Recall this version of the PCP theorem from last lecture.

**Theorem 1.** There exists an alphabet  $\Sigma$  and a constant  $\varepsilon > 0$  for which the following task is NP-hard: Given a satisfiable 2CSP instance over  $\Sigma$ , find an assignment that satisfies a  $1 - \varepsilon$  fraction of constraints.

In a general 2CSP instance, a variable may be present in an arbitrary number of constraints. What if we restrict our attention to instances where every variable appears in at most d constraints, where d is small compared to the number of variables? When d=1, every variable appears in one constraint and finding a satisfying assignment is easy. When d=2, the task is a bit harder but still solvable in time linear in n. On the other hand, when d is as large as the number of constraints the problem becomes NP-hard. This suggests that the problem may become gradually harder as d gets larger.

It turns out that this intuition is incorrect:

**Theorem 2.** There exists an alphabet  $\Sigma$  and constants d and  $\varepsilon$  such that given a satisfiable 2CSP instance over  $\Sigma$  where every variable appears in at most d constraints, it is NP-hard to satisfy a  $1 - \varepsilon$  fraction of the constraints.

We prove this statement by reduction from Theorem 1. Let  $\Phi$  be the 2CSP instance in question. We want to construct a new instance  $\Phi'$  which is as hard as  $\Phi$ , but every variable appears in at most d constraints. Some of the variables in  $\Phi$  may appear in more constraints. If variable  $x_i$  appears in  $n_i$  different constraints it is natural to replace it with  $n_i$  new variables  $x'_{i1}, \ldots, x'_{in_i}$  and impose some additional constraints that force all of  $x'_{i1}, \ldots, x'_{in_i}$  to take the same value.

The first thing we may try is to add the constraints  $x'_{i1} = x'_{i2}, x'_{i2} = x'_{i3}, \dots, x'_{i(n_i-1)} = x'_{in_i}$  to  $\Phi'$ . Then if  $\Phi$  has a satisfying assignment, the assignment obtained by setting  $x'_{i1} = \dots = x'_{in_i} = x_i$  will be satisfying for  $\Phi'$ . Suppose that we could then find an assignment x' that satisfies a  $1 - \varepsilon'$  fraction of its constraints of  $\Phi'$ . Can we use x' to obtain an assignment that satisfies most constraints in  $\Phi$ ?

It is not hard to see that if  $\varepsilon' = 0$ , the assignment  $x_i = x'_{i1} = \cdots = x'_{in_i}$  is satisfying for  $\Phi$ . However, even if one of the equality constraints is violated, the values of  $x'_{ij}$  could split into two equally sized sets. Then it is not clear which value to assign to  $x_i$  and it is possible to come up with examples where no matter which value we assign, a large fraction of the constraints of  $\Phi$  will be violated.

So we need to make the equality constraints more robust: If there is no clear majority among the values  $x'_{i1}, \ldots, x'_{in_i}$ , then not one but many of the equality constraints should be violated. One way to do so is to impose the equality constraint  $x'_{ij} = x'_{ij'}$  for every pair j < j'; but then we have done nothing about reducing the number of constraints a variable appears in.

In general the equality constraints we are looking for can be described by an undirected graph G on the vertices  $\{1,\ldots,t\}$ . An assignment to  $x'_{i1},\ldots,x'_{in_i}$  can be viewed as a partition of the vertices into sets  $A_{\sigma} = \{j : x'_{ij} = \sigma\}$ , where  $\sigma$  ranges over  $\Sigma$ . On the one hand, we want the degree of this graph to be constant. On the other hand, we want that a partition  $\{A_{\sigma}\}$  splits many of the edges of G, unless one of the sets  $A_{\sigma}$  contains most of the vertices.

Both of these properties are achieved by expander graphs. To understand expander graphs and their properties we first need to take a detour into random walks, adjacency matrices, and eigenvalues.

In what follows we will assume the graph G is undirected, connected, and d-regular.

## 1 Adjacency matrix and eigenvalues

Suppose a particle sits at a vertex s of some graph G. At every step, s moves to a random one of its neighbors. How long will it take s to reach a vertex in G that looks random and independent of s?

To answer this question, it will be helpful to represent the random walk by a sequence of probability distributions  $\mathbf{p}^0, \mathbf{p}^1, \ldots$  on the vertices of G, with the following interpretation: At each step t,  $\mathbf{p}^t(u)$  is the probability of the particle ending up at vertex u after t steps of the walk. Initially, we have  $\mathbf{p}^0$  assign probability 1 to vertex s, and probability 0 to all the other vertices. The distribution  $\mathbf{p}^{t+1}$  can be calculated from  $\mathbf{p}^t$  via the formula

$$\mathbf{p}^{t+1}(u) = \sum_{v:(v,u) \text{ is an edge}} \frac{1}{d} \cdot \mathbf{p}^t(v). \tag{1}$$

We are now interested in the following question: When t gets large, how close does the distribution  $\mathbf{p}^t$  get to the uniform distribution  $\mathbf{u}$  on the set of vertices? To answer this question, we need some way of measuring how "close" two distributions are. In our setting the most convenient measure is the  $\ell_2$  norm. The  $\ell_2$  norm of a vector  $\mathbf{v}$  is the quantity

$$\|\mathbf{v}\| = \left(\sum_{i} \mathbf{v}_{i}^{2}\right)^{1/2}$$

and the  $\ell_2$  distance between two vectors  $\mathbf{v}$  and  $\mathbf{v}'$  is the  $\ell_2$  norm of  $\mathbf{v} - \mathbf{v}'$ . We will think of probability distributions as vectors in  $\mathbb{R}^n$  (with one entry for each vertex in the graph), and we will say that two distributions  $\mathbf{p}$  and  $\mathbf{p}'$  are  $\epsilon$ -close (in  $\ell_2$  distance) if  $\|\mathbf{p} - \mathbf{p}'\| \le \epsilon$ .

The (normalized) adjacency matrix of G is an  $n \times n$  matrix A defined as follows:

$$A_{u,v} = \frac{\text{number of edges between } u \text{ and } v \text{ in } G}{d}$$

This matrix is symmetric and the entries in each row add up to one. Using A, we can write equation 1 in matrix form as  $\mathbf{p}^{t+1} = \mathbf{p}^t A$  (it is customary to represent  $\mathbf{p}^t$  as row vectors) and so we immediately obtain that  $\mathbf{p}^t = \mathbf{p}^0 A^t$ .

The eigenvalues and eigenvectors of A play a significant role in determining the behavior of random walks on G. Recall that an eigenvalue-eigenvector pair is a complex number  $\lambda$  and a vector  $\mathbf{v}$  such that  $\mathbf{v}A = \lambda \mathbf{v}$ . It is a basic theorem in linear algebra that symmetric matrices have an orthonormal basis of eigenvectors with real eigenvalues. Let's denote these pairs by  $(\lambda_1, \mathbf{v}_1), \ldots, (\lambda_n, \mathbf{v}_n)$  where  $\lambda_1 \geq \lambda_2 \geq \ldots \geq \lambda_n$ . (Some of the  $\lambda_i$  may be negative.)

What is the meaning of this? Initially the position of our particle is determined by the distribution  $\mathbf{p}^0$ . Since the vectors  $\mathbf{v}_1, \dots, \mathbf{v}_n$  form an orthonormal basis we can decompose  $\mathbf{p}^0$  in the form

$$\mathbf{p}^0 = \alpha_1 \mathbf{v}_1 + \dots + \alpha_n \mathbf{v}_n$$

where  $\alpha_i = \langle \mathbf{p}^0, \mathbf{v}_i \rangle$  and  $\alpha_1^2 + \cdots + \alpha_n^2 = 1$ .

After one step of the random walk, the distribution becomes

$$\mathbf{p}^1 = \mathbf{p}^0 A = \alpha_1 \mathbf{v}_1 A + \dots + \alpha_n \mathbf{v}_n A = \alpha_1 \lambda_1 \mathbf{v}_1 + \dots + \alpha_n \lambda_n \mathbf{v}_n$$

and after t steps

$$\mathbf{p}^t = \mathbf{p}^0 A^t = \alpha_1 \lambda_1^t \mathbf{v}_1 + \dots + \alpha_n \lambda_n^t \mathbf{v}_n. \tag{2}$$

Let's think of what happens when t becomes large. We will assume the values  $\alpha_i$  are nonzero since the initial position of the particle can be arbitrary. Eventually the right hand side of the expression will be dominated by the term in which  $\lambda_i$  has largest absolute value; this is either  $|\lambda_1|$  or  $|\lambda_n|$ . This absolute value cannot exceed 1, because  $\mathbf{p}^t$  would then become very large, but its norm is bounded since it is a probability distribution. Similarly, the absolute value cannot be less than 1 because then  $\mathbf{p}^t$  would become very small when t gets large. Finally, the largest  $\lambda_i$  in absolute value cannot be -1, because  $\mathbf{p}^t$  would then eventually be shifting signs; since it is a vector of probabilities, its entries must always be nonnegative.

Therefore, it must be the case that  $\lambda_1 = 1$ , and

$$\max\{|\lambda_i|: 2 \le i \le n\} = \max(\lambda_2, -\lambda_n) \le 1.$$

The quantity on the left side is denoted by  $\lambda = \lambda(G)$  and plays a very important role. Because  $\mathbf{u}A = \lambda_1 \mathbf{u}$ , so the eigenvector  $\mathbf{v}_1$  associated to  $\lambda_1 = 1$  equals  $\sqrt{n} \cdot \mathbf{u}$ . Now from (2) we have that

$$\|\mathbf{p}^t - \alpha_1 \mathbf{v}_1\|^2 = \alpha_2^2 \lambda_2^{2t} + \dots + \alpha_n^2 \lambda_n^{2t} \le \lambda^{2t}.$$

The left hand side has a natural interpretation. Recall that  $\alpha_1 = \langle \mathbf{p}^0, \mathbf{v}_1 \rangle = 1/\sqrt{n}$ , so  $\alpha_1 \mathbf{v}_1$  equals the uniform distribution  $\mathbf{u}$ . Thus  $\lambda^t$  measures how close  $\mathbf{p}^t$  gets to the uniform distribution after t steps of the walk:

$$\|\mathbf{p}^t - \mathbf{u}\| \le \lambda^t. \tag{3}$$

Another way of saying this is that  $\lambda$  determines the *rate* at which  $\mathbf{p}^t$  converges to the uniform distribution: The smaller  $\lambda$  is, the faster we will get to a uniformly random vertex.

## 2 Expander graphs

To get some intuition about equation (3), notice that in t steps the particle can reach at most  $1 + (d-1) + \cdots + (d-1)^t \le (d-1)^{t+1}$  vertices of the graph. This value is attained when the t-neighborhood of s is a d-regular tree. Let t be the largest value for which  $(d-1)^{t+1}$  is at most n/2. Then at least half the entries of  $\mathbf{p}^t$  are zero and

$$\lambda^t \ge \|\mathbf{p}^t - \mathbf{u}\| \ge (n/2 \cdot (0 - 1/n)^2)^{1/2} = \frac{1}{\sqrt{2n}} \ge \frac{1}{\sqrt{2(d-1)^{t+2}}}$$

from where  $\lambda \ge (1/\sqrt{d-1}) \cdot (2(d-1)^2)^{-1/2t}$ . As n gets larger, the second term approaches 1 and  $\lambda$  must be at least as large as  $1/\sqrt{d-1}$ .

A more precise analysis shows that for every graph,  $\lambda \geq 2\sqrt{d-1}/d - o_n(1)$ , where  $o_n(1)$  is quantity that converges to zero as n gets large. There exist graphs such that  $\lambda = 2\sqrt{d-1}/d$  for infinitely many values of n. Such graphs are called  $Ramanujan\ graphs$ .

<sup>&</sup>lt;sup>1</sup>This is not quite right: The correct way to say it is that for every index i there exists an initial position for the particle that makes  $\alpha_i \neq 0$ .

<sup>&</sup>lt;sup>2</sup>Ramanujan graphs are known to exist for every d such that d+1 is a power of a prime larger than two.

For our purposes, it will be enough to consider graph families for which as n grows,  $\lambda$  stays bounded away from one. If this is the case, then after only  $t = \Theta(\log n)$  steps of the random walk, we have that

$$\|\mathbf{p}^t - \mathbf{u}\| \le \lambda^{\Theta(\log n)} = n^{-\Theta(1)} \tag{4}$$

so  $\mathbf{p}^t$  gets very close to the uniform distribution, and in fact all vertices of G are reached with probability  $1/n \pm o(1/n)$  for the correct choice of  $\Theta(\cdot)$  constant.

**Definition 3.** A family of graphs  $\{G_n\}$ , where  $G_n$  has n vertices and is d-regular, is called an expander family if there is a constant  $\epsilon > 0$  such that  $\lambda(G_n) \leq 1 - \epsilon$  for every sufficiently large n.

## 3 Edge expansion

Suppose you start at a random vertex of some set S that is not too large and you take a random edge out of this vertex. How likely are you to get out of S? If a random walk out of any vertex s approaches the uniform distribution quickly, we would expect such a walk to avoid "getting stuck" in any set S. The following claim makes this intuition precise. The probability is taken over a pair of endpoints (u, w) of a random directed edge of G.

**Theorem 4.** For every set S of vertices,

$$\Pr_{(u,w)}[u \in S \text{ and } w \notin S] \ge (1 - \lambda_2) \Pr_u[u \in S] \Pr_w[w \notin S].$$

To prove this theorem it is useful to describe the eigenvalues of A, the normalized adjacency matrix of G, in an alternative way. We look at the value of the expression  $\mathbf{v}A\mathbf{v}^{\mathrm{T}}$  as  $\mathbf{v}$  ranges over all vectors of norm 1. We expand  $\mathbf{v}$  in the basis of eigenvectors

$$\mathbf{v} = \alpha_1 \mathbf{v}_1 + \dots + \alpha_n \mathbf{v}_n$$

where  $\alpha_1^2 + \cdots + \alpha_n^2 = 1$ . Then

$$\mathbf{v}A\mathbf{v}^{\mathrm{T}} = \left(\sum_{i=1}^{n} \alpha_{i} \mathbf{v}_{i}\right) A\left(\sum_{i=1}^{n} \alpha_{j} \mathbf{v}_{j}\right) = \sum_{i,j=1}^{n} \alpha_{i} \alpha_{j} \cdot \mathbf{v}_{i} A\mathbf{v}_{j}^{\mathrm{T}}.$$

Since  $\mathbf{v}_i A \mathbf{v}_j^{\mathrm{T}} = \lambda_i \mathbf{v}_i \mathbf{v}_j^{\mathrm{T}} = \lambda_i \langle \mathbf{v}_i, \mathbf{v}_j \rangle$  takes value  $\lambda_i$  when i = j and zero otherwise, we get

$$\mathbf{v}A\mathbf{v}^{\mathrm{T}} = \alpha_1^2 \lambda_1 + \alpha_2^2 \lambda_2 + \dots + \alpha_n^2 \lambda_n.$$

It follows that  $\mathbf{v}A\mathbf{v}^{\mathrm{T}}$  can be at most  $\lambda_1$  and this value is attained when  $\alpha_1 = 1$ , namely when  $\mathbf{v} = \mathbf{v}_1$ . So we can describe  $\lambda_1$  as

$$\lambda_1 = \max_{\|\mathbf{v}\|=1} \mathbf{v} A \mathbf{v}^T.$$

Similarly, we can describe  $\lambda_2$  as the maximum of  $\mathbf{v}A\mathbf{v}^T$  but taken only over those  $\mathbf{v}$  for which  $\alpha_1 = 0$ , namely those  $\mathbf{v}$  that are perpendicular to  $\mathbf{v}_1$ . In our case  $\mathbf{v}_1$  is parallel to  $\mathbf{u}$  so we can write

$$\lambda_2 = \max_{\|\mathbf{v}\| = 1, \mathbf{v} \perp \mathbf{u}} \mathbf{v} A \mathbf{v}^{\mathrm{T}}.$$

We will now give a probabilistic interpretation to the quantity  $1 - \lambda_2$ . Fix **v** such that  $\|\mathbf{v}\| = 1$  and notice that

$$\sum_{u,w=1}^{n} A_{uw}(\mathbf{v}(u) - \mathbf{v}(w))^{2} = \sum_{u,w=1}^{n} A_{uw}\mathbf{v}(u)^{2} + \sum_{u,w=1}^{n} A_{uw}\mathbf{v}(w)^{2} - 2\sum_{u,w=1}^{n} A_{uw}\mathbf{v}(u)\mathbf{v}(w).$$

Since each row and each column of A adds up to one, each of the first two sums equals the sum of squares of the entries of  $\mathbf{v}$ , which is 1. The third sum equals  $\mathbf{v}A\mathbf{v}^{\mathrm{T}}$ . Therefore we can write

$$1 - \lambda_2 = \frac{1}{2} \min_{\|\mathbf{v}\| = 1, \mathbf{v} \perp \mathbf{u}} \sum_{u,v=1}^n A_{uw} (\mathbf{v}(u) - \mathbf{v}(w))^2 = \frac{1}{2} \min_{\mathbf{v} \perp \mathbf{u}} \frac{\sum_{u,v=1}^n A_{uw} (\mathbf{v}(u) - \mathbf{v}(w))^2}{\sum_{u=1}^n \mathbf{v}(u)^2}.$$

Since there are dn directed edges of G, we get that  $\sum A_{uw}(\mathbf{v}(u) - \mathbf{v}(w))^2 = n \, \mathbf{E}_{(u,w)}[(\mathbf{v}(u) - \mathbf{v}(w))^2]$ . We also have  $\sum \mathbf{v}(u)^2 = n \, \mathbf{E}_u[\mathbf{v}(u)^2]$ . Therefore

$$1 - \lambda_2 = \frac{1}{2} \min_{\mathbf{v} \perp \mathbf{u}} \frac{E_{(u,w)}[(\mathbf{v}(u) - \mathbf{v}(w))^2]}{E_u[\mathbf{v}(u)^2]}.$$
 (5)

Proof of Theorem 4. Let S be any set of vertices,  $\alpha = \Pr[u \in S] = |S|/n$  and set

$$\mathbf{v}(u) = \begin{cases} 1 - \alpha, & \text{if } u \in S \\ -\alpha, & \text{if } u \notin S. \end{cases}$$

Notice that  $\mathbf{v} \perp \mathbf{u}$ , and that  $(\mathbf{v}(u) - \mathbf{v}(w))^2$  is 1 when exactly one of u and w is in S and the other is in  $\overline{S}$ , and 0 otherwise. In the first case we will say (u, w) crosses  $(S, \overline{S})$ . Plugging into (5) we obtain

$$1 - \lambda_2 \le \frac{1}{2} \frac{\Pr_{(u,w)}[(u,w) \text{ crosses } (S,\overline{S})]}{\mathbb{E}_u[\mathbf{v}(u)^2]}.$$

where

$$\Pr_{(u,w)}[(u,w) \text{ crosses } (S,\overline{S})] = 2\Pr[u \in S \text{ and } w \not\in S]$$

and

$$E_u[\mathbf{v}(u)]^2 = \alpha(1-\alpha)^2 + (1-\alpha)\alpha^2 = \alpha(1-\alpha) = \Pr[u \in S] \Pr[w \notin S].$$

## 4 Proof of Theorem 2

We now show how to deduce Theorem 2 from Theorem 1. Let  $\Phi$  be a 2CSP with no restrictions on the number of occurrences of each variable. We show how to get a new instance  $\Phi'$  out of  $\Phi$  where every variable occurs at most d times.

Each variable  $x_i$  in  $\Phi$  gives rise to  $n_i$  variables  $x'_{i1}, \ldots, x'_{in_i}$  in  $\Phi'$ . For each constraint  $\phi_{ii'}(x_i, x_{i'})$  in  $\Phi$  we assign unique copies  $x'_{ij}, x'_{i'j'}$  in  $\Phi'$  and add d/2 copies of the constraint  $\phi_{ii'}(x'_{ij}, x'_{i'j'})$  in  $\Phi'$ . Finally, for every i we fix a d/2-regular graph  $G_i$  on  $n_i$  vertices with edge expansion  $\lambda(G_i) \geq 1/2$  and introduce equality constraints  $x'_{ij} = x'_{ij'}$  for every edge (i, i') of  $G_i$ . We will call these the equality constraints for i. We will talk about how to construct such an expander in the next two lectures.

If  $\Phi$  has m/2 constraints, then  $\Phi'$  will have m variables and dm constraints. If  $\Phi$  is satisfiable, then  $\Phi'$  is clearly satisfiable. Now suppose we could find an assignment x' that satisfies a  $1 - \varepsilon$  fraction of the constraints of  $\Phi'$ . Then the following claim allows us to convert x' into an assignment that satisfies a  $1 - 18\varepsilon$  fraction of the constraints of  $\Phi$ :

Claim 5. If some assignment x' the violates at most an  $\varepsilon$ -fraction of contraints in  $\Phi'$ , then the assignment x where

$$x_i = plurality \ (most \ frequent) \ value \ among \ x'_{i1}, \ldots, x'_{in_i}$$

violates at most a  $34\varepsilon$  fraction of constraints in  $\Phi$ .

By Theorem 4, within every graph  $G_i$ 

$$|E(S, \overline{S})| \ge \frac{d|S||\overline{S}|}{4n_i}$$

for every subset S of vertices in  $G_i$ , where  $E(S, \overline{S})$  is the number of edges from a vertex in S to a vertex outside S.

Let  $S_i$  be the set of variables  $x'_{ij}$  that agree with the plurality value  $x_i$ . Let  $\varepsilon_i$  be the fraction of the  $dn_i/4$  equality constraints for i violated by the assignment x'. We will argue that  $|\overline{S_i}| \leq 8\varepsilon_i n_i$ :

- If  $|S_i| > n_i/2$ , then  $|E(S_i, \overline{S}_i)| \ge d|\overline{S}_i|/8$ . Since all the equality constraints for i between  $S_i$  and  $\overline{S}_i$  are violated by x',  $\varepsilon_i(dn_i/4) \ge |E(S_i, \overline{S}_i)|$ , so  $|\overline{S}_i| \le 2\varepsilon_i n_i$ .
- If  $n_i/4 \le |S_i| \le n_i/2$ , then  $|E(S_i, \overline{S}_i)| \ge d|S_i|/8 \ge dn_i/32$ . Since all the equality constraints for i between  $S_i$  and  $\overline{S}_i$  are violated by x', it follows that  $\varepsilon_i \ge 1/8$ , so  $|\overline{S}_i| \le n_i \le 8\varepsilon_i n_i$ .
- If  $|S_i| < n_i/4$ , then no value in  $\Sigma$  is taken by more than a 1/4-fraction of the  $x'_{ij}$ s, so there must exist some subset of values  $\Sigma' \subseteq \Sigma$  so the number of  $x'_{ij}$  taking values in  $\Sigma'$  is between  $n_i/4$  and  $n_i/2$ . Just like in the previous case, we get  $|\overline{S_i}| \le n_i \le 8\varepsilon_i n_i$ .

Now consider what happens in  $\Phi'$  when we replace the assignment x' with the plurality assignment  $x'_{\text{plur }ij} = x_i$  for every j. Replacing x' by  $x'_{\text{plur }}$  may cause the violation of at most  $(d/2)|\overline{S_i}|$  non-equality constraints for every i. If x' violates  $\varepsilon dm$  constraints,  $x'_{\text{plur }}$  will then violate at most

$$\varepsilon dm + \sum_{i=1}^{n} (d/2)|\overline{S_i}| \le \varepsilon dm + \sum_{i=1}^{n} (d/2)(8\varepsilon_i n_i) = \varepsilon dm + 16\sum_{i=1}^{n} \varepsilon_i dn_i/4 \le 17\varepsilon dm$$

constraints of  $\Phi'$ . This is a  $17\varepsilon$ -fraction of all the constraints in  $\Phi'$ . Since exactly half the constraints in  $\Phi'$  are equality constraints, x cannot violate more than a  $34\varepsilon$  fraction of constraints in  $\Phi$ .