b\rangle$  is linear but not unitary so it cannot be implemented as a quantum operation. To achieve it probabilitistic we will use the phase estimation alogirhtm for  $U=e^{iA}$  with the exponential (and all its powers need for PE) implemented by Hamiltonain simuluation. Then a post-selected measurement process will achieve the descired unitary transformation.

PE with hamiltonian simulation (both assumed to work perfectly) give

$$\left|b\right\rangle \left|0\right\rangle \rightarrow\sum_{j}\beta_{j}\left|u_{j}\right\rangle \left|\lambda_{j}\right\rangle$$

Next we will adjoin an extra ancilla qubit  $|0\rangle$  and apply the controlled rotation (controlled by  $|j\rangle$  register)

C-rotation does

$$|\lambda_j\rangle |0\rangle \rightarrow \sqrt{1 - \frac{c^2}{\lambda_j^2}} |\lambda_j\rangle |0\rangle + \frac{c}{\lambda_j} |\lambda_j\rangle |1\rangle = |j\rangle (\cos\theta_j |0\rangle + \sin\theta_j |1\rangle)$$

here c is chosen to have  $c \leq \min_j |\lambda_j|$   $(c = \frac{1}{k} \text{ for definteness}).$ 

 $\theta_i \in (-\frac{\pi}{2}, \frac{\pi}{2})$  and  $\theta_i = \arcsin \frac{c}{\lambda_i}$  which is determined by the content of the  $|\lambda_j\rangle$  register

Controlled roation is a fixed operation on n + 1 qubits (independent of A or b) that can be implemented by a poly(n) circuit of 1 and 2 qubit gates.

Apply controlled rotation:

$$\sum_{j=0}^{N-1}\beta_{j}\sqrt{1-\frac{c^{2}}{\lambda_{j}^{2}}}\left|u_{j}\right\rangle\left|u_{j}\right\rangle\left|0\right\rangle+\beta_{j}\frac{c}{\lambda_{j}}\left|u_{j}\right\rangle\left|\lambda_{j}\right\rangle\left|1\right\rangle$$

We want the part associated with ancilla in state  $|1\rangle$ . This step is the post-selection step.

We measure the ancilla (hoping to get the result 1) with probability

$$p = ||\sum_j \beta_j \frac{c}{\lambda_j} |u_j\rangle |_j\rangle ||^2 = \sum_j |\beta_j \frac{c}{j}|^2 = \frac{1}{k^2} \sum_j |\frac{\beta_j}{\lambda_j}|^2 \ge \frac{1}{k^2}$$

as  $\sum_{j} |\beta_{j}|^{2} = 1$  and  $\frac{1}{\lambda_{j}} \geq 1$ . In this case the post-measurment state is:

$$|\gamma\rangle = \frac{1}{\sqrt{p}}c\sum_{j=1}^{N-1}\frac{\beta_j}{j}|u_j\rangle|\lambda_j\rangle$$

To mitigate the probabilistic nature of this process and obtain the state with any fixed proability with  $1-\eta$  for  $\eta>0$  and  $\eta$  small. We can repeat the whole process  $\frac{\log\frac{1}{\eta}}{p}=O(k^2)$  times and the outcome 1 will occur at least once with probability  $1-\eta$ .

**Fact**: If a signle trial has success probability p and we repeat the trial M times independently then for any  $0 < 1 - \epsilon < 1$  Prob( at least on trial in M succeeds)  $> 1 - \epsilon$  if  $M \ge \frac{-\log \epsilon}{p}$ .

Remark (amplitude amplification): the  $O(k^2)$  repetitions needed to generate  $|\gamma\rangle$  can be improve to O(k) by using AA instead of repreated measurements.

Having obtained  $|\gamma\rangle$  we run the inverse phase estimation process to "uncompute" or "erase" the  $|\lambda_j\rangle$  register resetting it to  $|0\rangle$  and get  $|\hat{x}\rangle = \frac{c}{\sqrt{p}} \sum_{j=0}^{N-1} \frac{\beta_j}{j} |u_j\rangle = \frac{c}{\sqrt{p}} |x\rangle$ 

 $\frac{c}{\sqrt{p}}\left|x\right\rangle \\ Finally, we will perform a measure mnt of observable \mathbf{Mon}|\hat{x}\rangle \text{ to estimate its mean value } \frac{c^2}{p}\left\langle x\right|M\left|x\right\rangle.$ 

According to Chernoff-Hopffling bound (look this up in lecture notes)  $O(\frac{\log \frac{1}{\ell}}{\mathcal{E}^2})$ measurements will suffice to estimate the meant to any desired accuracy  $\xi$  with any probability of success  $1 - \eta$ . Similarly we can estiamte the prob p in the post-selection step to any accuarcy  $\xi$  by applying C-H bound to the ancilla measurment outcome. As the ancilla measurement outcome is a random variable with two outomces 0 and 1, so finding the mean of this variable gives p. Using these values of  $\frac{c^2}{n}\langle x|M|x\rangle$  and  $p,c=\frac{1}{k}$  we get out our value  $M=\langle x|M|x\rangle$ .

#### 16 Lecture 15

Let's stop and think about HHL a little bit more. We have considered the ideal situation but we now want to consider how it is effected by approximation errors

#### 16.0.1Summary analysis of approxiamtion erros and runtime in HHL (non-examinable)

A full analysis of runtime and approximation errors can be found in Physical Review Letters vol 103.150502 (2005).

We assume we can make  $|b\rangle$  correctly (the reason we don't want to focus on this is the issue of creating  $|b\rangle$  is an orthonorous process to HHL and should be considered separately so introduces no errors). Similarly, the controlled rotation C-ROT is a final unitary (it can be efficiently implemented using some universal gate set by Solovay-Kitaev theorem)

We will assume that C-rotations can be done exactly

Convsider the phase estimation part. PE for a unitary U and n qubit lines gives an estimate of  $\lambda$  up to n bits of precision. Denote this estimate  $\lambda'$  and assume that it is the closest n-bit approximation  $\lambda'$  is up to additive error  $\eta = \frac{1}{2^n}$ 

The PE process uses Hamalitonian simulation that requires execution of con-

trolled 
$$U, U^2, U^4, ..., U^{2^{n-1}}$$
 If  $U = e^{iA}$  then  $C - U = e^{i\tilde{A}}$  with  $\tilde{A} = \begin{pmatrix} 0 & | & 0 \\ - & & - \\ 0 & | & A \end{pmatrix}$ .

 $\hat{A}$  remains s-sparse and row-computable if A was.

Thus we need to implement the following  $e^{i\bar{A}t}$  for  $t=1,2,4,...,2^{n-1}$  which is done using hamiltonian simulation for s-spare matrices with circuit size  $O(\log Nts^2)$ . So the total circuit size  $O(\log Nt_0s^2)$  where  $t_0 = 1 + 2 + ... + 2^{n-1} \sim 2^n = \frac{1}{n}$ .

After the controlled rotation  $C_{rot}$  using the  $\lambda$ 's produced and the corresponding measurement in post-selction step the final state is:

$$|\hat{x}'\rangle = \frac{1}{D'} \sum_{j} \frac{c}{\lambda'_{j}} |u_{j}\rangle, D' = \sqrt{p'} = \sqrt{\sum_{j} |\beta_{j}|^{2} \frac{c^{2}}{j}}$$

- $|\hat{x}\rangle$ , D, p will denote the expressions with  $\lambda'$ s replaced by true  $\lambda$  values. Our requirement is that  $|\hat{x}\rangle$  should be wihtin  $\epsilon$  of  $|\hat{x}\rangle$  so we need to choose a suitably large  $t_0 = \frac{1}{\eta}$  ( number of qubit lines). To establish dependenace of  $t_0$  on we need the following facts:
- a) If  $\lambda$  has additive error  $\delta()=\eta$  then  $\frac{1}{\lambda}$  has  $\delta(\frac{1}{\lambda})\sim\frac{1}{\lambda}\delta(\lambda)$  so  $\frac{1}{\lambda}$  has relative error  $\frac{\delta(\frac{1}{\lambda})}{1}=\frac{\eta}{}$
- b) If A' and B' approximate A and B to relative error  $\xi$  then  $\frac{A'}{B'}$  approximates  $\frac{A}{B}$  to relative error  $\xi$  too  $\frac{A'}{B'} = \frac{A(1+O(\epsilon))}{B(1+O(\epsilon))} = \frac{A}{B}(1+O(\epsilon))$
- By (a)  $\frac{1}{\lambda'_j}$  approximates  $\frac{1}{\lambda_j}$  to relative error  $O(\frac{q}{\lambda_j}) = O(Kq)$ . Similarly D' approximates D to relative error O(Kq). Then by (b)  $A' = \frac{\beta_j c}{\lambda'_j}$ ,  $\beta' = D'$  the amplitudes of  $|\hat{x}'\rangle$  approximate those of  $|\hat{x}\rangle$  up to multiplicative error O(Kq).

$$|||\hat{x}'\rangle - |\hat{x}\rangle|| = ||(1 + O(k\eta)|\hat{x}\rangle - |\hat{x}\rangle|| = O(kq) = \epsilon$$

for a given  $\epsilon$  we choose  $kq = \epsilon, t_0 = \frac{1}{\eta}$  we gget  $t_0 = \frac{k}{\epsilon}$  and the PE step (run once) has runtime  $O(\log Nt_0s^2) = O(\log N\frac{k}{q}s^2)$  for  $k = poly\log N$  so runtime is  $O(poly\log N)\frac{1}{\epsilon}$ ).

#### Further comments on erros and runtime:

In the above analysis we assumed that the PE process (with n qubits outputs the closest n-bit approximation  $\lambda'$  to  $\lambda$  whereas it outputs a superposition of all  $\lambda_k = \lambda' \pm \frac{k}{2^n}$  (peaked around k = 0) which is another source of error. TO mitigate this we use PE(b), use more qubit lines to get clsoer approximation with any high probability.

In the HHL paper they use a slighly improved PE which uses non-uniform superposition state at the start. It may be shown that an optimal choice of this non-uniform superposition can lead to a quadratic reduction in the variance of the k-value outcomes. In other works  $\delta(\lambda)$  is more concentrated on smaller vlaues.

In the HHL paper, the authors present an algorithm that can handle all conditioned matrices A by introducing a further qubit "flag" and filtering functions to separate well and ill conditioned parts of A and processing the former.

In the HHL algorithm the PE step is the soruce of the worst accuracy depdnance  $O(\frac{1}{\epsilon}$ . This was ubsequently improved by Berry et all arXiv 1212.1414 from  $O(\frac{1}{\epsilon} \to O(\log(\frac{1}{\epsilon}))$  this technique invovles rpelacing PE process by an entirely different method for implementing  $A^{-1}$ . This new method is called "linear combinations of unitaries" (LCU). Implementing  $A^{-1}$  as a linear combination of unitaries each of which is individually implemented. Then we represent the linear sum in terms of another unitary V on a larger space with extra ancillar qubits.

# 17 Lecture 16

# 17.1 Clifford computations (and classical simulation of QC)

Vague question, what is the key "quantum effect/resource/ingredient" giving quantum computing its benefits or advantage over classical computation.

## 17.1.1 Classical simulation of quantum computation

We are given classical description of the following:

- 1) description of a quantum circuit C which is just a collection of 1 or 2 qubit gates on n qubits with size N = poly(n).
- 2) description of the input (e.g.  $|l_1...l_n\rangle$ ,  $l_k \in \{0,1\}$  or  $|\alpha_1\rangle ... |\alpha_n\rangle$  a general product state)
- 3) cominated output lliines (e.g. 1st line to read off the answer)

We want to by classical means only:

Weak simulation: sample the output distribution (using classical computer with access to calssical randomness)

**Strong** simulation: calculate the output probabilities of  $p_0$  and  $p_1$  (probability of getting 0 or 1 on the first line).

We want to do it efficiently in classical poly(n) time known as "efficient classical simulation".

Note, the quantum process itself gives only a weak simulation "self-simulation" as it only gives a sample of the output distribution.

"Direct" strong classical simulation is always possible but not efficent in general.

However, if we are promised that our state is a product state at every stage then we can give an efficent simulation. Each step we only involve 1 or 2 qubits in known state  $|\alpha_i\rangle |\alpha_j\rangle$ . Consider the 2-qubit case: U acting on 2 qubits with  $|\psi\rangle = \sum_{i_1..i_n} C_{i_1...i_n} |l_1...l_n\rangle$  then in general you get  $U|\psi\rangle = \sum_{i_1}^n V_{l_1 l_2}^{k_1 k_2} C_{k_1 k_2 l_3...l_n} |\beta_{k_1}\beta_{k_2} l_3'...l_n'\rangle$  (let  $\tilde{C}_{i_1...i_n} = V_{l_1 l_2}^{k_1 k_2} C_{k_1 k_2 l_3...l_n}$ . When  $|\psi\rangle$  is a product state we have  $C_{i_1...i_n} = a_{i_1}b_{i_2}c_{i_3}...c_{i_n}$  so  $\tilde{C}_{i_1...i_n} = (\sum_{i_1i_2} U_{i_1i_2}^{k_1 k_2} a_{k_1}b_{k_2})c_{i_3}...c_{i_n}$ . so we can factorise back into the form  $f_1...f_n$ . Sometimes it was claimed that entanglement was the secret ingredient as it looks like if you take away entanglement then you can classically simulate the system. However, this argument shows that if we have no entanglement then we get no quantum advantage. This shows that entanglement is necessary but not sufficent for quantum advantage. "Most entangled states are too entangled to be useful as quantum computational resources" Fross, Flanimma, Eisert Phys Rev Letter 107 1905 01 '99. Universal quantum computation with little entanglement Physical Rev Letter 110 060504 '13. So the story with entanglement is complicated. so we are going to consider something else called Clifford computations.

#### Clifford Computation

Prelminary defintions: n-qubit Pauli group  $\mathcal{P}_n = \{\pm 1, \pm iP_1 \otimes ... \otimes P_n\}$  with each  $P_i \in \{I, X, Y, Z\}$ .

Define a Clifford operation C on n qubits as such that conjugates Paulis to paulis.  $\forall P \in \mathcal{P}_n CPC^{\dagger} \in \mathcal{P}_{\backslash}$ . So the Clifford group is the normalizer of subgroup  $\mathcal{P}_n$  in  $U(2^n)$ .

Clifford circuits found many applications:

- 1) the theory of quantum error correction ('stabilizer codes')
- 2) insights into quantum computational power
- 3) verification of quantum computations

Examples: All Paulis are clifford

$$H, S = \begin{pmatrix} 1 & 0 \\ 0 & i \end{pmatrix}$$
 phase  $(\frac{\pi}{2})$  gate  $CX$  2 qubit gate  $Z = S^2$  and  $X = HZH$ 

 $SWAP_{12} = CX_{12}CX_{21}CX_{12}$ 

Full explict characterisation of Clifford gates:

**Theorem:** C on n qubits is Clifford iff C is a circuit of H, S, CX gates

Clifford computation circuits that only involve clifford gates.

Note: Generating entanglement  $|0\rangle |0\rangle \to^H \frac{1}{\sqrt{2}} (|0\rangle + |1\rangle) |0\rangle \to^{CX_{12}} \frac{1}{\sqrt{2}} (|0\rangle |0\rangle +$  $|1\rangle |1\rangle$ 

**Grottesman-Knull Theorem**: Let C be any Clifford circuit on n qubits, N = poly(n) with input any product state  $|a_1\rangle \dots |a_n\rangle$  and output a measurment on the 1st line. Then the output can always be classically strongly efficently simulated.

#### 18 Lecture 17

**Proof**: The idea is that instead of evolving the input state  $|\psi_{in}\rangle$  to a final state  $|\psi_{final}\rangle$  for measurement we will back propagate the final Z measurement.

Let's write the measurement  $Z_1 = Z \otimes B \otimes ... \otimes B$  and recall that Z = $|0\rangle\langle 0|-|1\rangle\langle 1|=\Pi_0-\Pi_1$ . So we can write the difference of the probabilities as  $p_0 - p_1 = \langle \psi_{final} | Z_1 | \psi_{final} \rangle = \langle \psi_{in} | C^{\dagger} Z_1 C | \psi_{in} \rangle$ . Let  $C = C_N C_{N-1} ... C_1$  with  $C_i \in \{H, S, CX\}$  so:

$$p_0 - p_1 = \langle \psi_{in} | C_1^{\dagger} ... C_N^{\dagger} Z_1 C_N ... C_1 | \psi_{in} \rangle$$

Each conjugation gives us a Pauli, and thus we get:

$$p_0 - p_1 = \langle \psi_{in} | \tilde{P}_n \otimes ... \otimes \tilde{P}_n | \psi_{in} \rangle$$

All update rules for Pauli products by H, S, CX conjugations are easy to compute ( we have only 1 and 2 gate computations so it is a constant size computation)

Finally we make use of the fact that our input state is a product state. Therefore:

$$p_0 - p_1 = \prod_{i=1}^n \langle a_i | \tilde{P}_i | a_i \rangle$$

Each of these terms is a product of n2x2 matrix calculations so is efficiently computable. We just need to remember that they are normalised so  $p_0 + p_1 = 1$  so we get  $p_0, p_1$  values via efficient classical computation. So Cliffords are easy.

Let us extend unitary Clifford circuits to allow the inclusion of intermediate measurements (1-qubit) measurements

Distinguish two scenarios, non-adaptive or adaptive (if we change future gates or measurements depending on the outcome of measurements). Non-adaptive: Gates cannot depend on earlier measurement outcomes. Adaptive: Gates can depend on earlier measurement outcomes.

**Theorem** Let C be any Clifford circuit with intermediate measurements with an arbitary product state input and a single line output. Then:

- a) If C is non-adaptive then the output is classically efficiently strongly simulatable
- b) If C is adaptive then full universal quantum computation is possible

Note: if input is restricted to the computational basis states then (b) remains classically weakly efficiently simulatable (so b) relies on the arbitary product state input for its universality). R. Sazsen M van den Nest arXic 1305.6190 (2013).

Proof: for a) we observe that non-adaptive Clifford circuits are reducible to the full unitary case, and then use GK. To see how this reduction works consider a measurement between  $C_1$  and  $C_2$  on one of the qubit lines with  $C_2$  independant of measurement outcome. This is equivalent to introducing an extra qubit ancilla which the measurement is swap for a CX on the ancilla, as then measuring the ancilla at any point would give the original circuit but you don't even need to measure it you can just discard it.

for b) we will use a Bravyi-Kitaev idea of so called magic states. Fact: Clifford gates with 1-qubit (non-Clifford) T gate with

$$T = \begin{pmatrix} 1 & 0 \\ 0 & e^{-\frac{\pi}{4}} \end{pmatrix} = \sqrt{s} = \text{phase } \frac{\pi}{4} \text{gates}$$

is universal. In fact if you add anything to the Clifford group that is non clifford you regain universal quantum computation. Let  $|A\rangle$  denote 1-qubit state  $|A\rangle = \frac{1}{\sqrt{2}}(|0\rangle + e^{i\frac{\pi}{4}}|1\rangle)$ .

We will implement a T gate using  $|A\rangle$  ancilla and an adaptive Clifford circuit

(T-gadget).

If we state with any n-qubit state plus the ancilla  $|A\rangle$ . If we apply a CX from any line to the ancilla and then measure and discard the ancilla. If we get 1 in this measurement we apply an S gate to the line we applied the CX from. This generates a T gate (up to overall phase). Easy to check (ES3) this gives an adaptive Clifford circuit that always implements T gate (up to overall phase that depends on m) on line k and probabilities are  $\frac{1}{2}$ ,  $\frac{1}{2}$  regardless of output  $|\psi\rangle$ 

Remark: To give increase in computational power we could:

- 1) allow inclusion of new gates (T-gate)
- 2) add extra "magic" states as extra inputs

For Clifford circuits with product state inputs, single line outputs and intermediate measurements

- a) non-adaptive -; classicaly strongly efficiently simulatable
- b) adaptive -; full universal quantum computing power
- a)  $C_0M(i_1,y_1)C_1M(i_2,y_2)C_2...$ b)  $C_0M(i_1,y_1)C_1(y_1)M(i_2,y_2)C_2(y_1,y_2)$  ...

# 19 Lecture 17

# 19.1 Quantum Error Correction/Stabilizer formalism

Classical computers are incredibly reliable about  $10^{-9}$  accuracy. This is not the case for quantum hardware, the gate fidelities are about 99.9% nowhere near classical computers. So if you ran a thousand gates in a row then as the errors are multiplicative the final state gets very perturbed, so we need to think about error correction. Our go to place is classical intituion, how is this done classically?

Consider a single classical bit x that we want to store. Let's assume that it gets flipped with probability p. One way to improve the probability of reading the correct value is instead of storing x we store xxx. Readout  $y = y_1y_2y_3$ . The wrong answer will be returned wih probability  $3p^2(1-p) + p^3 = p^2(3-2p) \sim O(p^2)$  so for  $p \in (0, \frac{1}{2})$   $p^2(3-2p) < p$ . So this map  $x \to x_1x_2x_3$  is an error correcting code.

Now let's talk about quantum errors. Not talking about a classical value but rather a state. We want to preserve a qubit  $|\psi\rangle$ . So the process of error correction can be described as:

$$|\psi'\rangle = DNE |\psi\rangle |0\rangle^{\otimes n}$$

with E is the unitary encoding, N is the noise operation, and D is decoding operation. The aim is to obtain  $|\psi'\rangle \approx |\psi\rangle$  for a given set of noise operations.

**Attempt 1**: Measure  $|\psi\rangle$  in the computational basis to get 0 or 1. Then encode using classical repetition code. This is not a good ida as we destroy superposiiton/entanglement

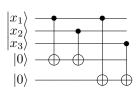
**Attempt 2**: Engineer a map  $|\psi\rangle \rightarrow |\psi\rangle |\psi\rangle |\psi\rangle$  this violates the no-cloning so isn't possible

Different approach (that works!). This is a 2-step process:

1. Encode  $|\psi\rangle=\alpha\,|0\rangle+\beta\,|1\rangle$  as  $|E(\psi)\rangle=\alpha\,|000\rangle+\beta\,|111\rangle$ . This is NOT the same as cloning. Encoding:

$$\begin{array}{c|c} \alpha \mid 0 \rangle + \beta \mid 1 \rangle & & \\ \mid 0 \rangle & & \\ \mid 0 \rangle & & \\ \end{array}$$

Decoding:



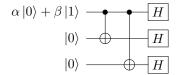
For any input superposition of the form  $\alpha |x_1x_2x_3\rangle + \beta |x_1x_2x_3 \oplus 111\rangle$  the decoding circuit performs the map  $(\alpha |x_1x_2x_3\rangle + \beta |x_1x_2x_3 \oplus 111\rangle) |0\rangle |0\rangle \rightarrow (\alpha |x_1x_2x_3\rangle + \beta |x_1x_2x_3 \oplus |11\rangle) |x_1 \oplus x_2\rangle |x_2 \oplus x_3\rangle$ 

If we measure two output qubits we learn  $x_1 \oplus x_2$ ,  $x_1 \oplus x_3$  without disturbing the original state. The result of measuring these output qubits is called the syndrome. If we assume only one bit has been flipped the syndrome tells us which one as they have all distinct syndromes. After we dectect a bit-flip error (on a single qubit) we apply the same bit-flip operation to that qubit and restor the original encode state  $\alpha |000\rangle + \beta 111 \rightarrow \alpha |0\rangle + \beta |1\rangle$  by reverse the original encoding circuit.

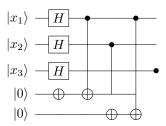
Consider the effect of a Z (or "phase") error. It maps the encoded state on a single qubit:

$$\alpha |000\rangle + \beta 111 \rightarrow \alpha |000\rangle - \beta |111\rangle$$

would have Syndrome measurement 00, so we cant even detect a Z error. However, they can be detected using a different code. If we notice that Z = HXH, so Z is effectively X but in a different basis the  $|+\rangle$   $|-\rangle$  basis. We can use a different code  $|\psi\rangle = \alpha |+++\rangle + \beta - --$  so what is our new encoding circuit:



Decoding:



A code of this kind can protect against Z errors. But there is a problem as this code can no longer protect against X errors.:

$$\alpha \left| 0 \right\rangle + \beta \left| 1 \right\rangle \rightarrow \frac{1}{2\sqrt{2}}\alpha(\left| 0 \right\rangle + \beta \left| 1 \right\rangle)(\left| 0 \right\rangle + \left| 1 \right\rangle))(\left| 0 \right\rangle + \left| 1 \right\rangle) + \beta(\left| 0 \right\rangle - \left| 1 \right\rangle)(\left| 0 \right\rangle - \left| 1 \right\rangle)(\left| 0 \right\rangle - \left| 1 \right\rangle) \rightarrow \frac{1}{2\sqrt{2}}\alpha(\left| 000 \right\rangle + \left| 111 \right\rangle)(\left| 0 \right\rangle - \left| 1 \right\rangle$$
(\left| 0 \right\rangle - \left| 1 \right\rangle

This is a 9-qubit code.

To decode, we first apply the decoding circuit for bit-flip error to each block. Assuming that at most one bit flip error has occured in each block the resulting state will be  $\alpha \mid +++\rangle + \beta \mid ---\rangle$  (with at most one phase flip error). We map this state to  $\alpha \mid 0 \rangle + \beta \mid 1 \rangle$  using the decoding algorithm for the phase flip code.

So with a great amount of effort we managed to protect against only two types of errors. It turns out this approach works for all errors, this is because the matrices  $\{I, X, Y, Z\}$  with Y = iXZ form a basis for the complex vector space of 2x2 matrices, so an arbitary error acting on a single qubit can be written as a linear combination of I, X, Z and iXZ.

Theorem: Quantum error correcting criteria

# 20 Example sheet 3

$$||A|| = \max_{|\psi\rangle} |A|\psi\rangle|$$

Here  $|A|\psi\rangle$  | is not  $\bar{\psi}A|\psi\rangle$  but just the magnitude of the vector created. So try  $|\psi\rangle=\begin{pmatrix}x\\y\end{pmatrix}$  with  $|x|^2+|y|^2=1$  so:

$$||A|| = \max_{x^2 + y^2 = 1} |A \begin{pmatrix} x \\ y \end{pmatrix}|$$

A can be expressed as a linear combination of unitary matrices of the form  $U_{ab} = X^{d-b}Z^{d-a}$  with  $X|k\rangle = |k+1 mod d\rangle$  and  $Z = e^{\frac{2\pi i k}{d}}|k\rangle$ , as these are linearly independant and unitary.

There is a different method for finding a linear combination of unitaries to represent a non-unitary matrix (elaborated on in the solutions). Use SVD:

$$A = UDV$$

with D real and diagonal. It suffices to make any such D as  $D = \sum_i c_i U_i$ .

Need to check A is well conditioned when performing HHL on Ax = b as we need to take its inverse (so need to check it has no zero eigenvalues). In example sheet this is shown by knowing  $||C|| \leq \frac{1}{2}$  and A = I - C so eigenvalues must lie in  $[\frac{1}{2}, \frac{3}{2}]$ . For HHL we also require A to be row-spares and row computable, and for  $|b\rangle$  to be implemented able in poly(log N) time.

If you ever need to find out how close to states are, you use the swap test!!!

# Chernofl-Hoelfing bound

If we westimate Pr(0) as frequency f of 0's seen in k samples of a 0/1 distribution, then:

$$|f - Pr(0)| < \xi$$
, with probability  $\geq 1 - \epsilon$ 

if  $k \geq \frac{\log \frac{1}{\epsilon}}{2\xi^2}$ 

# 21 Lecture 21

# 21.1 Alternative approach to proving GK theorem

We will discuss unitary part in lectures. The part that deals with measurements will be discussed in ES4.

We will start with the stabiliser state and update the list of generators for each Clifford gate applied to it.

Let;s start with something concrete. Starting with the list of zero states  $|0\rangle^{\otimes n}$  we can write the list of generators of this state as  $\langle Z|I,IZII,...,I|Z\rangle$ .

We will introduce the stabilizer tableau, giving two  $n \times n$  matrices + extra space



 $3) Apply {\tt CNOT} (i-the Sitwise {\tt XOR} the i-the olumn of the X matrix into the j-the olumn of the X matrix, and then bitwise {\tt XOR} the i-the olumn of the interval into the j-the olumn of the interval interv$ 

 $Bitwise ext{XOR} the ext{i}-th column of$ 

Let's define a model:

Input: a) Classical part description of a quantum circuit, input state and measurement lines, b) Quantum part input  $|0\rangle^{\otimes n}$  Computation described as:

$$M_2C_2M_1C_1\ket{0}^{\otimes n}$$

 $M_1$  gives a classical string  $x_1, ..., x_n$  that we feed into classical processing, and the output of the classical processing affect the block of gates  $C_2$ . This is repeated again and again for  $C_3$ ,  $M_3$  etc. etc.

 $C_1, C_2, \dots$  are just sequences of 1 and 2 qubit gates.

We will distinguish intermediate states after each measurement  $M_n$  as  $|\psi_n\rangle$  with input state as  $|\psi_0\rangle$ .

We will distinguish two flavours of this type of computation, which essentially depends on what happens with the intermediate states.

Type 1: The input state is carried through the computation. In other words we have  $|\psi_{out}\rangle = M_N C_N ... M_1 C_1 |\psi_0\rangle$ 

Type 2: Every round (consisting of application of  $C_i$  followed by  $M_i$ ) we make use of  $|\psi_{i-1}\rangle = |0\rangle^{\otimes n}$ .

Example of hybrid quantum-classical computation of type 1) Pauli-based computation (Bravyi-smith, smoh '16) covered this lecture, type 2) variational quantum eigensolver (last lecture of this course).

#### Pauli-based computation

Start with an input state  $|\psi_{in}\rangle$  on t qubits

Apply a sequence of pairwise commuting Pauli-measurements  $P_1...P_s \in \mathcal{P}_t$  sequentially to the input state

The choice of each  $P_i$  can in general depend on previous measurement outcomes. When there is no such dependancy we say that PBC is non-adaptive.

Considering the adaptive case. There is a nice tree diagram in the notes that I can't replicate in latex.

**Theorem:** Let circuit C be any generally adaptive quantum circuit on n+tqubits with input state  $|\psi\rangle = |\sigma\rangle \otimes |\rho\rangle$  where  $|\sigma\rangle$  is a stabilizer state on n qubits,  $|\rho\rangle$  is any state on t qubits. Suppose that the unitary steps of C are all Clifford gates. Then:

- i) C may be weakly simulated by a generally adaptive PBC circuit  $\tilde{P}_1...\tilde{P}_s$  on t qubits and with  $s \leq t$  steps.
- ii) If C is non-adaptive (with final measurement outputs) then the corresponding PBC circuit  $P_1, ... P_S$  can be chosen to be non-adaptive
- iii) If some Z measurement in C are postseleted to outcome +1, then this circuit can be weakly simulated by a PBC circui in which some of the Pauli measurements are postselected to outcome +1

**Lemma A**: Let  $P,Q \in \mathcal{P}_n$  be anti-commuting Pauli operators. let  $|\psi\rangle$  be an eigenstate of P with  $P|\psi\rangle =_p |\psi\rangle$ ,  $\lambda_p = \pm 1$ . Then 1) a measurement of Q on  $|\psi\rangle$  yields  $\lambda_a=\pm 1$  with equal probabilities.

- 2) The operator  $V(\lambda_P, \lambda_Q) = \frac{\lambda_p P +_q Q}{\sqrt{2}}$  is always a unitary Clifford operation. 3)  $V(\lambda_p, \lambda_q)$  is the normalised projection of  $|\psi\rangle$  onto the Q eigenspace of Q.

Remark: Hence measurement of Q on  $|\psi\rangle$  is equivalent to classically chossing a uniformly random  $\in \{\pm 1\}$  and applying clifford unitary.

#### 22 Lecture 23

Lemma B: For any

$$P = \pm A_1 \otimes ... \otimes A - n \otimes B_1 \otimes ... \otimes B_t \in \mathcal{P}_{n+t}$$

, with  $A_1, B_1 \in \{X, Y, Z, I\}$  write:

$$\tilde{P} = \pm B_1 \otimes ... \otimes B_t \in P_t$$

(with the same oeveral sign as P).

If P commutes with  $Z_1...Z_n \in \mathcal{P}_{n+t}$ , then each  $A_i$  is either Z or I. If for all i, each  $A_i$  is either Z or I, then for any t-qubit state  $|\psi\rangle$  the measurement of P on  $|0\rangle^{\otimes n} |\psi\rangle$  and the measurement  $\tilde{P}$  on  $|\psi\rangle$  give the same output distribution, and corresponding post-measurement states:  $|0\rangle^{\otimes n} |\psi'\rangle, |\psi'\rangle$  respectively, with he same t-qubit state  $|\psi'\rangle$ .