# Near log-convexity of measured heat in (discrete) time and consequences 

Mert Sağlam<br>saglam@uw.edu


#### Abstract

Let $u, v \in \mathbb{R}_{+}^{\Omega}$ be positive unit vectors and $S \in \mathbb{R}_{+}^{\Omega \times \Omega}$ be a symmetric substochastic matrix. For an integer $t \geq 0$, let $m_{t}=\left\langle v, S^{t} u\right\rangle$, which we view as the heat measured by $v$ after an initial heat configuration $u$ is let to diffuse for $t$ time steps according to $S$. Since $S$ is entropy improving, one may intuit that $m_{t}$ should not change too rapidly over time. We give the following formalizations of this intuition.

We prove that $m_{t+2} \geq m_{t}^{1+2 / t}$, an inequality studied earlier by Blakley and Dixon (also Erdős and Simonovits) for $u=v$ and shown true under the restriction $m_{t} \geq e^{-4 t}$. Moreover we prove that for any $\epsilon>0$, a stronger inequality $m_{t+2} \geq t^{1-\epsilon} \cdot m_{t}^{1+2 / t}$ holds unless $m_{t+2} m_{t-2} \geq \delta m_{t}^{2}$ for some $\delta$ that depends on $\epsilon$ only. Phrased differently, $\forall \epsilon>0, \exists \delta>0$ such that $\forall S, u, v$ $$
\frac{m_{t+2}}{m_{t}^{1+2 / t}} \geq \min \left\{t^{1-\epsilon}, \delta \frac{m_{t}^{1-2 / t}}{m_{t-2}}\right\}, \quad \forall t \geq 2
$$ which can be viewed as a truncated log-convexity statement. Using this inequality, we answer two related open questions in complexity theory: Any property tester for $k$-linearity requires $\Omega(k \log k)$ queries and the randomized communication complexity of the $k$-Hamming distance problem is $\Omega(k \log k)$. Further we show that any randomized parity decision tree computing $k$-Hamming weight has size $\exp (\Omega(k \log k))$.


## 1 Introduction

Suppose that some initial heat configuration $u: \Omega \rightarrow \mathbb{R}_{+}$is given over a finite space $\Omega$ and the configuration evolves according to the map $w \mapsto S w$ in each time step $t=0,1, \ldots$, for some symmetric stochastic matrix $S: \Omega \times \Omega \rightarrow \mathbb{R}_{+}$. Assume that we are interested in the amount of heat contained in a certain region $R \subseteq \Omega$ and how this quantity changes over time. In notation, assuming $\|u\|_{2}=1$ for normalization purposes and $v(x):=\mathbb{1}_{x \in R} /|R|^{1 / 2}$ for $x \in \Omega$, we would like to understand how

$$
m_{t}:=\left\langle v, S^{t} u\right\rangle
$$

changes as a function of $t$. In this paper we derive local bounds that $\left\{m_{t}\right\}_{t=0}^{\infty}$ must obey for any $S, u$ and $v$ satisfying the symmetry, magnitude and positivity constraints above (in fact our bounds work for any countable $\Omega$, arbitrary non-negative unit vector $v$ and symmetric non-negative $S$ ). Further, we establish a tight connection between such bounds and the well-studied $k$-Hamming distance problem [PG86; Yao03; Cor+00; Bar+04; GKW04; Hua+06; BBM12; Buh+12; BBG14; Amb +15 ] and the $k$-Hamming weight problem [AFH12; BK12; Buh +12 ] and obtain the first tight bounds for respectively the communication complexity and parity decision tree complexity of them.

Our tight $\Omega(k \log (k / \delta))$ lower bound for the $\delta$-error communication complexity of the $k$-Hamming distance problem (that applies whenever $k^{2}<\delta n$ ) answers affirmatively a conjecture stated in [BBG14] (Conjecture 1.4). Prior to our work, the best impossibility results for this problem were an $\Omega\left(k \log ^{(r)} k\right)$ bits lower bound $\left(\log ^{(r)}\right.$ being the iterated logarithm) that applies to any randomized $r$-round communication protocol [ST13], and an $\Omega(k \log (1 / \delta))$ lower bound that applies to any $\delta$-error randomized protocol for $k<\delta n$ [BBG14].

Our parity decision tree lower bound shows that any $\delta$-error parity decision tree solving the $k$-Hamming weight problem has size $\exp \Omega(k \log (k / \delta))$, which directly implies an $\Omega(k \log (k / \delta))$ bound on the depth of any such decision tree. Previously no nontrivial lower bound was known for the parity decision tree size of this problem and an $\Omega(k \log (1 / \delta))$ bound on the parity decision tree depth followed from the communication complexity bound of [BBG14]. Prior to [BBG14], the best bound on the parity decision tree depth was $\Omega(k)$, derived in [BBM12] and [BK12].

Either by combining our communication complexity lower bound with the reduction technique developed in [BBM12] or by combining our parity decision tree lower bound with a reduction given in [BCK14], one obtains an $\Omega(k \log (k / \delta))$ bound for any (potentially adaptive) property tester for the $\delta$-error probability $k$-linearity testing problem. This establishes the correct bound for this problem which was studied extensively [Fis+02; Gol10; BK12; BCK14; Buh+12; BBM12] since [Fis+02] or earlier.

### 1.1 Motivating our bounds on $m_{t}$

We would like provide some intuition as to why one should expect

$$
\begin{align*}
& m_{t+2} \geq m_{t}^{1+2 / t}, \text { and }  \tag{1.1}\\
& m_{t+2} \geq m_{t}^{1+2 / t} \cdot \min \left\{t^{1-\epsilon}, \delta \frac{m_{t}^{1-2 / t}}{m_{t-2}}\right\} \tag{1.2}
\end{align*}
$$

to hold for appropriate $\epsilon, \delta$. Recalling that $S$ is a symmetric matrix with maximum eigenvalue 1 , we may write $S=Q D Q^{\top}$ for an orthonormal matrix $Q$ having columns $q_{x}, x \in \Omega$ and a diagonal matrix $D$ with entries $\lambda_{x} \leq 1, x \in \Omega$. Plugging this into $m_{t}=\left\langle v, S^{t} u\right\rangle$, we get

$$
\begin{equation*}
m_{t}=\sum_{x \in \Omega} \lambda_{x}^{t}\left\langle u, q_{x}\right\rangle\left\langle v, q_{x}\right\rangle \tag{1.3}
\end{equation*}
$$

For sake of analogy let us drop our assumption that $S, u, v$ are coordinate-wise nonnegative for a moment but instead assume that each summand in the right hand side of Eq. (1.3) is nonnegative by some coincidence. In this case we can consider $\left\{m_{t}\right\}_{t=0}^{\infty}$ as the moment sequence of a random variable supported on $[0,1]$ that takes the value $\left|\lambda_{x}\right|$ with probability $\left|\left\langle u, q_{x}\right\rangle\left\langle v, q_{x}\right\rangle\right|$ and the value 0 with probability $1-\sum_{x}\left|\left\langle u, q_{x}\right\rangle\left\langle v, q_{x}\right\rangle\right|$ (which is nonnegative by Cauchy-Schwarz inequality). This would imply that $\left\{m_{t}\right\}_{t=0}^{\infty}$ is completely monotone by Hausdorff's characterization [Hau21] and therefore log-convex (e.g., [NP05], Section 2.1, Example 6).

One particular implication of the log-convexity of $\left\{m_{t}\right\}_{t=0}^{\infty}$, that $\frac{1}{t} \log m_{t}+\frac{t-1}{t} \log m_{0} \geq \log m_{1}$, when combined with the fact $0 \leq m_{0} \leq 1$ (that follows from our assumption on the terms of Eq. (1.3)), leads to $m_{t} \geq m_{1}^{t}$. In 1958, Mandel and Hughes showed that if $u=v$, rather surprisingly, one can trade the assumption that the summands of Eq. (1.3) are nonnegative with the assumption that $S$ and $u=v$ are coordinate-wise nonnegative and still obtain the conclusion $m_{t} \geq m_{1}^{t}$ :

Theorem 1.1 (Mandel and Hughes [MH58]). Let $u$ be a nonnegative unit vector and $S$ be $a$ symmetric matrix with nonnegative entries. For an integer ${ }^{1} t \geq 1$ we have $\left\langle u, S^{t} u\right\rangle \geq\langle u, S u\rangle^{t}$.

A more general implication of the log-convexity of $\left\{m_{t}\right\}_{t=0}^{\infty}$ and that $m_{0} \leq 1$ is that for $k \geq t$, $\frac{t}{k} \log m_{k}+\frac{k-t}{k} \log m_{0} \geq \log m_{t}$, therefore $m_{k}^{t} \geq m_{t}^{k}$. In 1966, Blakley and Dixon [BD66] investigated whether $m_{k}^{t} \geq m_{t}^{k}$ holds in the case $u=v$ when the nonnegativity assumption on the summands of Eq. (1.3) is replaced by the coordinate-wise nonnegativity of $S, u=v$. They note that the inequality $m_{k}^{t} \geq m_{t}^{k}$ fails when $k$ and $t$ have different parity and otherwise holds true under the restriction $m_{t} \geq e^{-4 t}$. While the following is not explicitly stated as a conjecture in [BD66], they write
if $t>1,[\ldots]$ we cannot show that the inequality Eq. (1.1) holds for each nonnegative $|\Omega|$-vector $u$ if $S$ is nonnegative.

[^0]so with the earlier caveat we attribute the following to Blakley and Dixon [BD66]:
Conjecture 1.2 (Blakley and Dixon [BD66]). Let $S: \Omega \times \Omega \rightarrow \mathbb{R}_{+}$be a symmetric matrix with nonnegative entries and let $u: \Omega \rightarrow \mathbb{R}_{+}$be a nonnegative unit vector. For positive integers $k \geq t$ of the same parity, we have
$$
\left\langle u, S^{k} u\right\rangle^{t} \geq\left\langle u, S^{t} u\right\rangle^{k}
$$

In Section 3 we prove the following theorem which shows that a generalization of Conjecture 1.2 holds true.

Theorem 1.3. Let $S: \Omega \times \Omega \rightarrow \mathbb{R}_{+}$be a symmetric matrix with nonnegative entries and $u, v: \Omega \rightarrow$ $\mathbb{R}_{+}$be nonnegative unit vectors. For positive integers $k \geq t$ of the same parity, we have

$$
\left\langle v, S^{k} u\right\rangle^{t} \geq\left\langle v, S^{t} u\right\rangle^{k}
$$

It goes without saying that Eq. (1.1) is equivalent to Theorem 1.3 as we can rearrange Eq. (1.1) to $m_{t+2}^{1 \rho(t+2)} \geq m_{t}^{1 / t}$ and apply it iteratively to obtain $m_{k}^{1 / k} \geq \cdots \geq m_{t+2}^{1 /(t+2)} \geq m_{t}^{1 / t}$ whenever $k \geq t$ and $k, t$ have the same parity. Moreover, while defining Eq. (1.1) we assumed $S$ to be substochastic only to illustrate our interpretation of the inequality: indeed any nonnegative $S$ can be scaled to be substochastic as both sides of Eq. (1.1) are $(t+2)$-homogeneous in $S$.

In Theorem 1.1 and Theorem 1.3 we observed that increasingly more general implications of the log-convexity of $\left\{m_{t}\right\}_{t=0}^{\infty}$ can be derived by only assuming the coordinate-wise nonnegativity of $S, u$ and $v$. One may naturally wonder if the coordinate-wise nonnegativity of $S, u$ and $v$ implies the log-convexity of $\left\{m_{t}\right\}_{t=0}^{\infty}$ in its entirety. Unfortunately the following example shows that this is far from the truth.


Figure 1: $\Omega=\{0,1, \ldots, t\}, S(i, i+1)=S(i+1, i)=\epsilon$ for $i=0, \ldots, t-1$.
Consider the transition matrix $S$ on $\Omega=\{0,1, \ldots, t\}$ such that $S(i, i+1)=S(i+1, i)=\epsilon$ for $i=0,1, \ldots, t-1$ and $S(i, j)=0$ otherwise. Let $u$ and $v$ be the point masses respectively on states 0 and $t$; namely $u=[1,0, \ldots, 0]^{\top}$ and $v=[0,0, \ldots, 1]^{\top}$. We have $m_{t-2}=0, m_{t}=\epsilon^{t}$ and $m_{t+2}=t \epsilon^{t+2}$. Therefore $m_{t-2} m_{t+2}=0 \nsupseteq \epsilon^{2 t}=m_{t}^{2}$. In this example the log-convexity breaks (in the strongest possible way) because the states 0 and $t$ are separated by $t$ hops according to $S$ and the point mass at state 0 cannot reach state $t$ before the $t$ th time step.

Our next theorem shows that such reachability issues are essentially the only way the logconvexity property can fail to hold:

Theorem 1.4. For every $\epsilon>0$ there is a $\delta>0$ such that for any symmetric matrix $S: \Omega \times \Omega \rightarrow \mathbb{R}_{+}$ and unit vectors $u, v: \Omega \rightarrow \mathbb{R}_{+}$with nonnegative entires, defining $m_{t}$ as before, we have

$$
\begin{equation*}
\frac{m_{t+2}}{m_{t}^{1+2 / t}} \geq \min \left\{t^{1-\epsilon}, \delta \frac{m_{t}^{1-2 / t}}{m_{t-2}}\right\}, \quad \forall t \geq 2 \tag{1.4}
\end{equation*}
$$

In other words, Theorem 1.4 shows that one can recover a truncated version of the log-convexity of $\left\{m_{t}\right\}_{t=0}^{\infty}$ from just the coordinate-wise nonnegativity assumption of $S, u$ and $v$. We stress that Theorem 1.4 is tight up to the appearance of $\epsilon$ and the choice of $\delta=\delta(\epsilon)$. A direct calculation on Figure 1 for time steps $t, t+2, t+4$ shows that Eq. (1.4) cannot be improved to

$$
\frac{m_{t+2}}{m_{t}^{1+2 / t}} \geq \min \left\{t^{1-2 / t},\left(\frac{1+\eta}{2}\right) \frac{m_{t}^{1-2 / t}}{m_{t-2}}\right\}
$$

for $\eta>0$.

### 1.2 Related work on $m_{t}$

Almost simultaneously with the work of Mandel and Hughes [MH58], Mulholland and Smith also prove Theorem 1.1 in [MS59] and moreover they characterize the equality conditions of the inequality. Independently, in 1965, Blakley and Roy [BR65] prove the same inequality and characterize the equality conditions and [Lon66] provides an alternative proof to that of [MS59] in 1966. We remark that Theorem 1.1 is most commonly referred to as the Blakley-Roy bound or "Sidorenko's conjecture for paths". Note these results show that Conjecture 1.2 is true whenever $t$ divides $k$. Finally in 2012, Pate shows that $m_{t} \geq m_{1}^{t}$ without the restriction $u=v$ :
Theorem 1.5 (Pate [Pat12]). Let $S: \Omega \times \Omega \rightarrow \mathbb{R}_{+}$be a symmetric matrix with nonnegative entries and let $u, v: \Omega \rightarrow \mathbb{R}_{+}$be nonnegative unit vectors. It holds that

$$
\left\langle v, S^{2 t+1} u\right\rangle \geq\langle v, S u\rangle^{2 t+1}
$$

with equality if and only if $\left\langle v, S^{2 t+1} u\right\rangle=0$ or $S u=\lambda v$ and $S v=\lambda u$ for some $\lambda \in \mathbb{R}_{+}$.
This result already shows that Theorem 1.3 is true when $t$ divides $k$ but such a bound does not have any implications for our applications in complexity theory. In [ES82], Erdős and Simonovits conjecture the following.

Conjecture 1.6 (Erdős and Simonovits [ES82], Conjecture 6). For a graph $G=(V, E)$, let $w_{k}(G)$ be the number length $k$ walks in $G$ divided by $|V|$. For an undirected graph $G$, we have $w_{k}(G)^{t} \geq w_{t}(G)^{k}$ for $k>t$ of the same parity.

Note that Conjecture 1.6 is a specialization of Conjecture 1.2 to $S$ having $0-1$ entries and $u=\mathbf{1} / \sqrt{|V|}$ therefore our Theorem 1.3 verifies Conjecture 1.6 as well.

### 1.3 Our results in complexity theory

Here we list our result in complexity theory; see Section 5 for the definition of the models and the problems. The following theorem (which was already known [BBG14]) is a consequence of Theorem 1.3 and uses the standard corruption technique in communication complexity.

Theorem 1.7. Any two party $\delta$-error randomized protocol solving the $k$-Hamming distance problem over length-n strings communicates at least $\Omega(k \log (1 / \delta))$ bits for $k^{2} \leq \delta n$.

The next is our main result for the communication complexity of the $k$-Hamming distance problem and is a consequence of Theorem 1.4. This result cannot be obtained by the standard corruption technique and requires a suitable modification similar to [She12].
Theorem 1.8. Any two party $\delta$-error randomized protocol solving the $k$-Hamming distance problem over length-n strings communicates at least $\Omega(k \log (k / \delta))$ bits for $k^{2} \leq \delta n$.
Theorem 1.9. Any $\delta$-error parity decision tree deciding the $k$-Hamming weight predicate over length-n strings has size $\exp \Omega(k \log (k / \delta))$ for $k^{2}<\delta n$.
Corollary 1.10. Any $\delta$-error probability property tester for $k$-linearity requires $\Omega(k \log (k / \delta))$ queries.

Note the bound $m_{k}^{t} \geq m_{t}^{k}$ obtained in [BD66] under the condition $m_{t} \geq e^{-4 t}$ does not have any implications for the communication complexity of the $k$-Hamming distance problem as our reduction crucially uses the fact that $u$ and $v$ are arbitrary, however it does lead to an $\exp \Omega(k)$ lower bound for the parity decision tree size of the $k$-Hamming weight problem when combined with our reduction.
Remark. Note that in Theorem 1.3, when $u=v$ and either both $k, t$ are even or $S$ is positive semidefinite, the summands of Eq. (1.3) become nonnegative and the inequality holds trivially. For our application in communication complexity we crucially use the fact that $u$ and $v$ are arbitrary and for our application in parity decision trees, one can do away with $u=v$ but only at the expense of having to choose $k, t$ odd. In both results the $S$ we choose has eigenvalues 1 and -1 with equal multiplicities and therefore far away from being positive semidefinite. In either case, the implications of Theorem 1.3 in complexity theory follow from the interesting cases of this theorem.

## 2 Preliminaries

We denote by $[n]$ the set $\{1,2, \ldots, n\}$. We take $\exp$ and $\log$ functions to the base 2 . Let $\Omega$ be a countable set. For a function $\mu: \Omega \rightarrow \mathbb{R}_{+}$and a set $\Psi \subseteq \Omega$, we use the shorthand

$$
\mu(\Psi):=\sum_{x \in \Psi} \mu(x) .
$$

A function $\mu: \Omega \rightarrow \mathbb{R}_{+}$is said to be a distribution on $\Omega$ if $\mu(\Omega)=1$ and a subdistribution if $\mu(\Omega) \leq 1$. For a function $\mu$ on $\Omega$, we define

$$
\operatorname{supp}(\mu):=\{x \in \Omega \mid \mu(x)>0\} .
$$

For two distributions $\mu: \Omega_{1} \rightarrow \mathbb{R}_{+}$and $\nu: \Omega_{2} \rightarrow \mathbb{R}_{+}$, let us denote by $\mu \nu$ the distribution on $\Omega_{1} \times \Omega_{2}$ given by $(\mu \nu)\left(x_{1}, x_{2}\right)=\mu\left(x_{1}\right) \nu\left(x_{2}\right)$.

For a discrete random variable $X$, we denote by $\operatorname{dist}(X)$ the distribution function of $X$ and we define $\operatorname{supp}(X):=\operatorname{supp}(\operatorname{dist}(X))$. If $X$ is so that $\operatorname{dist}(X): \Omega \rightarrow \mathbb{R}_{+}$, then we say that $X$ has sample space $\Omega$. Two random variables $X$ and $Y$ are said to be independent if $\operatorname{dist}(X Y)=\operatorname{dist}(X) \operatorname{dist}(Y)$.

Lemma 2.1 (Jensen [Jen06], Formula (5)). Let $X$ be a real-valued random variable and $f$ be a convex function. We have $\mathbb{E}[f(X)] \geq f(\mathbb{E}[X])$. When $f$ is strictly convex, the inequality holds with equality if and only if $X$ is constant with probability 1.

### 2.1 Information theory

In this section we review the definitions and facts we use from information theory. Let $\mu$ and $\nu$ be two nonnegative functions on $\Omega$. The Kullback-Leibler divergence [Wal45; KL51] of $\mu$ from $\nu$, denoted $\mathbf{D}(\mu \| \nu)$, is defined by

$$
\begin{equation*}
\mathbf{D}(\mu \| \nu):=\sum_{x \in \Omega} \mu(x) \log \frac{\mu(x)}{\nu(x)} \tag{2.1}
\end{equation*}
$$

Here, if $\mu(x)=0$ for some $x$, then its contribution to the summation is taken as 0 , even when $\nu(x)=0$. The divergence is undefined if there is an $x \in \Omega$ such that $\mu(x)>0$ and $\nu(x)=0$. It can be shown that if the related series converges for the right hand side of Eq. (2.1), it converges absolutely, which justifies leaving the summation order unspecified. A fundamental property of $\mathbf{D}(\cdot \| \cdot)$ is that the divergence of a distribution from a subdistribution is always nonnegative.

Lemma 2.2 (Gibbs [Gib02], Theorem VIII). Let $\mu, \nu: \Omega \rightarrow \mathbb{R}$ be such that $\mu$ is a distribution and $\nu$ is a subdistribution. We have $\mathbf{D}(\mu \| \nu) \geq 0$ with equality if and only if $\mu=\nu$.

Proof. This follows from the next lemma by noting that $-\log \nu(\Omega) \geq 0$ for $\nu$ a subdistribution.
Lemma 2.3 (Kullback and Leibler [KL51], Lemma 3.2). Let $\mu, \nu: \Omega \rightarrow \mathbb{R}_{+}$be so that $\mu$ is a distribution on $\Omega$ and $\operatorname{supp}(\mu)=\Psi \subseteq \Omega$. We have

$$
\mathbf{D}(\mu \| \nu) \geq-\log \nu(\Psi)
$$

with equality if an only if $\mu(x)=\nu(x) / \nu(\Psi)$ for $x \in \Psi$ and $\mu(x)=0$ for $x \notin \Psi$.
Proof. By Eq. (2.1) we write

$$
\mathbf{D}(\mu \| \nu)=-\sum_{x \in \Psi} \mu(x) \log \frac{\nu(x)}{\mu(x)} \geq-\log \sum_{x \in \Psi} \nu(x)=-\log \nu(\Psi)
$$

where the inequality follows from Lemma 2.1 and concavity of $z \mapsto \log z$ on $\mathbb{R}_{+}$. If $\mu(x)=\nu(x) / \nu(\Psi)$ for $x \in \Psi$, we have $\mathbf{D}(\mu \| \nu)=-\log \nu(\Psi)$ by direct computation. Otherwise $\mathbf{D}(\mu \| \nu)>-\log \nu(\Psi)$ by strict concavity of $z \mapsto \log z$.

We extend the divergence notation $\mathbf{D}(\cdot \| \cdot)$ to apply to random variables as follows. Let $X, Y$ be discrete random variables on the same sample space $\Omega$. Define

$$
\begin{equation*}
\mathbf{D}(X \| Y):=\mathbf{D}(\operatorname{dist}(X) \| \operatorname{dist}(Y)) \tag{2.2}
\end{equation*}
$$

With this notation in hand, we are ready to define the conditional divergence. Let $X_{1} X_{2}$ and $Y_{1} Y_{2}$ be random variables defined on the sample space $\Omega_{1} \times \Omega_{2}$. The divergence of $X_{1} \mid X_{2}$ from $Y_{1} \mid Y_{2}$ is defined by

$$
\begin{equation*}
\mathbf{D}\left(X_{1}\left|X_{2} \| Y_{1}\right| Y_{2}\right):=\underset{x_{2} \sim X_{2}}{\mathbb{E}} \mathbf{D}\left(X_{1}\left|X_{2}=x_{2} \| Y_{1}\right| Y_{2}=x_{2}\right) \tag{2.3}
\end{equation*}
$$

Here, for each $x_{2} \in \operatorname{supp}\left(X_{2}\right), X_{1} \mid X_{2}=x_{2}$ and $Y_{1} \mid Y_{2}=x_{2}$ are random variables on the sample space $\Omega_{1}$ obtained from, respectively $X_{1} X_{2}$ and $Y_{1} Y_{2}$, by conditioning on the second coordinate equaling $x_{2}$.

Lemma 2.4 (e.g., [CT06]). Let $X_{1} X_{2}$ and $Y_{1} Y_{2}$ be random variables, both on the sample space $\Omega_{1} \times \Omega_{2}$. We have

$$
\mathbf{D}\left(X_{1} X_{2} \| Y_{1} Y_{2}\right)=\mathbf{D}\left(X_{1} \| Y_{1}\right)+\mathbf{D}\left(X_{2}\left|X_{1} \| Y_{2}\right| Y_{1}\right)
$$

The mutual information Let $X$ and $Y$ be jointly distributed random variables. The mutual information of $X$ and $Y$, denoted $\mathbf{I}(X: Y)$, is defined as

$$
\begin{equation*}
\mathbf{I}(X: Y):=\mathbf{D}(\operatorname{dist}(X, Y) \| \operatorname{dist}(X) \operatorname{dist}(Y)) \tag{2.4}
\end{equation*}
$$

The mutual information of a random variable with itself, i.e., the quantity $\mathbf{I}(X: X)$ is called the Shannon entropy of $X$.

## 3 Monotonicity of $t \mapsto m_{2 t}^{1 /(2 t)}$ and $t \mapsto m_{2 t+1}^{1 /(2 t+1)}$

Let $S: \Omega \times \Omega \rightarrow \mathbb{R}_{+}$be a symmetric matrix with nonnegative entries, $u, v: \Omega \rightarrow \mathbb{R}_{+}$nonnegative unit vectors and $m_{t}=\left\langle v, S^{t} u\right\rangle$. In this section we prove Theorem 1.3 which we restate here (with additional equality conditions) for the convenience of the reader. Recall this theorem confirms Conjecture 1.2 and Conjecture 1.6.
Theorem 1.3 (restated). Let $S: \Omega \times \Omega \rightarrow \mathbb{R}_{+}$be a symmetric matrix with nonnegative entries and $u, v: \Omega \rightarrow \mathbb{R}_{+}$be nonnegative unit vectors. For positive integers $k \geq t$ of the same parity, we have

$$
\begin{equation*}
\left\langle v, S^{k} u\right\rangle^{t} \geq\left\langle v, S^{t} u\right\rangle^{k}, \tag{3.1}
\end{equation*}
$$

with equality if and only if $\left\langle v, S^{k} u\right\rangle=0$ or $S u=\lambda v$ and $S v=\lambda u$ for some $\lambda \in \mathbb{R}_{+}$when $t$ is odd and $u=v$ is an eigenvector of $S^{2}$ when $t$ is even.

We prove Theorem 1.3 by an information theoretic argument. Define the distributions $\mu:=$ $u /\|u\|_{1}$ and $\nu:=v /\|v\|_{1}$. Since either side of Eq. (3.1) is $k t$-homogeneous in $S$, we may assume that $S$ is substochastic by scaling as needed. Having fixed this normalization, we view Eq. (3.1) as a statement about random walks on $\Omega$ that start from a state sampled according to $\mu$ or $\nu$ and evolve according to the transition matrix $S$.

### 3.1 Reference random walks

Let $\Omega_{0}=\Omega \cup\{r\}$ for some state $r \notin \Omega$ and $t$ be a positive integer. Recall that $\mu=u /\|u\|_{1}$ and $\nu=v /\|v\|_{1}$. We start by defining random walks $F^{t}, B^{t}$ on $\Omega$ 。 that evolve in discrete time steps $-1,0,1, \ldots, t, t+1$.

The random walk $F^{t}$ starts at $r$ and transitions to a state $x \in \Omega$ with probability $\mu(x)$ at time step -1 . In steps $0,1, \ldots, t-1$, the random walk proceeds according to the transition matrix $S$. At the time step $t$, each state $x \in \Omega$ transitions to $r$ with probability $\nu(x)$ and transitions to an
arbitrary state in $\Omega$ with probability $1-\nu(x)$ (say, all of them to the same arbitrary state). We view $F^{t}$ as a joint random variable $F^{t}=\left(F_{-1}^{t}, F_{0}^{t}, \ldots, F_{t+1}^{t}\right)$, where $F_{i}^{t}$ is the location of the walk in time step $i$.

The random walk $B^{t}$ proceeds backwards in time. At time step $t+1$ the walk $B^{t}$ starts at $r$ and transitions to a state $x \in \Omega$ with probability $\nu(x)$. In time steps $t, t-1, \ldots, 1$, the random walk proceeds as prescribed by $S$. At time step 0 , each state $x \in \Omega$ transitions to $r$ with probability $\mu(x)$ and to an arbitrary state in $\Omega$ with probability $1-\mu(x)$. Similarly, $B^{t}$ denotes the joint random variable $B^{t}=\left(B_{-1}^{t}, B_{0}^{t}, \ldots, B_{t+1}^{t}\right)$, where $B_{i}^{t}$ is the location of the walk at time step $i$.

The following facts about $F^{t}$ and $B^{t}$ are immediate. The random variables $F_{-1}^{t}$ and $B_{t+1}^{t}$ are fixed to a single value $r$. The random variables $F^{t}, B^{t}$ are Markovian, namely, $\operatorname{dist}\left(F_{i}^{t} \mid F_{i-1}^{t}, \ldots, F_{-1}^{t}\right)=$ $\operatorname{dist}\left(F_{i}^{t} \mid F_{i-1}^{t}\right)$ and $\operatorname{dist}\left(B_{i-1}^{t} \mid B_{i}^{t}, \ldots B_{t+1}^{t}\right)=\operatorname{dist}\left(B_{i-1}^{t} \mid B_{i}^{t}\right)$ for $i \in\{0, \ldots, t+1\}$.

### 3.2 Random walks returning to the origin

Assume that $\operatorname{Pr}\left[F_{t+1}^{t}=r\right]>0$. Let $X$ be the walk $F^{t}$ conditioned on $F_{t+1}^{t}=r$. Note that $X$ is a random variable on the sample space $\Omega_{\circ}^{t+3}$. The next two lemmas explicitly calculate the distribution of $X$.

For a matrix $M: \Omega \times \Omega \rightarrow \mathbb{R}_{+}$, functions $f, g: \Omega \rightarrow \mathbb{R}_{+}$, and $x, y \in \Omega$ we use the shorthands

$$
\begin{aligned}
& M(f, y):=\sum_{x \in \Omega} f(x) M(x, y)=\left(M^{\top} f\right)(y) \\
& M(x, g):=\sum_{y \in \Omega} M(x, y) g(y)=(M g)(x) \\
& M(f, g):=\sum_{x, y \in \Omega} f(x) M(x, y) g(y)=f^{\top} M g
\end{aligned}
$$

where the last expression in each line is understood as a matrix vector multiplication.
Lemma 3.1. Under our assumption $S^{t}(\mu, \nu)>0$,
(i) we have $\operatorname{Pr}\left[X_{i}=x\right]=\frac{S^{i}(\mu, x) S^{t-i}(x, \nu)}{S^{t}(\mu, \nu)}$, and
(ii) if $S^{t-i}(x, \nu)>0$, we have $\operatorname{Pr}\left[X_{i+1}=y \mid X_{\leq i}=x_{\leq i}\right]=\frac{S\left(x_{i}, y\right) S^{t-i-1}(y, \nu)}{S^{t-i}\left(x_{i}, \nu\right)}$.

Proof. From the definition of $F^{t}$ (cf. Section 3.1), we have

$$
\begin{align*}
\operatorname{Pr}\left[F_{i}^{t}=x\right] & =S^{i}(\mu, x)  \tag{3.2}\\
\operatorname{Pr}\left[F_{t+1}^{t}=r \mid F_{i}^{t}=x\right] & =S^{t-i}(x, \nu)  \tag{3.3}\\
\operatorname{Pr}\left[F_{t+1}^{t}\right. & =r] \tag{3.4}
\end{align*}=S^{t}(\mu, \nu),
$$

Using Bayes' rule with Eq. (3.2), (3.3) and (3.4) gives (i). Combining Eq. (3.3), (3.5) and the observation that $F^{t}$ is Markovian gives (ii).

With Lemma 3.1 we confirm that the random variable $X=\left(X_{-1}, X_{0}, \ldots, X_{t+1}\right)$ is Markovian; in particular a time inhomogeneous random walk on $\Omega_{0}$. Next we observe that the random variable $B^{t}$ conditioned on $B_{-1}^{t}=r$ is precisely $X$ also.

Lemma 3.2. Under our assumption $S^{t}(\mu, \nu)>0$,
(i) we have $\operatorname{dist}(X)=\operatorname{dist}\left(B^{t} \mid B_{-1}^{t}=r\right)$, and
(ii) if $S^{i}(\mu, x)>0$, we have $\operatorname{Pr}\left[X_{i-1}=y \mid X_{\geq i}=x_{\geq i}\right]=\frac{S\left(x_{i}, y\right) S^{i-1}(y, \mu)}{S^{i}\left(x_{i}, \mu\right)}$.

Proof. For any $x \in \Omega_{\circ}^{t+3}$ with $x_{t+1}=r$,

$$
\begin{array}{rlr}
\operatorname{Pr}[X=x] & =\frac{\mu\left(x_{0}\right) \prod_{i=1}^{t} S\left(x_{i-1}, x_{i}\right) \nu\left(x_{t}\right)}{S^{t}(\mu, \nu)} \\
& =\frac{\nu\left(x_{t}\right) \prod_{i=1}^{t} S\left(x_{i}, x_{i-1}\right) \mu\left(x_{0}\right)}{S^{t}(\mu, \nu)} & \\
& =\frac{\operatorname{Pr}\left[B^{t}=x\right]}{\operatorname{Pr}\left[B_{t+1}^{t}=r\right]}=\operatorname{Pr}\left[B^{t}=x \mid B_{t+1}^{t}=r\right] & \text { (as } S \text { is symmetric) }
\end{array}
$$

This proves (i). Given (i), the proof of (ii) is the same as Lemma 3.1(ii).
Lemma 3.3. We have $\mathbf{D}\left(X \| F^{t}\right)=\mathbf{D}\left(X \| B^{t}\right)=-\log S^{t}(\mu, \nu)$.
Proof. Recall that $\operatorname{Pr}\left[F_{t+1}^{t}=r\right]=S^{t}(\mu, \nu)$. Since $X$ is obtained from $F^{t}$ by conditioning on $F_{t+1}^{t}=r$, the equality criteria of Lemma 2.3 are fulfilled and thus $\mathbf{D}\left(X \| F^{t}\right)=-\log S^{t}(\mu, \nu)$. The derivation of $\mathbf{D}\left(X \| B^{t}\right)$ is identical as per Lemma 3.2(i).

### 3.3 Longer random walks

Let $J$ be an integer valued random variable taking the values $\{1,2, \ldots, t\}$, each with equal probability. For each fixing $j$ of $J$ we perform a random walk $Z \mid J=j$ on $\Omega_{0}$ that evolves in time steps $-1,0,1, \ldots, t, t+1, t+2, t+3$ as follows.

The random walk starts at $r$ and for each time step $-1 \leq i<j$, proceeds according to the transition kernel $\operatorname{dist}\left(X_{i+1} \mid X_{i}\right)$. At time step $j$, the random walk proceeds according to $\operatorname{dist}\left(X_{j-1} \mid X_{j}\right)$ and in time steps $j<i \leq t+3$ proceeds according to the transition kernel $\operatorname{dist}\left(X_{i-1} \mid X_{i-2}\right)$. We view $Z$ as a joint random variable $Z=\left(Z_{-1}, Z_{0}, \ldots, Z_{t+3}\right)$, where $Z_{i}$ denotes the location of the random walk at time step $i$.
Lemma 3.4. For $-1 \leq i \leq j$, we have $\operatorname{dist}\left(Z_{i} \mid J=j\right)=\operatorname{dist}\left(X_{i}\right)$ and for $j<i \leq t+3$, $\operatorname{dist}\left(Z_{i} \mid J=j\right)=\operatorname{dist}\left(X_{i-2}\right)$.
Proof. This follows from the fact that

$$
\operatorname{dist}\left(F^{t} \mid F_{t+1}^{t}=r\right)=\operatorname{dist}(X)=\operatorname{dist}\left(B^{t} \mid B_{-1}^{t}=r\right)
$$

and that $X$ is an actual random walk (i.e., Markovian) on $\Omega_{0}$.
To be more explicit, we have $\operatorname{dist}\left(X_{i}\right)=\operatorname{dist}\left(Z_{i}\right)$ for $i \leq j$ since both $X$ and $Z$ start at $r$ in time step -1 and evolve according to the transition kernel $\operatorname{dist}\left(X_{i+1} \mid X_{i}\right)$ for $i=-1, \ldots, j-1$. Since $\operatorname{dist}\left(X_{j}\right)=\operatorname{dist}\left(Z_{j}\right)$ and $Z$ proceeds according to $\operatorname{dist}\left(X_{i-1} \mid X_{j}\right)$ at time step $j$, by Lemma 3.2, $\operatorname{dist}\left(X_{j-1}\right)=\operatorname{dist}\left(Z_{j+1}\right)$. Finally in time steps $i>j$, we have $\operatorname{dist}\left(Z_{i}\right)=\operatorname{dist}\left(X_{i-2}\right)$ since $\operatorname{dist}\left(X_{j-1}\right)=\operatorname{dist}\left(Z_{j+1}\right)$ and $Z$ proceeds according to $\operatorname{dist}\left(X_{i-1} \mid X_{i-2}\right)$.

From this we can deduce that $Z$ always ends up in $r$ at time step $t+3$. We next argue that if $X$ does not diverge too much from the reference random walk $F^{t}$, then $Z$ does not diverge too much from $F^{t+2}$.
Lemma 3.5. We have

$$
\begin{aligned}
\mathbf{D}\left(Z \mid J \| F^{t+2}\right)=\frac{t+2}{t} \mathbf{D}\left(X \| F^{t}\right) & -\frac{1}{t}\left(\mathbf{D}\left(X_{0} \| F_{0}^{t}\right)+\mathbf{D}\left(X_{t+1}\left|X_{t} \| F_{t+1}^{t}\right| F_{t}^{t}\right)\right) \\
& -\frac{1}{t}\left(\mathbf{D}\left(X_{t} \| B_{t}^{t}\right)+\mathbf{D}\left(X_{-1}\left|X_{0} \| B_{-1}^{t}\right| B_{0}^{t}\right)\right)
\end{aligned}
$$

Proof. For a fixing $j$ of $J$, we have

$$
\begin{aligned}
\mathbf{D}\left(Z \mid J=j \| F^{t+2}\right)= & \sum_{i=-1}^{j-1} \mathbf{D}\left(X_{i+1}\left|X_{i} \| F_{i+1}^{t+2}\right| F_{i}^{t+2}\right)+\mathbf{D}\left(X_{j-1}\left|X_{j} \| F_{j+1}^{t+2}\right| F_{j}^{t+2}\right) \\
& +\sum_{i=j+1}^{t+2} \mathbf{D}\left(X_{i-1}\left|X_{i-2} \| F_{i+1}^{t+2}\right| F_{i}^{t+2}\right)
\end{aligned}
$$

where we have used the chain rule for divergence (cf. Lemma 2.4), the fact that $Z \mid J=j$ and $F^{t+2}$ are Markovian and Lemma 3.4. Recalling that $F^{t}$ and $B^{t}$ evolve according to $S$ in time steps $0,1, \ldots, t-1$, and $\operatorname{dist}\left(F_{t+1}^{t} \mid F_{t}^{t}\right)=\operatorname{dist}\left(F_{t+3}^{t+2} \mid F_{t+2}^{t+2}\right)$, we write

$$
\begin{aligned}
\mathbf{D}\left(Z \mid J=j \| F^{t+2}\right)= & \sum_{i=-1}^{t} \mathbf{D}\left(X_{i+1}\left|X_{i} \| F_{i+1}^{t}\right| F_{i}^{t}\right) \\
& \quad+\mathbf{D}\left(X_{j-1}\left|X_{j} \| B_{j-1}^{t}\right| B_{j}^{t}\right)+\mathbf{D}\left(X_{j}\left|X_{j-1} \| F_{j}^{t}\right| F_{j-1}^{t}\right) \\
= & \mathbf{D}\left(X \| F^{t}\right)+\mathbf{D}\left(X_{j-1}\left|X_{j} \| B_{j-1}^{t}\right| B_{j}^{t}\right)+\mathbf{D}\left(X_{j}\left|X_{j-1} \| F_{j}^{t}\right| F_{j-1}^{t}\right)
\end{aligned}
$$

again by the chain rule for divergence (Lemma 2.4) and the fact that $X$ and $F^{t}$ are Markovian. Now taking an expectation over all $j \in \operatorname{supp}(J)$, we have

$$
\begin{aligned}
\mathbf{D}\left(Z \mid J \| F^{t+2}\right)= & \frac{1}{t} \sum_{j} \mathbf{D}\left(Z \mid J=j \| F^{t+2}\right) \\
= & \mathbf{D}\left(X \| F^{t}\right)+\frac{1}{t} \sum\left(\mathbf{D}\left(X_{j-1}\left|X_{j} \| B_{j-1}^{t}\right| B_{j}^{t}\right)+\mathbf{D}\left(X_{j}\left|X_{j-1} \| F_{j}^{t}\right| F_{j-1}^{t}\right)\right) \\
= & \mathbf{D}\left(X \| F^{t}\right)+\frac{\mathbf{D}\left(X \| B^{t}\right)-\mathbf{D}\left(X_{t} \| B_{t}^{t}\right)-\mathbf{D}\left(X_{-1}\left|X_{0} \| B_{-1}^{t}\right| B_{0}^{t}\right)}{t} \\
& +\frac{\mathbf{D}\left(X \| F^{t}\right)-\mathbf{D}\left(X_{0} \| F_{0}^{t}\right)-\mathbf{D}\left(X_{t+1}\left|X_{t} \| F_{t+1}^{t}\right| F_{t}^{t}\right)}{t}
\end{aligned}
$$

Since $\mathbf{D}\left(X \| B^{t}\right)=\mathbf{D}\left(X \| F^{t}\right)$ by Lemma 3.3, collecting the $\mathbf{D}\left(X \| F^{t}\right)$ terms we finish the proof.
Finally, we lower bound the negative terms in the statement of Lemma 3.5.
Lemma 3.6. We have

$$
\begin{align*}
\mathbf{D}\left(X_{0} \| F_{0}^{t}\right)+\mathbf{D}\left(X_{-1}\left|X_{0} \| B_{-1}^{t}\right| B_{0}^{t}\right) & \geq \mathrm{H}_{2}(\mu):=-\log \|\mu\|_{2}^{2}, \text { and }  \tag{3.6}\\
\mathbf{D}\left(X_{t} \| B_{t}^{t}\right)+\mathbf{D}\left(X_{t+1}\left|X_{t} \| F_{t+1}^{t}\right| F_{t}^{t}\right) & \geq \mathrm{H}_{2}(\nu):=-\log \|\nu\|_{2}^{2}, \tag{3.7}
\end{align*}
$$

where $\mathrm{H}_{2}(\cdot)$ denotes the second order Rényi entropy.
Proof. We only prove the first inequality as the second one is symmetric. By Lemma 3.1, we have

$$
\begin{aligned}
\mathbf{D}\left(X_{0} \| F_{0}^{t}\right) & =\sum_{x \in \Omega} \frac{\mu(x) S^{t}(x, \nu)}{S^{t}(\mu, \nu)} \log \frac{S^{t}(x, \nu)}{S^{t}(\mu, \nu)} \quad \text { and } \\
\mathbf{D}\left(X_{-1}\left|X_{0} \| B_{-1}^{t}\right| B_{0}^{t}\right) & =\sum_{x \in \Omega} \frac{\mu(x) S^{t}(x, \nu)}{S^{t}(\mu, \nu)} \log \frac{1}{\mu(x)}
\end{aligned}
$$

Let $\Psi=\operatorname{supp}\left(X_{0}\right)$. By adding the two terms we get

$$
\begin{align*}
\mathbf{D}\left(X_{0} \| F_{0}^{t}\right)+\mathbf{D}\left(X_{-1}\left|X_{0} \| B_{-1}^{t}\right| B_{0}^{t}\right) & =-\sum_{x \in \Psi} \frac{\mu(x) S^{t}(x, \nu)}{S^{t}(\mu, \nu)} \log \frac{\mu(x) S^{t}(\mu, \nu)}{S^{t}(x, \nu)} \\
& \geq-\log \sum_{x \in \Psi} \frac{\mu(x)^{2} S^{t}(x, \nu) S^{t}(\mu, \nu)}{S^{t}(\mu, \nu) S^{t}(x, \nu)}  \tag{3.8}\\
& =-\log \sum_{x \in \Psi} \mu(x)^{2} \\
& \geq-\log \sum_{x \in \Omega} \mu(x)^{2} \tag{3.9}
\end{align*}
$$

where the first inequality is by concavity of $z \mapsto \log z$ and the second inequality is true as the summands are nonnegative.

### 3.4 Combining the inequalities

Proof of Theorem 1.3. Note that $Z_{-1}$ is fixed to $r$ by definition (cf. Section 3.3) and $Z_{t+3}$ is fixed to $r$ by Lemma 3.4. Therefore by Lemma 2.3 we have

$$
\begin{align*}
-\log S^{t+2}(\mu, \nu) & \leq \mathbf{D}\left(Z \| F^{t+2}\right)  \tag{3.10}\\
& =\mathbf{D}\left(Z \mid J \| F^{t+2}\right)-\mathbf{I}(J: Z)  \tag{3.11}\\
& \leq \mathbf{D}\left(Z \mid J \| F^{t+2}\right)  \tag{3.12}\\
& \leq \frac{t+2}{t} \mathbf{D}\left(X \| F^{t}\right)+\frac{\log \|\mu\|_{2}^{2}+\log \|\nu\|_{2}^{2}}{t}
\end{align*}
$$

Here Eq. (3.11) follows from the chain rule for the divergence (Lemma 2.4) and the definition of mutual information (cf. Eq. (2.4)), Eq. (3.12) follows from the nonnegativity of mutual information and the last line follows from Lemma 3.5 and Lemma 3.6. Plugging in $\mathbf{D}\left(X \| F^{t}\right)=-\log S^{t}(\mu, \nu)$, provided by Lemma 3.3, we obtain

$$
-\log S^{t+2}(\mu, \nu) \leq-\frac{t+2}{t} \log S^{t}(\mu, \nu)+\frac{\log \|\mu\|_{2}^{2}+\log \|\nu\|_{2}^{2}}{t}
$$

Arranging, we get

$$
\|\mu\|_{2}^{2}\|\nu\|_{2}^{2}\left\langle\nu, S^{t+2} \mu\right\rangle^{t} \geq\left\langle\nu, S^{t} \mu\right\rangle^{t+2}
$$

and substituting $\mu=u /\|u\|_{1}, \nu=v /\|v\|_{1}$, and recalling that $u, v$ are unit vectors, we obtain

$$
\begin{align*}
\left\langle v, S^{t+2} u\right\rangle^{t} & \geq\left\langle v, S^{t} u\right\rangle^{t+2}, \text { i.e., } \\
m_{t+2} & \geq m_{t}^{1+2 / t} \tag{3.13}
\end{align*}
$$

By applying this inequality iteratively, we get $\left\langle v, S^{k} u\right\rangle^{t} \geq\left\langle v, S^{t} u\right\rangle^{k}$ or written differently $m_{k}^{1 / k} \geq$ $m_{t}^{1 / t}$ as long as $k>t$ and $k, t$ have the same parity.

Next we characterize the equality conditions of Eq. (3.13). Let us verify the 'if' direction of the statement. Clearly if $\left\langle v, S^{k} u\right\rangle=0$ then we have $\left\langle v, S^{k} u\right\rangle=\left\langle v, S^{t} u\right\rangle$ by the first part of the theorem and the fact that $\left\langle v, S^{t} u\right\rangle \geq 0$. If $S^{2} u=\lambda^{2} u$, then $m_{2 t}=\lambda^{2 t}$ and if $S u=\lambda v$ and $S v=\lambda u$, then $m_{2 t+1}=\lambda^{2 t+1}$, therefore in both $t$ even and $t$ odd cases the inequality holds with equality.

Conversely, if $0 \neq\left\langle v, S^{k+2} u\right\rangle=\left\langle v, S^{k} u\right\rangle$, then the inequalities (3.8), (3.9), (3.10), and (3.11) must hold with equality. Combining the assumption that Eq. (3.9) and (3.8) hold with equality with the strict concavity of $z \mapsto \log z$ and Jensen's lemma, we get that $S^{t} \nu=\lambda_{1} \mu+\sigma_{1}$ for some $\lambda_{1} \leq 1$ and $\sigma_{1} \in \mathbb{R}_{+}^{\Omega}$ satisfying $\operatorname{supp}\left(\sigma_{1}\right) \cap \operatorname{supp}(\mu)=\emptyset$. This also means that $\operatorname{Pr}\left[X_{0}=x\right]=\mu(x)^{2} /\|\mu\|_{2}^{2}$ for $x \in \Omega$ and a similar and symmetrical argument shows that $\operatorname{Pr}\left[X_{t}=x\right]=\nu(x)^{2} /\|\nu\|_{2}^{2}$ for $x \in \Omega$. Assuming Eq. (3.11) holds with equality leads to $\mathbf{I}(Z: J)=0$, which in turn shows that $\operatorname{dist}\left(X_{i}\right)=\operatorname{dist}\left(X_{i+2}\right)$ for $i=0, \ldots, t-2$. Let $X^{k+2}:=\left(F^{t+2} \mid F_{t+3}^{t+2}=r\right)$. From our assumption $0 \neq\left\langle v, S^{k+2} u\right\rangle=\left\langle v, S^{k} u\right\rangle$ we conclude that $\operatorname{dist}(Z)=\operatorname{dist}\left(X^{t+2}\right)$ as $X^{t+2}$ is the minimally divergent distribution from $F^{t+2}$ among distributions on walks ending at state $r$. Since $\operatorname{dist}\left(X_{2}^{t+2}\right)=\operatorname{dist}\left(X_{0}^{t+2}\right)=\operatorname{dist}\left(X_{0}\right)$, and $S^{t} \nu=\lambda_{1} \mu+\sigma_{1}$ we get that $S^{2} \mu=\lambda_{2} \mu+\sigma_{2}$ for some $\lambda_{2} \leq 1$ and $\sigma_{2} \in \mathbb{R}_{+}^{\Omega}$ satisfying $\operatorname{supp}\left(\sigma_{2}\right) \cap \operatorname{supp}(\mu)=\emptyset$. Now we will show that it must be that $\sigma_{2}=0$. Suppose for sake of contradiction that $\sigma_{2}(z)>0$ for some $z \notin \operatorname{supp}(\mu)$. There exists $x, y \in \Omega$ so that $\mu(x) S(x, y) S(y, z)>0$. If $y \in \operatorname{supp}\left(X_{1}\right)$ then adding to a walk $w \in \operatorname{supp}(X)$ with $w_{1}=y$ the loop $(y, z)(z, y)$ we obtain a length $t+2$ walk which is not in the support of $X^{t+2}$ as $\operatorname{dist}\left(X_{2}^{t+2}\right)=\operatorname{dist}\left(X_{2}\right)=\operatorname{dist}\left(X_{0}\right)$, which contradicts the fact that $X^{t+2}$ is defined as $F^{t+2} \mid F_{t+3}^{t+2}=r$. If on the other hand $y \notin \operatorname{supp}\left(X_{1}\right)$, adding the loop $(x, y)(y, x)$ to a walk with $w_{0}=x$ leads to a walk which is not in the support of $X^{t+2}$, which is a contradiction. Having established $S^{2} \mu=\lambda \mu$, we complete the proof for even $t$ by recalling that $\operatorname{dist}\left(X_{0}\right)=\operatorname{dist}\left(X_{t}\right)$ therefore $\mu=\nu$. For $t$ odd, the last argument shows that $S \mu=\lambda_{3} \nu+\sigma_{3}$ for some $\lambda_{3} \leq 1$ and $\sigma_{3} \in \mathbb{R}_{+}^{\Omega}$ satisfying $\operatorname{supp}\left(\sigma_{3}\right) \cap \operatorname{supp}(\mu)=\emptyset$. It remains to show that $\sigma_{3}=0$ by using the assumption
that $\operatorname{dist}(Z)=\operatorname{dist}\left(X^{t+2}\right)$. Suppose $\sigma_{3}(y)>0$ for some $y \notin \operatorname{supp}(\nu)$. There exists $x \in \Omega$ such that $\mu(x) S(x, y)>0$. Adding the loop $(x, y)(y, x)$ to a walk $w$ with $w_{t-1}=x$ leads to a length $t+2$ walk which is not in the support of $X^{t+2}$. A symmetrical argument shows that $\lambda^{\prime} \mu=S \nu$. This completes the proof.

## $4 \quad$ Near log-convexity of $t \mapsto m_{2 t}$ and $t \mapsto m_{2 t+1}$

In this section we would like to prove the following improvement to Eq. (3.13): for all $\epsilon>0$ there exists a $\delta>0$ such that

$$
\begin{equation*}
m_{t+2} \geq m_{t}^{1+2 / t} \cdot \min \left\{t^{1-\epsilon},\left\lceil\delta \frac{m_{t}^{1-2 / t}}{m_{t-2}}\right\rceil\right\}, \quad \forall t \geq 2 \tag{4.1}
\end{equation*}
$$

Recall that in proving Eq. (3.13), in line (3.12), we used the relaxation $\mathbf{I}(J: Z) \geq 0$. Note that $J$ is uniformly distributed on $[t]$ therefore has $\log t$ bits of entropy and provided that it is possible to infer $J$ from $Z$ (i.e., it is possibly to locate the time reversal we have inserted in $Z$ ) the $\mathbf{I}(J: K)$ term appears to be large enough to recover the factor

$$
\min \left\{t^{1-\epsilon},\left\lceil\delta \frac{m_{t}^{1-2 / t}}{m_{t-2}}\right\rceil\right\}
$$

Note moreover that intuitively we are able to infer $J$ from $Z$ better when $\frac{m_{t}^{1-2 / t}}{m_{t-2}}$ is high, as in such cases on average for a time step $i \in[t]$ and a typical $x \sim X_{i}$, the distributions $\operatorname{dist}\left(X_{i-1} \mid X_{i}=x\right)$ and $\operatorname{dist}\left(X_{i+1} \mid X_{i}=x\right)$ should be far from each other, as otherwise we can argue that there should be many $t-2$ walks as follows. If $\operatorname{dist}\left(X_{i-1} \mid X_{i}=x\right)$ and $\operatorname{dist}\left(X_{i+1} \mid X_{i}=x\right)$ are close to each other, there should be many $p \in \Omega$ which has high probability in both these distributions. Sample such a $p$, and attach to it a walk sampled from $X_{-1} X_{1} \ldots X_{i-1} \mid X_{i-1}=p$ and another walk sampled from $X_{i+1} X_{i+2} \ldots X_{t+1} \mid X_{i+1}=p$, which leads to a length $t-2$ walk returning to the origin. However if $m_{t-2}$ is low, this should not happen and therefore $\operatorname{dist}\left(X_{i-1} \mid X_{i}=x\right)$ and $\operatorname{dist}\left(X_{i+1} \mid X_{i}=x\right)$ on average should be far apart, which means that we can notice when we take a step backwards in time and therefore infer $J$. In particular, Figure 1 gives such an example where $m_{t-2}=0$ and we can always recover $J$ with certainty from a sample from $Z$ : whenever we take a step to the left, it must be that we are at time step $J$.

Given this discussion, a direct approach to proving Eq. (4.1) appears to bound

$$
\begin{equation*}
\mathbf{I}(Z: J) \geq \log \min \left\{t^{1-\epsilon},\left\lceil\delta \frac{m_{t}^{1-2 / t}}{m_{t-2}}\right\rceil\right\} \tag{4.2}
\end{equation*}
$$

Unfortunately, this approach does not seem to work as we demonstrate with an example in the full version of this paper. The problem here appears to be that we fix a single distribution $Z$ to explore the two cases of Eq. (4.1). In our final approach, we pick different distributions depending on the case we would like to prove. Namely, if $\mathbf{I}(J: Z) \geq(1-\epsilon) \log t$, then carrying out the calculations in Eq. (3.10) through (3.13) with the assumption $\mathbf{I}(J: Z) \geq(1-\epsilon) \log t$, we prove the first case, namely $m_{t+2} \geq t^{1-\epsilon} m_{t}^{1+2 / t}$ using the distribution $Z$. If $\mathbf{I}(J: Z)<(1-\epsilon) \log t$ on the other hand, we demonstrate two new random variables $W, Y$ which are distributed respectively on length $t+2$ and length $t-2$ paths so that $\mathbf{D}\left(W \| F^{t+2}\right)+\mathbf{D}\left(Y \| F^{t-2}\right) \leq-2 \log S^{t}(\mu, \nu)-\log \delta$, which implies that $m_{t-2} m_{t+2} \geq \delta m_{t}^{2}$. While $W$ and $Y$ are constructed by modifying $X$ in suitable ways, which is how $Z$ was constructed also, we do so with the hindsight of having inspected what causes $\mathbf{I}(J: Z)$ to be smaller than $(1-\epsilon) \log t$. It is precisely this adaptivity which enables this approach to overcome the difficulties encountered by the one suggested in Eq. (4.2).

If $\mathbf{I}(J: Z) \geq(1-\epsilon) \log t$, by plugging this into Eq. (3.12) and carrying out the following calculations, we get $m_{t+2} \geq t^{1-\epsilon} \cdot m_{t}^{1+2 / t}$. Therefore it remains to show there exists a $\delta>0$ such that assuming $\mathbf{I}(J: Z)<(1-\epsilon) \log t$, we have $m_{t+2} m_{t-2} \geq \delta m_{t}^{2}$. To do so, we will demonstrate
distributions $W$ and $Y$ on walks that start from $r \in \Omega_{\circ}$ and return to $r$ after spending respectively $t+2$ and $t-2$ time steps in $\Omega$ such that $\mathbf{D}\left(W \| F^{t+2}\right)+\mathbf{D}\left(Y \| F^{t-2}\right) \leq-2 \log S^{t}(\mu, \nu)-\log \delta$. Notice that by Lemma 2.1 this indeed implies that $m_{t+2} m_{t-2} \geq \delta m_{t}^{2}$. The distributions $W$ and $Y$ will be mixture of $\Theta(t)$ random walks, in particular, they are not Markovian in general.

For brevity let us set $\mu_{i}^{x}:=\operatorname{dist}\left(X_{i-1} \mid X_{i}=x\right)$ and $\nu_{i}^{x}:=\operatorname{dist}\left(X_{i+1} \mid X_{i}=x\right)$. Let $U$ be the unary encoding of $J$ : a length $t$ bit vector of which only the $J$ th coordinate is set. First we would like to understand the contribution of each bit of $U$ to $\mathbf{I}(Z: J)=\mathbf{I}(Z: U)$. Using the chain rule, we write

$$
\begin{align*}
&(1-\epsilon) \log t>\mathbf{I}(U: Z) \\
&=\sum_{i=1}^{t} \mathbf{I}\left(U_{i}: Z \mid U_{<i}\right)  \tag{4.3}\\
&=\sum_{i=1}^{t} \frac{t-i+1}{t} \mathbf{I}\left(U_{i}: Z \mid U_{<i}=0\right)  \tag{4.4}\\
& \geq \sum_{i=1}^{t} \frac{t-i+1}{t} \mathbf{I}\left(U_{i}: Z_{i} Z_{i+1} \mid U_{<i}=0\right)  \tag{4.5}\\
&=\sum_{i=1}^{t} \frac{1}{t} \underset{x \sim X_{i}}{\mathbb{E}} \mathbf{D}\left(\mu_{i}^{x} \| \lambda_{i} \mu_{i}^{x}+\left(1-\lambda_{i}\right) \nu_{i}^{x}\right) \\
& \quad+\sum_{i=1}^{t} \frac{t-i}{t} \underset{x \sim X_{i}}{\mathbb{E}} \mathbf{D}\left(\nu_{i}^{x} \| \lambda_{i} \mu_{i}^{x}+\left(1-\lambda_{i}\right) \nu_{i}^{x}\right) \tag{4.6}
\end{align*}
$$

where we set $\lambda_{i}:=1 /(t-i+1)$, which is the probability that $U_{i}=1 \mid U_{<i}=0$. Here, Eq. (4.3) follows from the chain rule, Eq. (4.4) is true because if $U_{<i} \neq 0$ then $U_{i}=0$ (as $U$ has a single coordinate that is one) and consequently the mutual information is zero, and Eq. (4.5) is the data processing inequality. Next we lower bound Eq. (4.6) by its first term (which is valid since $\mu_{i}^{x}, \nu_{i}^{x}$ are distributions hence the second term of Eq. (4.6) is nonnegative), obtaining

$$
\begin{equation*}
(1-\epsilon) \log t>\underset{i \sim J}{\mathbb{E}} \underset{x \sim X_{i}}{\mathbb{E}} \mathbf{D}\left(\mu_{i}^{x} \| \lambda_{i} \mu_{i}^{x}+\left(1-\lambda_{i}\right) \nu_{i}^{x}\right) \tag{4.7}
\end{equation*}
$$

To simplify the presentation, here we only provide the proof of Theorem 1.4 for $\epsilon>7 / 8$ which demonstrates the ideas in their simplest form. This bound already implies all our results in complexity theory, with a constant factor loss of no more than 8 . The proof for any $\epsilon>0$ can be found in the full version of this paper.

### 4.1 The bound for $\epsilon>7 / 8$

If we condition on the event $i \in\{1, \ldots,\lceil t / 2\rceil\}$, this expectation increases by a factor of at most 2 ; namely

$$
\underset{i \sim[t / 2]}{\mathbb{E}} \underset{x \sim X_{i}}{\mathbb{E}} \mathbf{D}\left(\mu_{i}^{x} \| \lambda_{i} \mu_{i}^{x}+\left(1-\lambda_{i}\right) \nu_{i}^{x}\right)<2(1-\epsilon) \log t
$$

By Markov's inequality

$$
\operatorname{Pr}_{i \sim[t / 2], x \sim X_{i}}\left[\mathbf{D}\left(\mu_{i}^{x} \| \lambda_{i} \mu_{i}^{x}+\left(1-\lambda_{i}\right) \nu_{i}^{x}\right) \geq 8(1-\epsilon) \log t\right]<1 / 4
$$

so it follows that there is a set $T \subseteq[[t / 2\rceil]$ of size at least $\lfloor t / 4\rfloor$ such that if $i \in T$ we have

$$
\begin{equation*}
\operatorname{Pr}_{x \sim X_{i}}\left[\mathbf{D}\left(\mu_{i}^{x} \| \lambda_{i} \mu_{i}^{x}+\left(1-\lambda_{i}\right) \nu_{i}^{x}\right) \geq 8(1-\epsilon) \log t\right]<1 / 2 \tag{4.8}
\end{equation*}
$$

For each $i \in T$ let $X_{i}^{\prime}$ be the random variable obtained from $X_{i}$ by conditioning on those $x \in \operatorname{supp}\left(X_{i}\right)$ satisfying $\mathbf{D}\left(\mu_{i}^{x} \| \lambda_{i} \mu_{i}^{x}+\left(1-\lambda_{i}\right) \nu_{i}^{x}\right)<8(1-\epsilon) \log t$. Furthermore, for each $i \in T$ and $x \in \operatorname{supp}\left(X_{i}^{\prime}\right)$, we construct distributions $\pi_{i}^{x}: \Omega \rightarrow \mathbb{R}_{+}$to be specified later. Let $P_{i}$ be sampled by $x \sim X_{i}^{\prime}$ first and then picking $p \sim \pi_{i}^{x}$.

### 4.2 The distributions $W$ and $Y$

Let $K$ be an integer sampled uniformly at random from the set $T$ (constructed in the previous section). For each fixing $k$ of $K$, the random variables $W \mid K=k$ and $Y \mid K=k$ are random walks (i.e., they are Markovian) constructed as follows. We first pick $x, p \sim X_{k}^{\prime} P_{k}$. The walk $Y \mid K=k$ is generated by concatenating a sample from $X_{-1} X_{0} \ldots X_{k-1} \mid X_{k-1}=p$ and an independent sample from $X_{k+1} \ldots X_{t+1} \mid X_{k+1}=p$. The walk $W$ is generated by concatenating a sample from $X_{-1} X_{0} \ldots X_{k} \mid X_{k}=x$, the path $(x, p)$ and $\left(p, x^{\prime}\right)$ for an independent sample $x^{\prime} \sim\left(X_{k}^{\prime} \mid P_{k}=p\right)$ and an independent sample from $\left.X_{k} \ldots X_{t+1}\right\rfloor X_{k}=x^{\prime}$.

For $k \in T$ we define another random walk $\check{X}^{k}=\left(\check{X}_{-1}^{k}, \ldots, \check{X}_{t+1}^{k}\right)$, only to be used in the analysis of $W$ and $Y$. We sample $x \sim X_{k}^{\prime}$ and set $\check{X}_{k}^{k}=x$. We pick the rest of the coordinates of $\check{X}^{k}$ according to the distribution $X \mid X_{k}=x$. Note that for any $k \in T$, we have

$$
\mathbf{D}\left(\check{X}^{k} \| X\right)=\mathbf{D}\left(X_{k}^{\prime} \| X_{k}\right) \leq 1
$$

by Eq. (4.8) and Lemma 2.3 and the fact that both $X$ and $\check{X}^{k}$ are Markovian.
Lemma 4.1. We have

$$
\begin{aligned}
\mathbf{D}\left(W \mid K=k \| F^{t+2}\right)+ & \mathbf{D}\left(Y \mid K=k \| F^{t-2}\right) \\
& \leq-2 \log S^{t}(\mu, \nu)+2+\underset{x \sim X_{k}^{\prime}}{\mathbb{E}} \mathbf{D}\left(\pi_{k}^{x} \| \mu_{k}^{x}\right)+\underset{x \sim X_{k}^{\prime}}{\mathbb{E}} \mathbf{D}\left(\pi_{k}^{x} \| \nu_{k}^{x}\right)
\end{aligned}
$$

Proof. We have

$$
\mathbf{D}\left(W \mid K=k \| F^{t+2}\right)=\mathbf{D}\left(\check{X}^{k} \| F^{t}\right)+\mathbf{D}\left(P_{k}\left|X_{k}^{\prime} \| F_{k+1}\right| F_{k}\right)+\mathbf{D}\left(X_{k}^{\prime}\left|P_{k} \| F_{k+1}\right| F_{k}\right)
$$

and further

$$
\begin{aligned}
\mathbf{D}\left(Y \mid K=k \| F^{t-2}\right)+ & \mathbf{D}\left(P_{k}\left|X_{k}^{\prime} \| F_{k+1}\right| F_{k}\right)+\mathbf{D}\left(X_{k}^{\prime}\left|P_{k} \| F_{k+1}\right| F_{k}\right) \\
& =\mathbf{D}\left(\check{X}^{k} \| F^{t}\right)+\underset{x \sim X_{k}^{\prime}}{\mathbb{E}} \mathbf{D}\left(\pi_{k}^{x} \| \mu_{k}^{x}\right)+\underset{x \sim X_{k}^{\prime}}{\mathbb{E}} \mathbf{D}\left(\pi_{k}^{x} \| \nu_{k}^{x}\right)
\end{aligned}
$$

Summing up the two inequalities and substituting $\mathbf{D}\left(\check{X}^{k} \| X\right) \leq 1$ we get the result.
At this point, in light of Lemma 4.1, we could pick each $\pi_{k}^{x}$ so that it minimizes $\mathbf{D}\left(\pi_{k}^{x} \| \mu_{k}^{x}\right)+$ $\mathbf{D}\left(\pi_{k}^{x} \| \nu_{k}^{x}\right)$ : the unique minimizer is given by $\pi_{k}^{x}=\sqrt{\mu_{k}^{x} \nu_{k}^{x}} /\left\langle\sqrt{\mu_{k}^{x}}, \sqrt{\nu_{k}^{x}}\right\rangle$. However doing so leads to $W, Y$ which diverge from the $F$ walk by more than a constant, and therefore is not good enough for our needs. To obtain better random variables $W$ and $Y$, we crucially use the fact that $W$ is a mixture of $\Theta(t)$ random walks. Namely, if we consider the entropy coming from the $\mathbf{I}(W: K)$ term also, a better strategy for picking the distributions $\pi_{k}^{x}$ becomes available. By contrast, we do not use the fact that $Y$ is a mixture and, in fact, it can be replaced by $Y \mid K=k_{0}$ where $k_{0}=\arg \min _{k} \mathbf{D}\left(Y \mid K=k \| F^{t-2}\right)$, however the averaged quantity $\mathbf{D}\left(Y \mid K \| F^{t-2}\right)$ is far more convenient to work with.

### 4.3 The contribution of $I(K: W)$

Similar to Eq. (4.7), we would like to understand the contribution of each time step $t \in T$ to $\mathbf{I}(K: W)$. Let $V$ be the unary encoding of $K$ : a length $t$ bit vector of which only the $V$ th coordinate is set. Using the chain rule for mutual information

$$
\begin{aligned}
\mathbf{I}(W: V) & =\sum_{i \in T} \mathbf{I}\left(V_{i}: W \mid V_{<i}\right) \\
& \geq \underset{k \sim K}{\mathbb{E}} \underset{x \sim X_{k}^{\prime}}{\mathbb{E}} \mathbf{D}\left(\pi_{k}^{x} \| \eta_{i} \pi_{k}^{x}+\left(1-\eta_{i}\right) \widetilde{v}_{k}^{x}\right),
\end{aligned}
$$

where $\eta_{k}=1 / \operatorname{rank}_{T}(k)$ and $\widetilde{v}_{k}^{x}:=\mathbb{E}_{j>k: j \in T} \operatorname{dist}\left(\check{X}_{k+1}^{j} \mid \check{X}_{k}^{j}=x\right)$. Here $\operatorname{rank}_{T}(i)$ denotes the position of $i \in T$ when the elements of $T$ are sorted in decreasing order. By Eq. (4.8), and the definition of $X_{k}^{\prime}$, we have $\widetilde{v}_{k}^{x}(y) \leq 2 \nu_{k}^{x}(y)$ for all $y \in \Omega$. Therefore we conclude that

$$
\begin{equation*}
\mathbf{I}(W: K) \geq \underset{k \sim K}{\mathbb{E}} \underset{x \sim X_{k}^{\prime}}{\mathbb{E}} \mathbf{D}\left(\pi_{k}^{x} \| \eta_{k} \pi_{k}^{x}+2\left(1-\eta_{k}\right) \nu_{k}^{x}\right) \tag{4.9}
\end{equation*}
$$

Note in the above divergence expression the reference measure is not a probability distribution, which our definition permits (cf. Eq. (2.1)).

Recall our goal in this section is to upper bound $\mathbf{D}\left(W \| F^{t+2}\right)+\mathbf{D}\left(Y \| F^{t-2}\right)+2 \log S^{t}(\mu, \nu)$ by $\log 1 / \delta$. Let us write

$$
\begin{align*}
\mathbf{D}\left(W \| F^{t+2}\right)+ & \mathbf{D}\left(Y \| F^{t-2}\right)+2 \log S^{t}(\mu, \nu) \\
& \leq \mathbf{D}\left(W \mid K \| F^{t+2}\right)+\mathbf{D}\left(Y \mid K \| F^{t-2}\right)-\mathbf{I}(K: W)+2 \log S^{t}(\mu, \nu) \\
& \leq 2+\underset{k \sim K, x \sim X_{k}^{\prime}}{\mathbb{E}} \mathbf{D}\left(\pi_{x}^{k} \| \mu_{x}^{k}\right)+\mathbf{D}\left(\pi_{x}^{k} \| \nu_{x}^{k}\right)-\mathbf{I}(K: W) \\
& \leq 2+\underset{k \sim K, x \sim X_{k}^{\prime}}{\mathbb{E}} \underset{y \sim \pi_{k}^{x}}{\mathbb{E}} \log \frac{\eta_{k} \pi_{k}^{x}(y)^{2}+2\left(1-\eta_{k}\right) \nu_{k}^{x}(y) \pi_{k}^{x}(y)}{\mu_{k}^{x}(y) \nu_{k}^{x}(y)} \tag{4.10}
\end{align*}
$$

where the second inequality follows from Lemma 4.1 and the last inequality is obtained by plugging in Eq. (4.9). Note that the function $z \mapsto z \log \left(a z^{2}+b z\right)$ is strictly convex in $\mathbb{R}_{+}$whenever $a b>0$, therefore for each $k, x$ there is a unique minimizer $\left(\pi_{k}^{x}\right)^{*}$ of Eq. (4.10), which can be calculated, say, using Lagrange multipliers. However, instead of the minimizer, we work with a simple approximation of it. For each $k \in T$ and $x \in \operatorname{supp}\left(X_{k}^{\prime}\right)$, we let

$$
\Psi_{k}^{x}:=\left\{y \in \Omega \left\lvert\, \nu_{k}^{x}(y) \geq \frac{\lambda_{k}}{1-\lambda_{k}} \mu_{k}^{x}(y)\right.\right\}
$$

By definition of $X_{k}^{\prime}$, we have $\mathbf{D}\left(\mu_{i}^{x} \| \lambda_{i} \mu_{i}^{x}+\left(1-\lambda_{i}\right) \nu_{i}^{x}\right)<8(1-\epsilon) \log t$. Let $\gamma=1-8(1-\epsilon)$, which is positive by our assumption $\epsilon>7 / 8$. By Markov's inequality, and the fact that $\lambda_{k} \leq 2 / t$, we get

$$
\mu_{k}^{x}\left(\Psi_{k}^{x}\right) \geq \gamma
$$

for large enough $t$. Let $\pi_{k}^{x}$ be $\mu_{k}^{x} \mid \Psi_{k}^{x}$, namely we have $\pi_{k}^{x}(y)=\mu_{k}^{x}(y) / \mu_{k}^{x}\left(\Psi_{k}^{x}\right)$ if $y \in \Psi_{k}^{x}$, and $\pi_{k}^{x}(y)=0$ otherwise. Continuing from Eq. (4.10), we have

$$
\begin{align*}
& \leq 2+\underset{k \sim K, x \sim X_{k}^{\prime}}{\mathbb{E}} \underset{y \sim \pi_{k}^{x}}{\mathbb{E}} \log \frac{\eta_{k} \pi_{k}^{x}(y)^{2}+2\left(1-\eta_{k}\right) \nu_{k}^{x}(y) \pi_{k}^{x}(y)}{\mu_{k}^{x}(y) \nu_{k}^{x}(y)} \\
& \leq 2+\underset{k \sim K}{\mathbb{E}} \log \left(\frac{\eta_{k}\left(1-\lambda_{k}\right)}{\lambda_{k} \gamma^{2}}+\frac{2}{\gamma}\right) \tag{4.11}
\end{align*}
$$

where the second inequality is true by definition of $\Psi_{k}^{x}$ and $\pi_{k}^{x}$. Now we argue that the expectation term in Eq. (4.11) is maximized when $T$ is the set containing the smallest $|T|$ elements of [ $[t / 2\rceil]$. To see this suppose there is an $i \notin T$ which is smaller than the maximum element of $T$. Let $j$ be the smallest item in $T$ which is greater than $i$. We see that the expectation term increases if we replace $T$ by $T \backslash\{j\} \cup\{i\}$ as $\log \left(\frac{C\left(1-\lambda_{k}\right)}{\lambda_{k} \gamma^{2}}+\frac{2}{\gamma}\right)$ is decreasing in $k$ and the ranks do not change
after swapping $j$ with $i$. Therefore,

$$
\begin{align*}
\mathbf{D}\left(W \| F^{t+2}\right)+ & \mathbf{D}\left(Y \| F^{t-2}\right)+2 \log S^{t}(\mu, \nu) \\
& \leq 2+\log \left(\prod_{i=1}^{|T|} \frac{t / 2+3 i}{i \gamma^{2}}\right)^{1 /|T|} \\
& =\log \frac{12}{\gamma^{2}}+\log \left(\prod_{i=1}^{|T|} \frac{t / 6+i}{i}\right)^{1 /|T|} \\
& \leq \log \frac{12}{\gamma^{2}}+\log \binom{2|T|}{|T|}^{1 /|T|}  \tag{4.12}\\
& \leq \log \frac{48}{\gamma^{2}} \tag{4.13}
\end{align*}
$$

where the $\binom{2|T|}{|T|}$ is the middle binomial coefficient, in the second inequality we use the fact $|T|>t / 6$, and the last inequality is true as $\binom{2 n}{n}<2^{2 n}$. Therefore it is enough to choose $\epsilon>7 / 8$ and $\delta \leq \frac{(1-8(1-\epsilon))^{2}}{48}=\frac{4}{3}(\epsilon-7 / 8)^{2}$. We have established the following.

Theorem 1.4 (restated). For any $\epsilon>7 / 8$ there is a $\delta>0$ such that $m_{t+2} \geq t^{1-\epsilon} m_{t}^{1+2 / t}$ unless $m_{t+2} m_{t-2} \geq \delta m_{t}^{2}$.

## 5 Randomized computational models

In this section we show the connection between Theorem 1.3, Theorem 1.4 and the randomized communication and query complexities of the Hamming distance problem.

### 5.1 Communication complexity

In a two player communication problem the players, named Alice and Bob, receive separate inputs, respectively $x \in \mathcal{X}$ and $y \in \mathcal{Y}$, and they communicate in order to compute the value $f(x, y)$ of a function $f: \mathcal{X} \times \mathcal{Y} \rightarrow\{0,1\}$ (known to both players). In an $r$-round protocol, the players can take at most $r$ turns alternately sending each other a message (that is, a bit string) and the last player to receive a message declares the output of the protocol. A protocol can be deterministic or randomized; in the latter case the players can base their actions on a common random source and we measure the error probability: the maximum over inputs $(x, y) \in \mathcal{X} \times \mathcal{Y}$, of the probability that the output of the protocol differs from $f(x, y)$. The communication cost of a protocol is the maximum, over the inputs and the random string, of the total number of bits sent between the players. For a function $f: \mathcal{X} \times \mathcal{Y} \rightarrow\{0,1\}$, an integer $r$ and $\delta \in[0,1]$, we denote by $R_{\delta}^{r}(f)$ the minimum over all protocols for $f$ having $r$-rounds and error probability at most $\delta$, of the communication cost incurred. We define $R_{\delta}(f)$ similarly, but we take the maximum over $\delta$ error protocols with no restriction on the number of rounds it uses.

In the $k$-Hamming distance problem, denoted $\operatorname{Ham}_{k}^{n}$, the players receive length- $n$ bit strings, respectively $x, y \in\{0,1\}^{n}$, and are required determine if $\|x-y\|_{1} \leq k$ or not. There is a well known one-round communication protocol which accomplishes this with error probability $\delta$ by communicating $O(k \log (k / \delta))$ bits.

Theorem 5.1 (e.g., Huang, Shi, Zhang and Zhu [Hua+06]). It holds that

$$
R_{\delta}^{1}\left(\operatorname{Ham}_{k}^{n}\right)=O(\min \{k \log (k / \delta), k \log (n / k)\})
$$

Highly related to the $\operatorname{Ham}_{k}^{n}$ is the $k$-disjointness problem $\operatorname{Disj}_{k}^{n}$, wherein the players each receive a $k$-subset of $[n]$ and their goal is to determine if their sets intersect. Notice that $\mathrm{Disj}_{k}^{n}$ can be seen as a promise version of $\operatorname{Ham}_{2 k-2}^{n}$ where each player is guaranteed to have a string with Hamming

| Problem | Upper bound | Rounds | Error | Lower bound | Reference |
| :---: | :---: | :---: | :---: | :---: | :---: |
| $\operatorname{Ham}_{k}^{n}$ | $O(k \log (k / \delta))$ | 1 | $\delta$ | Folklore, [Hua+06] |  |
|  | $\begin{gathered} \dot{2} \\ \stackrel{\rightharpoonup}{2} \\ \text { 完 } \end{gathered}$ | any | $\delta$ | $\Omega(k \log (1 / \delta))$ | [BBG14] |
|  |  | any | $\delta$ | $\Omega(k \log (k / \delta))$ | This work |
| $\operatorname{Disj}_{k}^{n}$ | $O(k \log (k / \delta))$ | 1 | $\delta$ |  | Folklore |
|  | $O(k)$ | $O(\log k)$ | $1 / 3$ |  | [HW07] |
|  | $O\left(k \log ^{(r)} k\right)$ | $r$ | $1 / \exp ^{(r)}\left(c \log ^{(r)} k\right)$ |  | [ST13] |
|  | $O(k)$ | $\log ^{*} k$ | $1 / \exp k$ |  | [ST13] |
|  |  | r | $1 / 3$ | $\Omega\left(k \log ^{(r)} k\right)$ | [ST13] |
|  |  | 1 | $1 / 3$ | $\Omega(k \log k)$ | [Buh+12; DKS12] |
|  |  | 1 | $\delta$ | $\Omega(k \log (1 / \delta))$ | [Sağ11; JW11] |
|  |  | any | $1 / 3$ | $\Omega(k)$ | [KS92] |

Table 1: Known bounds for $\operatorname{Disj}_{k}^{n}$ and $\operatorname{Ham}_{k}^{n}$.
weight $k$ : the sets are disjoint if and only if the Hamming distance between the characteristic vectors of the sets is more than $2 k-2$. Therefore any upper bound for the Ham $n k$ carries over to Disj ${ }_{k}^{n}$ and any lower bound for $\mathrm{Disj}_{k}^{n}$ carries over to $\operatorname{Ham}_{k}^{n}$. Around 1993, Håstad and Wigderson [HW07] showed that there is a more efficient protocol for $\mathrm{Disj}_{k}^{n}$ than that implied by Theorem 5.1, which communicates only $O(k)$ bits, but over $O(\log k)$ rounds.

On the lower bounds side, the result of [KS92] implies that $\Omega(k)$ bits is needed for these problems even if one uses arbitrarily large number of round protocols. In [Buh+12] it was shown that any 1-round protocol for $\operatorname{Disj}_{k}^{n}$ needs to communicate at least $\Omega(k \log k)$ bits when $k^{2}<n$ (this result was proven later in [DKS12] also). In Theorem 3.2 of [Sağ11], an $\Omega(k \log (1 / \delta))$ bound for 1-round complexity of $\mathrm{Disj}_{k}^{n}$ was shown even when Bob receives just one element (i.e., the indexing problem) for $k<\delta n$ and a slightly more general result was shown in [JW11]. Finally in [ST13] the communication complexity of $\operatorname{Disj}_{k}^{n}$ was settled:

$$
R_{1 / 3}^{r}\left(\operatorname{Disj}_{k}^{n}\right)=\Theta\left(k \log ^{(r)} k\right)
$$

for $1 \leq r<\log ^{*} k$ and $k<n^{2}$. Their upper bound solves the disjointness problem with error probability at most $1 / \exp k+1 / \exp ^{(r)}\left(c \log ^{(r)} k\right)$ for any $c>1$ by communicating $O\left(k \log ^{(r)} k\right)$ bits over $r$ rounds. In fact bulk of the bits is sent in the first round and the rest of the rounds amount to an $O(k)$ bits of communication. Taking $r=\log ^{*} k$, this lead to an $O(k)$ bits protocol with error probability that is exponentially small in $k$. Their lower bound shows that at least one message of size $\Omega\left(k \log ^{(r)} k\right)$ bits needs to be sent by any $r$-round protocol, even if it has $1 / 3$ error probability. Prior to this work, this lower bound provided the strongest lower bound for $\operatorname{Ham}_{k}^{n}$ also, along with the incomparable bound of $\Omega(k \log (1 / \delta))$ due to [BBG14] which holds for any number of rounds, which we discuss shortly.

To summarize the above results, the 1-round communication complexity of both $\mathrm{Disj}_{k}^{n}$ and $\operatorname{Ham}_{k}^{n}$ is $\Theta(k \log (k / \delta))$ by [Buh+12; Sağ11; JW11] and [Hua+06]. We know that Disj ${ }_{k}^{n}$ can be solved much more efficiently if one is allowed many rounds: firstly the $\log k$ factor can be removed [HW07] and secondly the error probability can be brought down to $\exp (-k)$ [ST13], by using no more than $\log ^{*} k$ rounds. It is an interesting question whether similar efficiency improvements can be obtained for $\operatorname{Ham}_{k}^{n}$ also, by using multiple rounds. The first separation of $\mathrm{Disj}_{k}^{n}$ and $\operatorname{Ham}_{k}^{n}$ was proven in [BBG14], which shows that $\Omega(k \log (1 / \delta))$ lower bound holds for any protocol solving $\operatorname{Ham}_{k}^{n}$. Therefore in $\operatorname{Ham}_{k}^{n}$, we get no improvements in error probability by interactive communication. It remained an open question whether any improvement can be made at all to the 1-round protocol by communicating interactively. In this work we answer this question negatively:

Theorem 1.8 (restated). For $k^{2}<\delta$ n we have $R_{\delta}\left(\operatorname{Ham}_{k}^{n}\right)=\Omega(k \log (k / \delta))$. The bound applies
even to protocols that may output an arbitrary answer when $\|x-y\|_{1} \notin\{k-2, k, k+2\}$.
Before we proceed with proving Theorem 1.8, let us first warm up by showing that Theorem 1.3 implies an $\Omega(k \log (1 / \delta))$ lower bound on $R_{\delta}\left(\operatorname{Ham}_{k}^{n}\right)$. To do so, let us review the so called corruption bound method. Let $f: \mathcal{X} \times \mathcal{Y} \rightarrow\{0,1\}$ be the function the players would like to compute with Alice having received $x \in \mathcal{X}$ and Bob $y \in \mathcal{Y}$. For a protocol $P$ for $f$, define the matrix $A_{P}: \mathcal{X} \times \mathcal{Y} \rightarrow[0,1]$ such that $A_{P}(x, y)$ is the probability that the protocol outputs 1 on input $(x, y)$. It is well known and not difficult to see that if $P$ has communication cost $c$, then $A_{P}$ is the average of matrices each of which is the sum of at most $2^{c}$ rank 1 matrices $u v^{\boldsymbol{\top}}$ with $u \in\{0,1\}^{\mathcal{X}}$ and $v \in\{0,1\}^{\mathcal{Y}}$. Therefore to show the communication cost of a protocol $P$ is more than $c$, it suffices to argue $A_{P}$ lies outside $2^{c}$ times the polytope

$$
\mathcal{T}:=\operatorname{conv}\left\{u v^{\boldsymbol{\top}} \mid u \in\{0,1\}^{\mathcal{X}}, v \in\{0,1\}^{\mathcal{Y}}\right\}
$$

where conv denotes the convex hull. By convexity, $A_{p}$ lies outside of $\mathcal{T}$ if and only if there is a hyperplane (with normal $H$ ) separating the two; namely that $\left\langle A_{P}, H\right\rangle>2^{c}\langle R, H\rangle$ for all vertices $R$ of the polytope $\mathcal{T}$.

Let $\mu_{k}:\{0,1\}^{n} \times\{0,1\}^{n} \rightarrow \mathbb{R}_{+}$be the distribution on pairs $(x, y)$ obtained as follows. Sample $x$ uniformly at random and obtain $y$ by flipping $k$ coordinates of $x$ chosen uniformly at random and with replacement (here if a coordinate gets flipped twice it reverts back to its initial value).

Theorem 1.7 (restated). For $k^{2}<\delta n$ we have $R_{\delta}\left(\operatorname{Ham}_{k}^{n}\right)=\Omega(k \log (1 / \delta))$. The bound applies even to protocols that may output an arbitrary answer when $\|x-y\|_{1} \notin\{k, k+2\}$.

Proof. Suppose we have a randomized protocol for $\operatorname{Ham}_{k}^{n}$ with error probability $\delta$. Form the matrix $A$, where $A(x, y)$ is the probability that the protocol reports $\|x-y\|_{1} \leq k$ on input $(x, y)$.

Set $H=\mu_{k}-\mu_{k+2} /(3 \delta)$. Let us first argue that $\langle A, H\rangle \geq 1 / 3$. We have $\left\langle A, \mu_{k}\right\rangle>(1-$ $\delta)\left(1-\binom{k}{2} / n\right)>1-3 \delta / 2$, and $\left\langle P, \mu_{k+2}\right\rangle \leq \delta+\binom{k+2}{2} / n \leq 3 \delta / 2$. Hence $\langle A, H\rangle \geq 1 / 3$ for $\delta \leq 1 / 9$.

Next we argue that $\langle R, H\rangle<(3 \delta)^{k / 2}$ for any $R=u v^{\top}$ with $u, v \in\{0,1\}^{n}$. If $\left\langle R, \mu_{k}\right\rangle<(3 \delta)^{k / 2}$, we are done as $\left\langle R, \mu_{k+2}\right\rangle \geq 0$ and is a negative term in $\langle R, H\rangle$. If $\left\langle R, \mu_{k}\right\rangle \geq(3 \delta)^{k / 2}$ on the other hand, observing $\left\langle R, \mu_{k}\right\rangle=\left\langle v, W^{k} u\right\rangle / 2^{n}$, where $W$ is the normalized adjacency matrix of the Hamming cube, we have by Theorem 1.3

$$
\left(\frac{\|u\|_{2}\|v\|_{2}}{2^{n}}\right)^{2 / k}\left\langle v, W^{k+2} u\right\rangle \geq\left\langle v, W^{k} u\right\rangle^{1+2 / k}
$$

Note $\|u\|_{2}\|v\|_{2} \leq 2^{n}$ since $u, v$ are $0-1$ vectors, therefore $\left\langle R, \mu_{k+2}\right\rangle \geq 3 \delta\left\langle R, \mu_{k}\right\rangle$ and hence $\langle R, H\rangle \leq$ 0 . In either case we have shown that $\langle R, H\rangle<(3 \delta)^{k / 2}$. This implies an $\log \left((3 \delta)^{-k / 2} / 3\right)=$ $\Omega(k \log (1 / \delta))$ bits lower bound on $R_{\delta}\left(\operatorname{Ham}_{k}^{n}\right)$.

Interestingly, Theorem 1.8 cannot be proved by a direct application of the corruption method described above. If we assume that the protocol is supposed to output 1 on inputs $\|x-y\|_{1} \leq k$, then there are vertices of the polytope $\mathcal{T}$ for which the $\Omega(k \log (1 / \delta))$ bound of Theorem 1.7 is tight. If we assume that the protocol is supposed to output 1 on inputs $\|x-y\|_{1}>k$ on the other hand, no bound above $\Omega(k)$ can be obtained, as there are vertices for which this is tight. If we insist however that the protocol outputs 1 for $\|x-y\|_{1}=k$ and 0 for $\|x-y\|_{1} \in\{k-2, k+2\}$ then a protocol with cost smaller than $O(k \log (k / \delta))$ would be in violation of the near log-convexity principle we established in Theorem 1.4 as we argue next. Of course, if we had a $\delta$-error randomized protocol $P$ for $\operatorname{Ham}_{k}^{n+2}$ outputting 1 when $\|x-y\|_{1} \in\{k-2, k\}$ and 0 if $\|x-y\|_{1}=k+2$ (but without any guarantees for other types of inputs), then given inputs $a, b \in\{0,1\}^{n}$ Alice and Bob can run $P$ (say, in parallel) on instances $(00 a, 00 b)$ and $(00 a, 11 b)$ and declare $\|a-b\|_{1}=k$ if $P$ returns 1 on $(00 a, 00 b)$ and 0 on $(00 a, 11 b)$. This would lead to a protocol with twice the error probability and communication cost of $P$, deciding between $\|a-b\|_{1}=k,\|a-b\|_{1}=k-2$ and $\|a-b\|_{1}=k+2$. The table below shows that $P$ outputting 1 on $(00 a, 00 b)$ and 0 on $(00 a, 11 b)$ implies $\|a-b\|_{1}=k$ or at least one invocation of $P$ erred.

| Input | $k-2$ | $k$ | $k+2$ |
| ---: | :---: | :---: | :---: |
| $(00 a, 00 b)$ | 1 | 1 | 0 |
| $(00 a, 11 b)$ | 1 | 0 | $?$ |

Proof of Theorem 1.8. Suppose we have a $\delta$-error randomized protocol that outputs 1 when $\|x-y\|_{1}=k$ and 0 when $\|x-y\|_{1}=k-2$ or $\|x-y\|_{1}=k+2$.

Form the matrix $A$, where $A(x, y)$ is the probability that the protocol reports that $\|x-y\|_{1}=k$ on input $(x, y)$. Let $\alpha_{1}, \alpha_{2}>0$ be some reals so that Theorem 1.4 implies $m_{t+2} \geq t^{\alpha_{1}} m^{1+2 / t}$ or $m_{t+2} m_{t-2} \geq \alpha_{2} m_{t}^{2}$ for $m_{t}$ defined in statement of this theorem.

Set $H=\mu_{k}-\left(\mu_{k-2}+\mu_{k+2}\right) /(6 \delta)$. Let us argue first that $\langle A, H\rangle \geq 1 / 3$. One can verify that $\left\langle A, \mu_{k}\right\rangle \geq(1-\delta)\left(1-\binom{k}{2} / n\right)>1-3 \delta / 2$ and $\left\langle A, \mu_{k-2}\right\rangle+\left\langle A, \mu_{k-2}\right\rangle<3 \delta$. Hence $\langle A, H\rangle \geq 1 / 3$ for $\delta \leq 1 / 9$.

We upper bound $\langle R, H\rangle$ for some rank-1 matrix $R=u v^{\top}$ with $0-1$ values. Let $W$ be the normalized adjacency matrix of the Hamming cube graph. Observe that $\left\langle R, \mu_{k}\right\rangle=\left\langle v, W^{k} u\right\rangle / 2^{n}$. By Theorem 1.4, either $\left\langle R, \mu_{k+2}\right\rangle\left\langle R, \mu_{k-2}\right\rangle \geq \alpha_{2}\left\langle R, \mu_{k}\right\rangle^{2}$ or

$$
\left(\frac{\|u\|_{2}\|v\|_{2}}{2^{n}}\right)^{2 / k}\left\langle R, \mu_{k+2}\right\rangle \geq k^{\alpha_{1}}\left\langle R, \mu_{k}\right\rangle^{1+2 / k}
$$

In the former case,

$$
\frac{\left\langle R, \mu_{k+2}\right\rangle+\left\langle R, \mu_{k-2}\right\rangle}{2} \geq \sqrt{\left\langle R, \mu_{k+2}\right\rangle\left\langle R, \mu_{k-2}\right\rangle} \geq \sqrt{\alpha_{2}}\left\langle R, \mu_{k}\right\rangle
$$

which implies that $\langle R, H\rangle<0$ whenever $\delta<2 \sqrt{\alpha_{2}} / 6$ (recall $\alpha_{2}$ is a constant). In the latter case, recalling $\|v\|_{2}\|u\|_{2} \leq 2^{n}$, we get $\langle R, H\rangle<0$ unless $6 \delta\left\langle R, \mu_{k+2}\right\rangle \leq\left\langle R, \mu_{k}\right\rangle$, which implies that $k^{\alpha_{1}}\left\langle R, \mu_{k}\right\rangle^{2 / k}<6 \delta$. From this we get

$$
\langle R, H\rangle \leq\left\langle R, \mu_{k}\right\rangle \leq\left(\frac{6 \delta}{k^{\alpha_{1}}}\right)^{k / 2}
$$

and hence $\langle R, H\rangle<\left(\frac{6 \delta}{k^{\alpha_{1}}}\right)^{k / 2}$ in every case and $R_{\delta}\left(\operatorname{Ham}_{k}^{n}\right)=\Omega(k \log (k / \delta))$ whenever $k^{2}<\delta n$.
For a protocol $P$, denote by $\Pi=\Pi(x, y)$ the random variable entailing all the messages communicated between the players on input $(x, y)$. So far we have considered the communication cost of a protocol which is the maximum length of $\Pi$ over all inputs and the configurations of the random source (these together determine the value of $\Pi$ ). When a distribution $\mu$ on the inputs is available, we may speak of a more refined notion of cost, the internal information cost, for a protocol $P$ which is defined as

$$
\mathrm{IC}_{\mu}(P):=\mathbf{I}(\Pi: Y \mid X)+\mathbf{I}(\Pi: X \mid Y)
$$

where $(X, Y) \sim \mu$. Combining our Theorem 1.8 with a result of [Ker +15 ] which relates information and communication costs of a protocol under suitable circumstances, one can conclude that any randomized protocol for $\operatorname{Ham}_{k}^{n}$ has information cost $\Omega(k \log k)$ as well, under the distribution $\mu=\left(\mu_{k}+\mu_{k-2}+\mu_{k+2}\right) / 3$. However we note that instead of using Theorem 1.4 black-box, taking a closer look at the proof of Theorem 1.3 and not performing the relaxation provided in Lemma 3.6, we get the following more directly.

Theorem 5.2. Let $P$ be a protocol outputting 1 on pairs ( $x, y$ ) having $\|x-y\|_{1}=k$ with probability $1-\delta$ and outputting 0 on pairs $(x, y)$ having $\|x-y\|_{1} \in\{k-2, k+2\}$ with probability $1-\delta$. We have $\mathrm{IC}_{\mu_{k}}(P)=\Omega(k \log (k / \delta))$.

Let us finally mention another highly related problem, the so called the gap Hamming distance problem. In $\mathrm{GHD}_{k}^{n}$, each of the players receive a bit string, respectively $x, y \in\{0,1\}^{n}$, with the promise that either $\|x-y\|_{1} \leq k$ or $\|x-y\|_{1} \geq k+\sqrt{k}$. Their goal is to determine which is the
case for any given input. In [CR12], an $\Omega(k)$ lower bound for this problem was shown, which applies to protocols with any number of rounds. Here we conjecture an improvement to this bound and argue that it would follow from a natural analogue of of Theorem 1.4 for continuous time Markov chains, which we discuss in Section 6.

Conjecture 5.3. For $k<\delta n$, we have $R_{\delta}\left(\mathrm{GHD}_{k}^{n}\right)=\Omega(k \log (1 / \delta))$.

### 5.2 Parity decision trees

In the parity decision tree model, we are given a string $x \in \mathbb{F}_{2}^{n}$ and our goal is to determine whether $x$ satisfies a fixed predicate $P: \mathbb{F}_{2}^{n} \rightarrow\{0,1\}$ by only making linear measurements of the form $\langle x, y\rangle$ for some $y \in \mathbb{F}_{2}^{n}$ we get to choose. Here, the inner product is over $\mathbb{F}_{2}^{n}$, and therefore we get a single bit answer for every measurement we make.

Such measurements can be identified by binary decision trees wherein each internal node is labeled by a $y \in \mathbb{F}_{2}^{n}$ denoting the linear measurement $\langle x, y\rangle$ we would make at that node and each leaf is labeled by a YES or a NO denoting the final decision we arrive. Given such a tree and an $x$, the output of the decision tree is obtained by a root to leaf walk, where at each internal node $v$ with label $y_{v}$, we perform the measurement $\left\langle x, y_{v}\right\rangle$ and walk to the left child of $v$ if $\left\langle x, y_{v}\right\rangle=0$ and to the right child if $\left\langle x, y_{v}\right\rangle=1$. If a leaf node is reached, the label of the node is taken as the answer of the decision tree. Two quantities we are concerned with are the depth and the size (i.e., the total number of nodes) of the tree.

A $\delta$-error randomized decision tree is a distribution $\nu$ over deterministic trees such that for any fixed $x$, the sampled decision tree outputs the correct answer with probability at least $1-\delta$, where the randomness is over the choice of the decision tree from $\nu$. The depth and the size of a randomized decision tree can be taken as the maximum over the decision trees in the support of $\nu$ (here, one can also take the average depth or size, however this choice leads to negligible changes in our bounds).

For a predicate $P: \mathbb{F}_{2}^{n} \rightarrow\{0,1\}$, let $\mathrm{PD}_{\delta}(P)$ be the minimum, over all randomized decision trees $T$ computing $P$ with probability $1-\delta$, of the depth of $T$. Let $\mathrm{PS}_{\delta}$ be the minimum, over all randomized decision trees $T$ computing $P$ with probability $1-\delta$, of the size of $T$. The following inequalities are immediate

$$
\begin{align*}
R_{\delta}(P \circ \oplus) & \leq 2 \operatorname{PD}_{\delta}(P)  \tag{5.1}\\
\log \operatorname{PS}_{\delta}(P) & \leq \operatorname{PD}_{\delta}(P)
\end{align*}
$$

where $P \circ \oplus$ is the two player communication game in which the two players are given strings $x, y \in \mathbb{F}_{2}^{n}$ and are required to calculate $P(x+y)$. We remark that $\log \mathrm{PS}_{\delta}$ is incomparable to $R_{\delta}$ in general.

Here we study the predicate $H_{k}^{n}$ which equals 1 if and only if the Hamming weight of its input is precisely $k$. By Eq. (5.1) and a padding arguments similar to the one we gave before the proof of Theorem 1.8, each lower bound for $\operatorname{Ham}_{k}^{n}$ listed in Table 1 applies to $\mathrm{PD}_{\delta}\left(H_{k}^{n}\right)$ as well. In [BK12] another direct $\Omega(k)$ bound for $\mathrm{PD}_{\delta}\left(H_{k}^{n}\right)$ was shown. In [BCK14], showing an $\Omega(k \log k)$ lower bound to a variant of $\mathrm{PD}_{\delta}\left(H_{k}^{n}\right)$ to obtain tight bounds for $k$-linearity problem (see Section 5.3) was suggested. Finally, our Theorem 1.8 shows that $\mathrm{PD}_{\delta}\left(H_{k}^{n}\right)=\Omega(k \log (k / \delta))$, which is tight. Next we show the same bound holds even for $\log \mathrm{PS}_{\delta}\left(H_{k}^{n}\right)$.
Theorem 5.4. For $k^{2}<\delta n, \log \mathrm{PS}_{\delta}\left(H_{k}^{n}\right)=\Omega(k \log (k / \delta))$.
Proof. The proof is very similar to that of Theorem 1.8, so we only describe the differences.
Let $T$ be a $\delta$-error randomized parity decision tree computing $H_{k}^{n}$. Form $A: \mathbb{F}_{2}^{n} \rightarrow[0,1]$ so that $A(x)$ is the probability $T$ outputs 1 on input $x \in \mathbb{F}_{2}^{n}$. Define the polytope

$$
\mathcal{P}:=\operatorname{conv}\left\{x \mapsto \mathbb{1}[B x=c] \mid B \in \mathbb{F}_{2}^{n \times n}, c \in \mathbb{F}_{2}^{n}\right\}
$$

whose vertices are indicator functions for affine subspaces of $\mathbb{F}_{2}^{n}$. Given a parity decision tree, the set of inputs that end up in a particular leaf of it is an affine subspace in $\mathbb{F}_{2}^{n}$. Therefore if $T$ has at
most $s$ leaves, then $A$ is inside $s \mathcal{P}$. It remains to demonstrate a hyperplane with normal $H$ so that $\langle A, H\rangle>s\langle V, H\rangle$ for any vertex $V$ of the polytope $\mathcal{P}$ for $s=\exp \Omega(k \log (k / \delta))$.

Let $\mu_{k}$ be a distribution on $\mathbb{F}_{2}^{n}$ obtained as follows. Start with the 0 vector, and flip a coordinate chosen uniformly at random with replacement $k$ times. Here, flipping a coordinate an even number of times leaves it seemingly intact. Set $H=\mu_{k}-\left(\mu_{k-2}+\mu_{k+2}\right) /(6 \delta)$.

First observe that $\left\langle A, \mu_{k}\right\rangle>(1-\delta)\left(1-\binom{k}{2} / n\right)>1-3 \delta / 2$ and $\left\langle A, \mu_{k+2}\right\rangle+\left\langle A, \mu_{k-2}\right\rangle<3 \delta$ so $\langle A, H\rangle \geq 1 / 3$ for $\delta \leq 1 / 9$. Next we would like to upper bound $\langle V, H\rangle$ for an indicator function $V$ of an affine subspace $\left\{x \in \mathbb{F}_{2}^{n} \mid B x=c\right\}$. The key observation is

$$
\begin{equation*}
\left\langle V, \mu_{k}\right\rangle=\left\langle\mathbb{1}_{c}, S^{k} \mathbb{1}_{0}\right\rangle \tag{5.2}
\end{equation*}
$$

where $S$ is a stochastic matrix describing the following transition: For any $x \in \mathbb{F}_{2}^{n}$, sample a column $y$ of $B \in \mathbb{F}_{2}^{n \times n}$ uniformly at random and transition to $x+y$. Namely, the right hand side of Eq. (5.2) describes the following probability. We start with the 0 vector in $\mathbb{F}_{2}^{n}$ and in each time step sample a uniform random column $y$ of $B$ and add $y$ to the current state. We measure the probability of reaching $c \in \mathbb{F}_{2}^{n}$ at time step $k$. Having observed Eq. (5.2), and that $\left\|\mathbb{1}_{0}\right\|_{2}=\left\|\mathbb{1}_{c}\right\|_{2}=1$, the rest of the proof is identical to that of Theorem 1.8: by Theorem 1.4, we either have

$$
\left\langle\mathbb{1}_{c}, S^{k+2} \mathbb{1}_{0}\right\rangle\left\langle\mathbb{1}_{c}, S^{k-2} \mathbb{1}_{0}\right\rangle \geq \alpha_{2}\left\langle\mathbb{1}_{c}, S^{k} \mathbb{1}_{0}\right\rangle^{2}
$$

or

$$
\left\langle\mathbb{1}_{c}, S^{k+2} \mathbb{1}_{0}\right\rangle \geq k^{\alpha_{1}}\left\langle\mathbb{1}_{c}, S^{k} \mathbb{1}_{0}\right\rangle^{1+2 / k}
$$

In either event, we conclude that $\langle V, H\rangle \leq\left(\frac{6 \delta}{k^{\alpha_{1}}}\right)^{k / 2}$. This completes the proof.
Note in Theorem 1.8, we use Theorem 1.4 with a simple and fixed $S$ (i.e., the standard random walk on the Hamming cube), but with complicated vectors $u, v$ that come from the particular communication protocol whose communication cost we would like to lower bound. By contrast, in Theorem 5.4 the vectors $u, v$ are simple point masses on states 0 and $c$ but the matrix $S$ is a convolution random walk on the Hamming cube that comes from the particular decision tree whose size we lower bound.

### 5.3 Property testing

In the property testing model, given black box access to an otherwise unknown function $f: \mathbb{F}_{2}^{n} \rightarrow \mathbb{F}_{2}$, our goal is to tell apart whether $x \in P$ for some fixed set of functions $P$ or $\|f-g\|_{1} \geq \epsilon 2^{n}$ for any $g \in P$. Here, the black box queries are done by providing an input $x \in \mathbb{F}_{2}^{n}$ to the function and observing $f(x)$.

A function $f: \mathbb{F}_{2}^{n} \rightarrow \mathbb{F}_{2}$ is called $k$-linear if $f$ is given by

$$
f(x)=\sum_{i \in S} x_{i}
$$

for some $S \subseteq[n]$ of size at most $k$. By combining our communication complexity lower bound Theorem 1.8 with the reduction technique developed in [BBM12] or by combining our parity decision tree lower bound Theorem 5.4 with a reduction given in [BCK14], one obtains the following.

Theorem 5.5. Any $\delta$-error property testing algorithm for $k$-linearity requires $\Omega(k \log (k / \delta))$ queries when $\epsilon=1 / 2$.

In fact through this, one obtains similar lower bounds to property testing for $k$-juntas, $k$-term DNFs, size- $k$ formulas, size- $k$ decision trees, $k$-sparse $\mathbb{F}_{2}$-polynomials; see [Bla09; CGM11].

## 6 Discussion

We showed that for a symmetric matrix $S: \Omega \times \Omega \rightarrow \mathbb{R}_{+}$and unit vectors $u, v: \Omega \rightarrow \mathbb{R}_{+}$, defining $m_{t}=\left\langle v, S^{t} u\right\rangle$ for $t=0,1, \ldots$, we have

$$
\begin{align*}
& m_{t+2} \geq m_{t}^{1+2 / t}, \text { and }  \tag{6.1}\\
& m_{t+2} \geq m_{t}^{1+2 / t} \cdot \min \left\{t^{1-\epsilon}, \delta \frac{m_{t}^{1-2 / t}}{m_{t-2}}\right\} \tag{6.2}
\end{align*}
$$

and argued that Eq. (6.2) and (6.1), in this order, are best viewed as gradual weakenings of the log-convexity of $\left\{m_{t}\right\}_{t=0}^{\infty}$. We conjecture that a similar principle holds true for continuous time Markov chains as well.

Call a function $f: \mathbb{R}_{+} \rightarrow[0,1]$, whose logarithm is continuously twice differentiable (i.e., $\log f \in C^{2}\left(\mathbb{R}_{+}\right)$, nearly-log-convex if $x^{2}(\log f)^{\prime \prime}(x) \geq 2 \log f(x)$ for $x \in \mathbb{R}_{+}$. Note that $\log f \leq 0$, therefore this is a weakening of the usual $\log$-convexity definition, which requires $(\log f)^{\prime \prime} \geq 0$.

Conjecture 6.1. Let $S: \Omega \times \Omega \rightarrow \mathbb{R}_{+}$be a symmetric substochastic matrix and $u, v: \Omega \rightarrow \mathbb{R}_{+}$be unit vectors. The function

$$
t \mapsto\left\langle v, e^{t(S-I)} u\right\rangle
$$

is nearly-log-convex.
By an argument similar to the proof of Theorem 1.8, one can show the following.
Theorem 6.2. Conjecture 6.1 implies Conjecture 5.3.

## Acknowledgments

We are indebted to Paul Beame, Shayan Oveis Gharan and Gábor Tardos for many valuable conversations and the anonymous reviewers for bringing Conjecture 1.6 to our attention and for many improvements to the presentation of the paper. Many thanks to Suvrit Sra for brining [Pat12] to our attention [Sra16] and its connections to Sidorenko's conjecture.

## References

[AFH12] Anil Ada, Omar Fawzi, and Hamed Hatami. "Spectral Norm of Symmetric Functions". In: Approximation, Randomization, and Combinatorial Optimization. Algorithms and Techniques. Ed. by Anupam Gupta et al. Berlin, Heidelberg: Springer Berlin Heidelberg, 2012, pp. 338-349. ISBN: 978-3-642-32512-0 (cited on page 1).
$[A m b+15]$ Andris Ambainis et al. "Lower Bounds on the Deterministic and Quantum Communication Complexity of Hamming-Distance Problems". In: ACM Trans. Comput. Theory 7.3 (June 2015), 10:1-10:10. ISSN: 1942-3454. DOI: 10.1145/2698587. URL: http://doi.acm.org/10.1145/2698587 (cited on page 1).
[Bar +04$]$ Z. Bar-Yossef et al. "Approximating edit distance efficiently". In: 45th Annual IEEE Symposium on Foundations of Computer Science. Oct. 2004, pp. 550-559. DOI: 10.110 9/FOCS.2004.14 (cited on page 1).
[BBG14] Eric Blais, Joshua Brody, and Badih Ghazi. "The Information Complexity of Hamming Distance". In: Approximation, Randomization, and Combinatorial Optimization. Algorithms and Techniques (APPROX/RANDOM 2014). Ed. by Klaus Jansen et al. Vol. 28. Leibniz International Proceedings in Informatics (LIPIcs). Dagstuhl, Germany: Schloss Dagstuhl-Leibniz-Zentrum fuer Informatik, 2014, pp. 465-489. ISBN: 978-3-939897-74-3. DOI: 10.4230/LIPIcs.APPROX-RANDOM.2014.465. URL: http : //drops.dagstuhl.de/opus/volltexte/2014/4717 (cited on pages 1, 2, 4, 16).
[BBM12] Eric Blais, Joshua Brody, and Kevin Matulef. "Property Testing Lower Bounds via Communication Complexity". In: computational complexity 21.2 (June 2012), pp. 311358. ISSN: 1420-8954. DOI: 10.1007/s00037-012-0040-x. URL: https://doi.org/10. 1007/s00037-012-0040-x (cited on pages 1, 2, 20).
[BCK14] Abhishek Bhrushundi, Sourav Chakraborty, and Raghav Kulkarni. "Property Testing Bounds for Linear and Quadratic Functions via Parity Decision Trees". In: Computer Science - Theory and Applications. Ed. by Edward A. Hirsch et al. Cham: Springer International Publishing, 2014, pp. 97-110. ISBN: 978-3-319-06686-8 (cited on pages 2, 19, 20).
[BD66] GR Blakley and RD Dixon. "Hölder Type Inequalities in Cones". In: Journal of Mathematical Analysis and Applications 14 (1966), pp. 1-4. URL: http ://www. sciencedirect.com/science/article/pii/0022247X66900564 (cited on pages 24).
[BK12] Eric Blais and Daniel Kane. "Tight Bounds for Testing k-Linearity". In: Approximation, Randomization, and Combinatorial Optimization. Algorithms and Techniques. Ed. by Anupam Gupta et al. Berlin, Heidelberg: Springer Berlin Heidelberg, 2012, pp. 435-446. ISBN: 978-3-642-32512-0 (cited on pages 1, 2, 19).
[Bla09] Eric Blais. "Testing juntas nearly optimally". In: Proceedings of the forty-first annual ACM symposium on Theory of computing. ACM. 2009, pp. 151-158 (cited on page 20).
[BR65] G. R. Blakley and Prabir Roy. "A Hölder Type Inequality for Symmetric Matrices with Nonnegative Entries". In: Proceedings of the American Mathematical Society 16.6 (1965), pp. 1244-1245. ISSN: 00029939, 10886826. URL: http://www.jstor.org/ stable/2035908 (cited on page 4).
[Buh+12] Harry Buhrman et al. The non-adaptive query complexity of testing k-parities. 2012. eprint: arXiv:1209. 3849 (cited on pages 1, 2, 16).
[CGM11] Sourav Chakraborty, David Garcia-Soriano, and Arie Matsliah. "Efficient Sample Extractors for Juntas with Applications". In: Proceedings of the 38th International Colloquim Conference on Automata, Languages and Programming - Volume Part I. ICALP'11. Zurich, Switzerland: Springer-Verlag, 2011, pp. 545-556. ISBN: 978-3-642-22005-0. URL: http://dl.acm.org/citation.cfm?id=2027127. 2027185 (cited on page 20).
[Cor +00$]$ Graham Cormode et al. "Communication Complexity of Document Exchange". In: Proceedings of the Eleventh Annual ACM-SIAM Symposium on Discrete Algorithms. SODA '00. San Francisco, California, USA: Society for Industrial and Applied Mathematics, 2000, pp. 197-206. ISBN: 0-89871-453-2. URL: http://dl.acm.org/citation. cfm?id=338219. 338252 (cited on page 1).
[CR12] Amit Chakrabarti and Oded Regev. "An optimal lower bound on the communication complexity of gap-hamming-distance". In: SIAM Journal on Computing 41.5 (2012), pp. 1299-1317 (cited on page 19).
[CT06] Thomas M. Cover and Joy A. Thomas. Elements of information theory (2. ed.) Wiley, 2006, pp. I-XXIII, 1-748. ISBN: 978-0-471-24195-9 (cited on page 6).
[DKS12] Anirban Dasgupta, Ravi Kumar, and D. Sivakumar. "Sparse and Lopsided Set Disjointness via Information Theory". In: APPROX-RANDOM. 2012, pp. 517-528 (cited on page 16).
[ES82] P. Erdős and M. Simonovits. "Compactness results in extremal graph theory". In: Combinatorica 2.3 (Sept. 1982), pp. 275-288. ISSN: 1439-6912. DOI: 10.1007/BF025792 34. URL: https://doi.org/10.1007/BF02579234 (cited on page 4).
[Fis+02] Eldar Fischer et al. "Monotonicity Testing over General Poset Domains". In: Proceedings of the Thiry-fourth Annual ACM Symposium on Theory of Computing. STOC '02. Montreal, Quebec, Canada: ACM, 2002, pp. 474-483. ISBN: 1-58113-495-9. DOI: 10.1 145/509907.509977. URL: http://doi.acm.org/10.1145/509907. 509977 (cited on page 2).
[Gib02] Josiah Willard Gibbs. Elementary Principles in Statistical Mechanics. Cambridge Books Online. Cambridge University Press, 1902. ISBN: 9780511686948. URL: http: //dx.doi.org/10.1017/CB09780511686948 (cited on page 5).
[GKW04] Dmitry Gavinsky, Julia Kempe, and Ronald de Wolf. "Quantum Communication Cannot Simulate a Public Coin". In: CoRR quant-ph/0411051 (2004). URL: http: //arxiv.org/abs/quant-ph/0411051 (cited on page 1).
[Gol10] Oded Goldreich. "On Testing Computability by Small Width OBDDs". In: Proceedings of the 13th International Conference on Approximation, and 14 the International Conference on Randomization, and Combinatorial Optimization: Algorithms and Techniques. APPROX/RANDOM'10. Barcelona, Spain: Springer-Verlag, 2010, pp. 574-587. ISBN: 3-642-15368-2. URL: http://dl.acm.org/citation.cfm?id=1886521.1886566 (cited on page 2).
[Hau21] F. Hausdorff. "Summationsmethoden und Momentfolgen I." ger. In: Mathematische Zeitschrift 9 (1921), pp. 74-109. URL: http://eudml.org/doc/167613 (cited on page 2).
[Hua+06] Wei Huang et al. "The communication complexity of the Hamming distance problem". In: Information Processing Letters 99.4 (2006), pp. 149-153. ISSN: 0020-0190. DOI: h ttp://dx.doi.org/10.1016/j.ipl.2006.01.014. URL: http://www.sciencedirect.com/ science/article/pii/S002001900600072X (cited on pages 1, 15, 16).
[HW07] Johan Håstad and Avi Wigderson. "The Randomized Communication Complexity of Set Disjointness". In: Theory of Computing 3.1 (2007), pp. 211-219 (cited on page 16).
[Jen06] J. Jensen. "Sur les fonctions convexes et les inégalités entre les valeurs moyennes". In: Acta Mathematica 30 (1 Dec. 1906), pp. 175-193. ISSN: 0001-5962. URL: http: //dx.doi.org/10.1007/BF02418571 (cited on page 5).
[JW11] T. S. Jayram and David P. Woodruff. "Optimal Bounds for Johnson-Lindenstrauss Transforms and Streaming Problems with Sub-Constant Error". In: SODA. 2011, pp. 1-10 (cited on page 16).
[Ker +15$] \quad$ Iordanis Kerenidis et al. "Lower bounds on information complexity via zero-communication protocols and applications". In: SIAM Journal on Computing 44.5 (2015), pp. 15501572 (cited on page 18).
[KL51] S. Kullback and R. A. Leibler. "On Information and Sufficiency". In: Ann. Math. Statist. 22.1 (Mar. 1951), pp. 79-86. DOI: 10.1214/aoms/1177729694. URL: http: //dx.doi.org/10.1214/aoms/1177729694 (cited on page 5).
[KS92] Bala Kalyanasundaram and Georg Schnitger. "The Probabilistic Communication Complexity of Set Intersection". In: SIAM J. Discrete Math. 5.4 (1992), pp. 545-557 (cited on page 16).
[Lon66] David London. "Inequalities in quadratic forms". In: Duke Math. J. 33.3 (Sept. 1966), pp. 511-522. DOI: 10.1215/S0012-7094-66-03360-6. URL: https://doi.org/10.1215/ S0012-7094-66-03360-6 (cited on page 4).
[MH58] SPH Mandel and IM Hughes. "Change in mean viability at a multiallelic locus in a population under random mating". In: Nature 182.4627 (1958), p. 63 (cited on pages 2, 4).
[MS59] H. P. Mulholland and C. A. B. Smith. "An Inequality Arising in Genetical Theory". In: The American Mathematical Monthly 66.8 (1959), pp. 673-683. ISSN: 00029890, 19300972. URL: http://www.jstor.org/stable/2309342 (cited on page 4).
[NP05] C. Niculescu and L.E. Persson. Convex Functions and their Applications: A Contemporary Approach. CMS Books in Mathematics. Springer New York, 2005. ISBN: 9780387243009. URL: https://books.google.com.tr/books?id=xpT\\_K190S24C (cited on page 2).
[Pat12] Thomas Pate. "Extending the Hölder type inequality of Blakley and Roy to nonsymmetric non-square matrices". In: Transactions of the American Mathematical Society 364.8 (2012), pp. 4267-4281 (cited on pages 4, 21).
[PG86] King F Pang and Abbas El Gamal. "Communication Complexity of Computing the Hamming Distance". In: SIAM J. Comput. 15.4 (Nov. 1986), pp. 932-947. Issn: 00975397. DOI: 10.1137/0215065. URL: http://dx.doi.org/10.1137/0215065 (cited on page 1).
[Sağ11] Mert Sağlam. "Tight bounds for data stream algorithms and communication problems". MA thesis. 8888 University Drive, Burnaby, B.C. Canada V5A 1S6: Simon Fraser University, School of Computing Science, Sept. 2011 (cited on page 16).
[She12] Alexander A. Sherstov. "The Communication Complexity of Gap Hamming Distance". In: Theory of Computing 8.8 (2012), pp. 197-208. DOI: 10.4086/toc.2012.v008a008. URL: http://www.theoryof computing.org/articles/v008a008 (cited on page 4).
[Sra16] Suvrit Sra. A conjecture of Blakley and Dixon about odd powers of positive matrices. https://mathoverflow.net/questions/236420/a-conjecture-of-blakley-and-dixon-about-odd-powers-of-positive-matrices. Accessed: 2018-08-15. 2016 (cited on page 21).
[ST13] Mert Sağlam and Gábor Tardos. "On the communication complexity of sparse set disjointness and exists-equal problems". In: CoRR abs/1304.1217 (2013) (cited on pages 1, 16).
[Wal45] A. Wald. "Sequential Tests of Statistical Hypotheses". In: Ann. Math. Statist. 16.2 (June 1945), pp. 117-186. DOI: 10.1214/aoms/1177731118. URL: http://dx.doi.org/ 10.1214/aoms/1177731118 (cited on page 5).
[Yao03] Andrew Chi-Chih Yao. "On the Power of Quantum Fingerprinting". In: Proceedings of the Thirty-fifth Annual ACM Symposium on Theory of Computing. STOC '03. San Diego, CA, USA: ACM, 2003, pp. 77-81. ISBN: 1-58113-674-9. DOI: 10.1145/780542.78 0554. URL: http://doi.acm.org/10.1145/780542.780554 (cited on page 1).


[^0]:    ${ }^{1}$ Since $u=v$ here, the summands inside Eq. (1.3) are nonnegative when $t$ is even so this theorem is most interesting for $t$ odd.

