Probabilistic Checking against Non-Signaling Strategies from Linearity Testing

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Abstract

Non-signaling strategies are a generalization of quantum strategies that have been studied in physics over the past three decades. Recently, they have found applications in theoretical computer science, including to proving inapproximability results for linear programming and to constructing protocols for delegating computation. A central tool for these applications is probabilistically checkable proofs (PCPs) that are sound against non-signaling strategies.

In this paper we prove that the exponential-length constant-query PCP construction due to Arora et al. (JACM 1998) is sound against non-signaling strategies.

Our result offers a new length-vs-query tradeoff when compared to the non-signaling PCP of Kalai, Raz, and Rothblum (STOC 2013 and 2014) and, moreover, may serve as an intermediate step to a proof of a non-signaling analogue of the PCP Theorem.

Keywords: probabilistically checkable proofs; linearity testing; non-signaling strategies
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1 Introduction

Probabilistically Checkable Proofs (PCPs) \cite{BFLS91, FGLSS96, AS98, ALMSS98} are proofs whose validity can be checked by a probabilistic verifier that accesses only a few locations of the proof. PCPs have numerous applications across the theory of computing, including to hardness of approximation \cite{FGLSS96} and delegation of computation \cite{Kib92, Mic00}.

A seminal result, known as the PCP Theorem \cite{AS98, ALMSS98}, states that every language in NTIME(T) can be probabilistically checked by a verifier that uses O(log T) random bits and makes O(1) queries to a proof of length \text{poly}(T).\footnote{In particular, for every language in NEXP = \bigcup_{k \in \mathbb{N}} \text{NTIME}(2^{n^k}), the verifier uses \text{poly}(n) random bits and makes O(1) queries to a proof of length 2^{\text{poly}(n)}.}

In this paper we study PCPs that are sound against non-signaling strategies (nsPCPs). These have recently found applications that appear out of the reach of (standard) PCPs, including 1-round delegation of computation from falsifiable assumptions \cite{KRR13, KRR14} and hardness of approximation for linear programming \cite{KRR16}. The efficiency measures achieved in known nsPCPs appear suboptimal, which affects the quality of the corresponding applications. We thus ask whether a non-signaling analogue of the PCP Theorem holds.

Below we explain the aforementioned notions, and then present our results in this direction.

Non-signaling strategies. Non-signaling strategies are a class of “non-local” correlations that strictly generalize quantum strategies, and capture the minimal condition that spatially-isolated parties cannot communicate instantaneously. They have been studied in physics for over three decades \cite{Ras85, KT92, PR94} in order to better understand quantum entanglement.

There are two definitions, corresponding to whether the strategy is meant to represent a function or isolated parties; the former is the relevant one for nsPCPs \cite{KRR13, KRR14}.$^2$ Given a locality parameter $k \in \mathbb{N}$, a $k$-non-signaling function $\mathcal{F}$ extends the notion of a function $f : D \rightarrow \{0, 1\}$ as follows: it is a collection $\{\mathcal{F}_S\}_{S \subseteq D, |S| \leq k}$ where each $\mathcal{F}_S$ is a distribution over $\{0, 1\}^S$ and, for every two subsets $S_1$ and $S_2$ each of size at most $k$, the restrictions of $\mathcal{F}_{S_1}$ and $\mathcal{F}_{S_2}$ to $S_1 \cap S_2$ are equal as distributions.$^3$ Note that if $k = |D|$ then $\mathcal{F}$ is a distribution over functions $f : D \rightarrow \{0, 1\}$.

Note that $k$-non-signaling functions are solutions to the linear program arising from the $k$-relaxation in the Sherali–Adams hierarchy \cite{SA90}. The variables are of the form $X_{S, \vec{b}}$ (for all $S \subseteq D$ of size at most $k$ and $\vec{b} \in \{0, 1\}^S$) and express the probability of $\vec{b}$ in the distribution $\mathcal{F}_S$; consistency across subsets $S$ and $T$ is expressed using the natural linear constraints.$^4$

Non-signaling PCPs. Recall that a classical PCP verifier is given oracle access to a proof represented as a function $f : D \rightarrow \{0, 1\}$. The verifier uses random bits, makes a few queries to $f$, and then accepts or rejects. Completeness requires that if the statement being checked is true then there is a function $f$ that makes the verifier always accept. Soundness requires that if the statement being checked is false then every function $f$ makes the verifier reject with high probability.

In the non-signaling setting, “proofs” are non-signaling functions rather than (classical) functions. Soundness is correspondingly stronger: given a locality parameter $k \in \mathbb{N}$, soundness requires that

\footnote{The other definition underlies the notion of multi-prover interactive proofs that are sound against non-signaling strategies (nsMIPs). Any nsPCP gives rise to an nsMIP with similar parameters. See \cite{KRR13, KRR14} for details.}

\footnote{A common relaxation of this condition only requires that the marginals $\mathcal{F}_{S_1|S_1 \cap S_2}$ and $\mathcal{F}_{S_2|S_1 \cap S_2}$ are statistically close, rather than equal; a further relaxation only requires these to be computationally close. While we only consider the standard definition (the marginals must equal) we note that \cite{CMS18} shows that this is almost without loss of generality, as every statistically or computationally non-signaling strategy is close to an (exact) non-signaling strategy.}

\footnote{In fact it suffices to only have variables of the form $X_{S_1|S}$ since all other probabilities can be computed from these.}
every $k$-non-signaling function $F$ makes the nsPCP verifier reject with high probability.

Efficiency measures of a nsPCP include familiar notions such as proof length (defined as $|D|$) and the verifier’s randomness and query complexity. In addition, the locality parameter $k$ controls how hard it is to attain soundness: the smaller $k$ is, the larger the set of non-signaling functions that the verifier could face. (Note that $k$-non-signaling implies $(k-1)$-non-signaling.)

There is a qualitative difference between the complexity classes captured by PCPs and by nsPCPs; namely, while PCPs capture non-deterministic time languages, nsPCPs capture deterministic ones. Indeed, the aforementioned PCP Theorem implies that it is $\text{NEXP}$-hard to approximate the maximum acceptance probability of a PCP verifier (that uses polynomial randomness). In contrast, computing the maximum acceptance probability of an nsPCP verifier that uses $r$ random bits reduces to a linear program with $2^{\text{poly}(rk)}$ variables and constraints, a problem solvable in $\text{EXP} = \cup_{c \in \mathbb{N}} \text{DTIME}(2^{n^c})$.

If $k = 2$, the linear program is solvable in PSPACE [Ito10], which is a tight upper bound [IKM09]. For $k > 2$ little is known, except for a seminal result of Kalai, Raz, and Rothblum [KRR13; KRR14], which shows that for $k = \text{poly}(n)$ it is EXP-hard to approximate a nsPCP verifier’s maximum acceptance probability. In more detail, every language in $\text{DTIME}(T)$ has a verifier that uses $\text{poly}(\log T)$ random bits and makes $\text{poly}(\log T)$ queries to a proof of length $\text{poly}(T)$; soundness holds against $\text{poly}(\log T)$-non-signaling functions; the verifier runs in time $n \cdot \text{poly}(\log T)$ and space $\text{poly}(\log T)$.

### The nsPCP Conjecture

The nsPCP construction behind the above result is a whitebox modification of early PCP constructions [BFL91; BFLS91], and achieves efficiency similar to those. However, modern “PCP technology” goes well beyond these early constructions, via tools such as proof composition [AS98] and proofs of proximity [DR04; BGHSV06], and enables better efficiency, including the PCP Theorem. Yet, current “nsPCP technology” is limited to the above results, and the question of whether a non-signaling analogue of the PCP Theorem holds remains open.

**Question 1.1.** Is it true that every language in $\text{DTIME}(T)$ has an nsPCP verifier that uses $O(\log T)$ random bits, makes $O(1)$ queries, and is sound against $O(1)$-non-signaling functions? (As above, we also require that the verifier runs in time $n \cdot \text{poly}(\log T)$ and space $\text{poly}(\log T)$.)

An affirmative answer to the above question would, e.g., improve the hardness result for linear programming in [KRR16], by yielding a reduction that outputs a linear program of polynomial, rather than a quasipolynomial, size. While we do not know if an affirmative answer exists (and we cannot prove that it does not exist), it is clear that the (very few) tools that we have to construct and analyze nsPCPs are far from this goal. In this paper we make headway towards this goal.

### 1.1 Towards a nsPCP Theorem

In [ALMSS98], a key step towards the PCP Theorem is to prove a weaker result in which the proof has exponential, rather than polynomial, size (and so the randomness complexity of the verifer is polynomial rather than logarithmic). Namely, one proves that every language in $\text{NTIME}(T)$ has a PCP verifier that uses $\text{poly}(T)$ random bits and makes $O(1)$ queries to a proof of length $2^{\text{poly}(T)}$.

In this paper we ask whether a non-signaling analogue of this result holds for the class $\text{DTIME}(T)$.

**Question 1.2.** Is it true that every language in $\text{DTIME}(T)$ has an nsPCP verifier that uses $\text{poly}(T)$ random bits, makes $O(1)$ queries, and is sound against $O(1)$-non-signaling functions?

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5Achieving time and space complexities that are $o(T)$ is important for applications. This is not surprising as every language in $\text{DTIME}(T)$ has a trivial nsPCP verifier that runs in time $T$: the verifier that simply decides the language, without asking any queries. This is unlike the case of PCPs for $\text{NTIME}(T)$, where time complexity is less critical.
We propose this question as a relaxation that, not only is interesting in its own right, but is likely to shed light on Question 1.1. However, one must be careful with the precise formulation of Question 1.2. If the verifier can use $\text{poly}(T)$ random bits then it can simply decide the language by running in time $T$, without making any queries. To recover a nontrivial question, we require that in order to decide whether an instance $x$ is in a language $L \in \text{DTIME}(T)$ the nsPCP verifier first generates queries via a $\text{poly}(T)$-time sampler that is input oblivious (knows the length of $x$ but not $x$ itself), and then rules according to a $o(T)$-time decision predicate that knows $x$. We stress that all PCP/nsPCP verifiers discussed in this paper are input oblivious.

In this paper we study Question 1.2 by analyzing a natural candidate construction, and ask:

Is the exponential-length $O(1)$-query PCP of [ALMSS98] sound against $O(1)$-non-signaling functions?

Hereafter, we consider the complexity class $\text{DSIZE}(S)$ (languages decidable by uniform circuits of size $S(n)$) instead of the class $\text{DTIME}(T)$ (languages decidable by machines in time $T(n)$) because our results, like their classical counterparts, are most easily stated in terms of uniform circuits. This change is only for simplicity, as $\text{DTIME}(T) \subseteq \text{DSIZE}(\text{poly}(T))$.

### 1.2 Main theorem

In this paper we prove that the exponential-length constant-query PCP construction of [ALMSS98] (without modifications) is sound against non-signaling functions. We obtain the following theorem.

**Theorem 1** (main theorem). Every language $L \in \text{DSIZE}(S)$ has an input-oblivious nsPCP verifier that uses $O(S^2)$ random bits, makes 11 queries, and is sound against $O(\log^2 S)$-non-signaling functions. The query sampler runs in time $O(S^2)$, and the decision predicate runs in time $O(n)$.

The theorem is close to answering Question 1.2, which asks for soundness against $O(1)$-non-signaling functions. (See Fig. 1 for a comparison with the classical result on nondeterministic languages.) At the same time, some may consider Ito’s algorithm [Ito10] as evidence that soundness against $O(1)$-non-signaling functions is too much to hope for. Understanding this gap needs further research.

Our result is incomparable to the nsPCP of [KRR13; KRR14], where the nsPCP verifier uses $\text{poly}(\log S)$ random bits to make $\text{poly}(\log S)$ queries. The fact that we prove soundness only against $O(\log^2 S)$-non-signaling functions (rather than $O(1)$-non-signaling functions) is somewhat undesirable, as this implies that the corresponding nsMIP requires $O(\log^2 S)$ provers. That said, the nsMIP of [KRR13; KRR14] requires many more provers: $\text{poly}(\log S)$ with the degree in the polynomial much larger than 2. Another feature of our result is that we have “room” to achieve smaller soundness error without using additional provers; for example, by asking more queries to the $O(\log^2 S)$ provers, we can achieve a sub-constant soundness error of $2^{-O(\log S)}$.

Finally, our result is the first to demonstrate that a classical PCP construction is secure against non-signaling functions, without any modifications. This should be compared to the construction considered in [KRR13; KRR14] that, while modeled after the PCP in [BFL91; BFLS91], includes several notable modifications that are needed in the soundness proof.

### 1.3 Main lemmas

We outline the ideas behind our theorem in Section 2. Concretely, we highlight several statements, which we deem of independent interest, that we prove on the way to the theorem.
Recall that the exponential-length constant-query PCP in [ALMSS98] is obtained in two steps. First, construct a constant-query verifier where soundness holds as long as the proof string is a linear function; this is known as a linear PCP. Second, use a linearity test [BLR93] and self-correction to compile this linear PCP into a (standard) PCP, where soundness holds against arbitrary proofs.

Our approach follows the same two steps, but adapted to the non-signaling setting. This also departs from the approach in [KRR14], which does not make use of any property testing results.

Note, however, that it is a priori not clear what is the non-signaling analogue of a linear function. A natural attempt would be to say that a non-signaling function \( F \) is linear iff it passes the BLR linearity test with probability 1 (where the probability is over the test and \( F \)). But this attempt is awkward, because the definition depends on a local test, and avoids discussing “global” structure.

A recent work [CMS18] tells us that the right definition is to say that \( F \) is linear iff it corresponds to a quasi-distribution over linear functions. A quasi-distribution is a probability distribution where the weights can be any real number and are not restricted to be in \([0, 1]\). Quasi-distributions over functions arise in this context because they are an equivalent description of non-signaling functions.

In light of the above, the notion of a non-signaling linear PCP (nsLPCP) is immediate: the definition requires soundness to hold against all linear non-signaling functions.

The first step in our proof is showing that the linear PCP verifier of [ALMSS98] (the “ALMSS verifier”), when used for deterministic computations, is sound against linear non-signaling functions.

**Theorem 2.** The (input oblivious) ALMSS verifier, for a given language \( L \in \text{DSIZE}(S) \), uses \( O(S) \) random bits, makes 4 queries, and is sound against linear \( O(\log S) \)-non-signaling functions.

See Fig. 1 for a comparison with the classical result showing soundness against linear functions. In order to “lift” Theorem 2 to Theorem 1, we need a suitable linearity test.

The linearity test of [BLR93] was recently analyzed in the non-signaling setting by [CMS18], who proved that any \( k \)-non-signaling function \( F \) that passes the linearity test with probability \( 1 - \varepsilon \) can be self-corrected to a \( \lfloor k/2 \rfloor \)-non-signaling function \( \hat{F} \) that is \( 2^{O(k)} \varepsilon \)-close to a linear \( \lfloor k/2 \rfloor \)-non-signaling function \( L \). (Self-correction and closeness have precise meanings, discussed later.) However, we cannot directly use [CMS18]’s result, because in our theorem the locality parameter \( k \) is required to be super-constant (\( k = O(\log S) \) in Theorem 2), and thus the bound on the distance between \( \hat{F} \) and \( L \) is too large, even when considering only query sets of small size. Specifically, we need the distance to be a sufficiently small constant on query sets of size 4 (the number of queries in Theorem 2).

We solve this problem by extending the result in [CMS18] in a black-box way and proving that the distance between \( \hat{F} \) and \( L \) on a query set \( Q \) is only \( O(|Q| \sqrt{\varepsilon}) \), provided that the error \( \varepsilon \) and \( L \)’s locality are sufficiently small. Crucially, if \( |Q| \) is constant, so is the distance between \( \hat{F} \) and \( L \).
The proof of this statement involves analyzing the repeated linearity test, whose behavior in the non-signaling setting is quite subtle when compared to the classical setting (see Section 2.6).

**Theorem 3.** Let $k, \bar{k} \in \mathbb{N}$ and $\varepsilon \in (0, 1/400]$ be such that $k = \Omega((\bar{k} + \log \frac{1}{\varepsilon}) \cdot \bar{k})$. If a $k$-non-signaling function $F: \{0, 1\}^n \rightarrow \{0, 1\}$ passes the linearity test with probability at least $1 - \varepsilon$ then there exists a linear $\bar{k}$-non-signaling function $L: \{0, 1\}^n \rightarrow \{0, 1\}$ such that for all query sets $Q \subseteq \{0, 1\}^n$ with size $|Q| \leq \bar{k}$ and for all events $E \subseteq \{0, 1\}^Q$ it holds that

$$\left| \Pr[\hat{F}(Q) \in E] - \Pr[L(Q) \in E] \right| \leq O(|Q| \sqrt{\varepsilon}) .$$

The above result on linearity testing enables us to transform our nsLPCP, and more generally any nsLPCP, into a corresponding nsPCP with minimal changes in parameters (the transformation is exactly the classical compiler). This is the last key statement in the proof of our main theorem.

**Lemma 1.3.** For every $\varepsilon \in [0, 1]$, if a language $L$ has an nsLPCP where the verifier uses $r$ random bits, makes $q$ queries, and has soundness error $1 - \varepsilon$ against linear $k$-non-signaling functions $L: \{0, 1\}^\ell \rightarrow \{0, 1\}$, then $L$ has an nsPCP where the verifier uses $r + O(q\ell)$ random bits, makes $O(q)$ queries, and has soundness error $1 - O_q(\varepsilon^2)$ against $O_\varepsilon(k^2)$-non-signaling functions $F: \{0, 1\}^\ell \rightarrow \{0, 1\}$. (Furthermore, if the former is input oblivious, so is the latter.)

### 1.4 Enriching the toolkit for non-signaling PCPs

Progress in our understanding of PCPs has typically moved hand in hand with progress in our understanding of low-degree testing. In particular, many PCP constructions follow this blueprint: (1) a low-degree test that, via only a few queries, ensures that a given proof conforms to a specified algebraic encoding; (2) a probabilistic test that, assuming the proof is (essentially) given in this encoding, ensures that the statement being checked is true with high probability.

In contrast, while the nsPCP in [KRR14] is reminiscent of this blueprint, its analysis does not follow it, despite the fact that the construction is modeled after the PCP in [BFL91; BFLS91], for which the two-step analysis is possible (in the classical setting). The lack of such general paradigms means that we lack general design principles to construct better nsPCPs.

This state of affairs raises the intriguing question of whether a theory of low-degree testing (and, more generally, property testing) is feasible in the non-signaling setting and, moreover, whether one can build on it to construct nsPCPs in order to make further progress towards Question 1.1.

An additional contribution of our work is to enrich the current “non-signaling toolkit”, by demonstrating an example where the aforementioned blueprint is both possible and useful.

Namely, building on the work of [CMS18] on linearity testing, our results provide a modular paradigm that not only simplifies the overall analysis, thereby enabling us to assert that the construction of [ALMSS98] with no modifications is sound against non-signaling strategies, but also (as discussed later) clarifies the technical barriers that separate us from answering Question 1.2. All this suggests that our techniques will be helpful for constructing more efficient nsPCPs.

We remark that an analogous modular approach has also been successful in the setting of quantum strategies. Indeed, Ito and Vidick [IV12] construct an MIP for NEXP that is sound against entangled provers, and their proof works by first analyzing the multilinearity test against quantum adversaries, and then building on their testing result to prove soundness of the proof system.
1.5 Concurrent work

In a concurrent work, Kiyoshima [Kiy18] studies the soundness of ALMSS-type PCPs against non-signaling strategies. Kiyoshima proves that, for a sufficiently large security parameter $t$ (at least logarithmic in the circuit size), the $t$-repetition of a $O(t)$-query modification of the ALMSS-verifier has soundness error $\text{negl}(t)$ against $O(t^2)$-non-signaling functions. In comparison, we prove that the unmodified 11-query ALMSS PCP has soundness error $O(1)$ against $O(\log^2 S)$-non-signaling strategies (and also that a modification of its $t$-repetition has soundness error $\exp(-t)$ for every $t = \Omega(\log S)$). While both our analysis and Kiyoshima’s analysis avoid the use of an augmented circuit (necessarily so as it would have had exponential size), our techniques differ. Kiyoshima conducts a direct analysis of the PCP verifier, while we adopt a modular approach in which we first prove soundness against linear non-signaling strategies (a simpler task), and then, building on a recent analysis of the linearity test [CMS18], we deduce soundness against all non-signaling strategies. We consider the modular and simple analysis in our work to be of independent interest. Kiyoshima additionally proves that soundness holds against computational non-signaling strategies, a relaxation where the marginal distributions on intersections are only required to be computationally close. Our results directly extend to computational non-signaling strategies as every computational non-signaling strategy is close to an exact non-signaling strategy (as proved in [CMS18]).

In another concurrent work, Holmgren and Rothblum [HR18] study the problem of constructing PCPs/MIPs in which the prover is very efficient in time and space [BC12], in the non-signaling setting. While they consider a construction that is more closely related to the PCP in [BFL91; BFLS91] (honest proofs are encoded via low-degree polynomials rather than linear functions), their soundness analysis also has the feature that it avoids the use of an augmented circuit.

1.6 Open problems

The question of whether the exponential-length constant-query PCP of [ALMSS98] is sound against $O(1)$-non-signaling functions remains open. A concrete approach to affirmatively answer this question is to prove that the linear PCP verifier of [ALMSS98] is sound against $k$-non-signaling functions for $k = O(1)$, rather than $k = O(\log S)$ as in Theorem 2. (Our generic compiler from Lemma 1.3 would then take care of the rest.) Another intriguing possibility is that an affirmative answer to Question 1.2 could come from a different exponential-size constant-query PCP. However, the result due to [Ito10] shows that the class of nsMIPs with 2 provers equals $\text{PSPACE}$, which possibly suggests that soundness against $O(1)$-non-signaling functions is too much to hope for.

Moreover, while our results can be interpreted as progress towards a non-signaling analogue of the PCP Theorem (Question 1.1), it remains unclear whether such an analogue holds, and more investigations in nsPCPs are needed. We believe that our work and our new techniques can inform such investigations.
2 Techniques

We outline the techniques used to prove our results. First, in Section 2.1, we explain the transformation from a nsLPCP to a corresponding nsPCP. Next, in Sections 2.2 to 2.5 we discuss the nsLPCP on which we apply this transformation, namely, the ALMSS verifier [ALMSS98]. Finally, in Section 2.6, we discuss linearity testing with low error, which underlies the transformation.

2.1 From nsLPCP to nsPCP

We discuss the transformation from nsLPCP to nsPCP (Lemma 1.3). We first recall the classical transformation from LPCP to PCP, and then explain how to achieve its non-signaling analogue.

**The classical case.** The classical transformation from LPCP to PCP relies on the following tools.

- **Testing linearity.** Given a boolean function \( f : \{0,1\}^\ell \rightarrow \{0,1\} \), the linearity test draws random \( x, y \in \{0,1\}^\ell \) and checks that \( f(x) + f(y) = f(x+y) \) [BLR93]. If the test passes with probability \( 1 - \varepsilon \), then \( f \) is \( \varepsilon \)-close to a linear function \( f^* : \{0,1\}^\ell \rightarrow \{0,1\} \) [BLR93] [BCHKS96].

- **Self-correction.** Given \( f \) that is \( \varepsilon \)-close to a linear function \( f^* \), one can create a probabilistic oracle \( O \) that, given any \( x \in \{0,1\}^\ell \), returns \( f^*(x) \) with probability \( 1 - 2\varepsilon \). Namely, \( O \) samples a random \( z \in \{0,1\}^\ell \), queries \( f \) on \( z+x \) and \( z \), and answers with \( f(z+x) - f(z) \).

The above tools imply a transformation from LPCP to PCP: given access to an arbitrary function \( f : \{0,1\}^\ell \rightarrow \{0,1\} \), the PCP verifier runs the linearity test and then runs the LPCP verifier by self-correcting each of its queries. If the LPCP verifier makes \( q \) queries and has soundness error \( \gamma \), then the resulting PCP verifier makes \( 3 + 2q \) queries and has soundness error \( \max\{1 - \varepsilon, \gamma + 2q\varepsilon\} \), where \( \varepsilon \) is (a bound on) the distance of \( f \) to linear functions. This soundness error is bounded by \( 1 - \frac{1}{2^q+1} \) (the maximum is when the two terms equal), which is bounded away from 1.

If desired, the soundness error can be made arbitrarily close to \( \gamma \) by repeating the linearity test. Given a parameter \( t \), the repeated linearity test samples \( x_i, y_i \in \{0,1\}^\ell \) for each \( i \in [t] \) and checks that \( f(x_i) + f(y_i) = f(x_i+y_i) \) for all \( i \in [t] \). Now, the PCP verifier makes \( 3t + 2q \) queries and has soundness error \( \max\{1 - \varepsilon^t, \gamma + 2q\varepsilon\} \), which for suitable \( \varepsilon \) and \( t = O_{\gamma,\varepsilon}(q) \) is arbitrarily close to \( \gamma \).

The non-signaling case. We follow the structure of the classical transformation. However, the non-signaling case not only calls for a different analysis but also raises a problem that we must solve.

The linearity test in the non-signaling setting has the following guarantee [CMS18]: if \( F \) is a \( k \)-non-signaling function such that \( \Pr_{x,y,F}[F(x) + F(y) = F(x+y)] \geq 1 - \varepsilon \) then \( F \) can be self-corrected (in the natural way) to a \((k/2)\)-non-signaling function \( \mathcal{F} \) that is \( 2^{O(k)}\varepsilon \)-close to a linear non-signaling function \( \mathcal{L} \). Note that self-correction is already part of the conclusion.

The above result appears sufficient for compiling a nsLPCP verifier into a corresponding nsPCP verifier. Namely, given a \( k \)-non-signaling function \( F : \{0,1\}^\ell \rightarrow \{0,1\} \), the nsPCP verifier checks that \( F(x) + F(y) = F(x+y) \) for random \( x, y \in \{0,1\}^\ell \) and also checks that the nsLPCP verifier accepts \( \mathcal{F} \). Analogously to before, if the nsLPCP verifier makes \( q \) queries and has soundness error \( \gamma \) against linear \((k/2)\)-non-signaling functions, then the resulting PCP verifier makes \( 3 + 2q \) queries and has soundness error \( \max\{1 - \varepsilon, \gamma + 2^{O(k)}\varepsilon\} \) against arbitrary \( k \)-non-signaling functions.

However, our analysis of the ALMSS verifier (the nsLPCP that we use) will require locality \( k = \Omega(\log N) \), which means that the additive term \( 2^{O(k)}\varepsilon \) grows with \( N \). This precludes achieving a constant soundness error with constant query complexity.
The foregoing motivates the problem of testing linearity of non-signaling functions with low error: how do we ensure that $\mathcal{F}$ is sufficiently close to a linear non-signaling function $L$? We stress that while in the classical case improving the “quality” of the self-correction has a straightforward solution (repeat the linearity test, and do self-correction), in the non-signaling case this problem is quite involved. Moreover, we do not wish to modify in any way the classical compiler, and thus relying on additional queries (even if only a constant number depending on $q$ and $\varepsilon$) is not an option.

We discuss our solution to this problem later on in Section 2.6, thereby providing the missing ingredient of our compiler from nsLPCP to nsPCP. In the meantime, in Sections 2.2 to 2.5, we discuss how we prove that the ALMSS verifier is secure against linear non-signaling functions.

### 2.2 The linear ALMSS verifier against linear non-signaling functions

Our goal is to establish that the linear PCP verifier of [ALMSS98] (the “ALMSS verifier”) is sound against linear non-signaling functions, and thus prove that every language $L \in \text{DSIZE}(S)$ has a constant-query nsLPCP verifier that is sound against linear $O(\log S)$-non-signaling functions. Note that we invoke the ALMSS verifier on deterministic (DSIZE) computations, rather than on nondeterministic (NSIZE) computations as in the classical case. We now recall the ALMSS verifier.

Let $L \in \text{DSIZE}(S)$ be a language, and let $\{C_i\}_{i \in \mathbb{N}}$ be a uniform boolean circuit family of size $N := S(n)$ that decides $L$ (for all $x \in \{0, 1\}^n$, $x \in L$ iff $C_n(x) = 1$). Hereafter we omit the subscript in $C_n$ as it is clear from context. Given an input $x$, one can express the condition “$C(x) = 1$” as a system of simple equations over $C$’s wires $W$; the variables are $w = (w_1, \ldots, w_N)$, one per wire. We use the convention that the input wires are $w_1, \ldots, w_n$ and the output wire is $w_N$. To ensure input consistency we need that $w_j = x_j$ for every $j \in \{1, \ldots, n\}$; to ensure correct gate computations we need that, for every $j \in \{n + 1, \ldots, N\}$, $w_j$ is the correct combination of the variables used to compute it (e.g., if $w_j$ is the output of an AND gate with inputs $w_{j_1}$ and $w_{j_2}$ then the equation is $w_j = w_{j_1} \cdot w_{j_2}$); to ensure that the output is 1 we need that $w_N = 1$. This can be summarized as a system of $M := N + 1$ equations $\{P_j(w) = c_j\}_{j \in [M]}$, where $P_1, \ldots, P_M$ are quadratic polynomials (each involving at most three variables in $w$) and $c_1, \ldots, c_M$ are boolean constants.

The ALMSS verifier is given below. We overload notation and use $P_j$ to also denote the upper triangular matrix in $\{0, 1\}^{N^2}$ such that $P_j(w) = \langle P_j, w \otimes w \rangle$; that is, if $P_j(w) = \sum_{i=1}^N a_i w_i + \sum_{1 \leq i' < i \leq N} a_{i,i'} w_i w_i'$, then $P_j$ has $a_i$ in the diagonal entry $(i, i)$ and $a_{i,i'}$ in the entry $(i, i')$, for $1 \leq i < i' \leq N$. Also, for $a \in \{0, 1\}^N$, $D_a$ is the diagonal matrix in $\{0, 1\}^{N^2}$ whose diagonal is $a$.

The ALMSS verifier, given input $x \in \{0, 1\}^n$ and oracle access to a linear non-signaling function $\mathcal{L} : \{0, 1\}^{N^2} \to \{0, 1\}$, works as follows:

1. Use the circuit $C$ and input $x$ to construct the matrices $P_1, \ldots, P_M \in \{0, 1\}^{N^2}$ and constants $c_1, \ldots, c_M \in \{0, 1\}$, which represent the computation of $C$ on $x$.
2. Draw random $s \in \{0, 1\}^M$, $u, v \in \{0, 1\}^N$, and query $\mathcal{L}$ on the set $\{\sum_{j=1}^M s_j P_j, D_u, D_v, u \otimes v\}$.
3. Check that $\mathcal{L}(\sum_{j=1}^M s_j P_j) = \sum_{j=1}^M s_j c_j$ and that $\mathcal{L}(D_u) \mathcal{L}(D_v) = \mathcal{L}(u \otimes v)$.

If $C(x) = 1$, the honest proof is the linear function $\pi : \{0, 1\}^{N^2} \to \{0, 1\}$ where $\pi(Z) := \langle w \otimes w, Z \rangle = \sum_{i,i' \in [N]} w_i w_i' \cdot Z_{i,i}$ where $w_i$ is now the value of the $i$-th wire in the computation of $C$ on $x$.

The challenge is to prove that the ALMSS verifier is sound against linear non-signaling functions. Namely, we must show that if there is a linear non-signaling function $L$ that is accepted with good probability then $x \in L$, or equivalently that $C(x) = 1$. We discuss this in the next sub-sections.
2.3 A linear local assignment generator suffices

The first step in our soundness analysis shows that, to establish that $C(x) = 1$, it suffices to construct a linear local assignment generator with sufficiently small error.

A linear $k$-local assignment generator for $(C, x)$ with error $\varepsilon$ is a linear $k$-non-signaling function $\mathcal{A} : \{0,1\}^N \to \{0,1\}$ that individually satisfies each of the $M$ constraints with probability $1 - \varepsilon$ (over the randomness of $\mathcal{A}$). Namely, (a) for each $i \in \{1, \ldots, n\}$, $\Pr[\mathcal{A}(e_i) = x_i] \geq 1 - \varepsilon$; (b) for each $i \in \{n+1, \ldots, N\}$, if $w_i$ is the output of a unary gate $g$ with input $w_j$ then $\Pr[\mathcal{A}(e_i) = g(\mathcal{A}(e_j))] \geq 1 - \varepsilon$, else if $w_i$ is the output of a binary gate $g$ with inputs $w_{j_1}, w_{j_2}$ then $\Pr[\mathcal{A}(e_i) = g(\mathcal{A}(e_{j_1}), \mathcal{A}(e_{j_2}))] \geq 1 - \varepsilon$; (c) $\Pr[\mathcal{A}(e_X) = 1] \geq 1 - \varepsilon$. (Here $e_i$ is the $i$-th vector in the standard basis.)

Lemma (informal). If there exists a $k$-local assignment generator for $(C, x)$ with error $\varepsilon$ for $k = \Omega(\log N)$ and $\varepsilon = O\left(\frac{1}{N \log N}\right)$, then $C(x) = 1$.

We sketch the proof of this lemma. The transcript of the computation of $C$ on $x$ is the unique correct assignment to all the wires. We say that a wire $w_i \in W$ of $C$ is correct whenever $\mathcal{A}(e_i)$ equals the value contained in this transcript; more generally, we say that a vector $z \in \{0,1\}^N$ is correct if $\mathcal{A}(z)$ equals the value of $z$ in the linear extension of the transcript. Below, we partition $C$’s wires $W$ into layers $W_1, \ldots, W_H$ according to depth. (We assume layered circuits; see Section 3.1.)

As a warmup, suppose for now that $k \geq N$. The probability that all wires in $W_1$ are correct is at least $1 - |W_1| \varepsilon$, and the probability that all the gates are correct is at least $1 - \sum_{h=2}^H |W_h| \varepsilon$. Therefore by union bound, the probability that all wires in the circuit are correct is at least $1 - \sum_{h=1}^H |W_h| \varepsilon$, because if all the input wires are correct and all the gates are computed correctly, then all the wires in the circuit are correct. In particular, we deduce that the output wire is correct with probability $1 - \sum_{h=1}^H |W_h| \varepsilon = 1 - |W| \varepsilon = 1 - N \varepsilon$. Since the output wire is 1 with probability $1 - \varepsilon$, and $\varepsilon = O\left(\frac{1}{N}\right)$, we conclude that $\Pr[C(x) = 1] > 0$, and thus $C(x) = 1$.

The above argument requires that $k \geq N$, because we have to simultaneously “view” assignments to all wires in the circuit. While the argument can be easily modified so that we only require $k$ to be at least twice the width of $C$, the latter may still be much larger than $O(\log N)$.

Using the linearity of $\mathcal{A}$, however, we can modify the argument to merely require $k = \Omega(\log N)$. For each layer $h$, we define an event $E_h$ such that if $E_h$ holds, then any wire in layer $h$ is correct with high probability. In the warmup above $E_h$ is the event “all wires in layer $h - 1$ are correct”; in our proof $E_h$ is the event “$t$ random linear combinations of wires in layer $h$ are correct”. Given a wire $w_i$ in layer $h$, we can bound the event “$\mathcal{A}(e_i)$ is incorrect and $E_h$ holds” as follows. If $\mathcal{A}(e_i)$ is incorrect, then all linear combinations of wires in layer $h$ can be split into pairs $z$ and $z + e_i$, and exactly one of $\mathcal{A}(z)$ and $\mathcal{A}(z + e_i)$ is incorrect. Hence, the probability that a random linear combination of wires in layer $h$ is correct, given that $\mathcal{A}(e_i)$ is incorrect, is at most $1/2$, and so $\Pr[E_h | \mathcal{A}(e_i) \text{ is incorrect}] \leq 2^{-t}$, since the $t$ random linear combinations are independent. Using Bayes’s rule (and an additional assumption that $\Pr[E_h] \geq 1/2$), we deduce that $\Pr[\mathcal{A}(e_i) \text{ is incorrect} | E_h]$ is small. We then proceed inductively on the layers as before.

The argument above requires that $\varepsilon = O\left(\frac{1}{N \log N}\right)$. One may wonder whether a similar result could be proved with, say, $\varepsilon = O(1)$. We additionally prove that our analysis is almost tight, in that an error of $\varepsilon = O\left(\frac{\log N}{N}\right)$ is necessary, regardless of how large the locality $k$ is.

See Section 7 for details.

Local assignment generators in prior works. Local assignment generators appear in prior works on nsPCPs [KRR14; PR17], but our notion is qualitatively different, as we now explain.
Prior works consider local assignment generators for an augmented circuit $C_{aug}$ rather than for $C$ itself. (See Section 1.5 for a discussion of concurrent work that avoids augmented circuits.) Informally, $C_{aug}$ not only contains $C$ as a sub-circuit but also low-degree extensions of $C$’s layers as well as subcircuits computing all low-degree tests on these. The wires contained in these additional subcircuits are what enables defining an event $E_h$ on which to condition for each layer.

The analogue of the augmented circuit $C_{aug}$ in our setting, however, has exponential size, and thus we cannot use it. Namely, we would have to encode each layer of $C$ via the Hadamard code (all linear combinations of wires in the layer) and then compute all possible linear tests on these.

Instead, our assumption that the local assignment generator is a linear non-signaling function implies that we do not have to construct an augmented circuit. Namely, the linear combinations that we use to define the event $E_h$ are implicitly available due this linearity, and so there is no need to augment $C$ (nor, in particular, to introduce any gates that evaluate linearity tests).

The assumption that the local assignment generator is linear is justified by the fact that a different part of our construction (the linearity test in our generic compiler) ensures the non-signaling function is (close to) linear. Overall, this separation not only avoids the aforementioned issues of using augmented circuits, but also simplifies the analysis of the local assignment generator.

### 2.4 Constructing the linear local assignment generator

Given a $k$-non-signaling function $L: \{0,1\}^{N^2} \to \{0,1\}$ that is accepted by the ALMSS verifier with probability at least $1 - \varepsilon$, we can obtain a linear $k$-local assignment generator $A: \{0,1\}^N \to \{0,1\}$ with error $O(\varepsilon)$ by “restricting $L$ to its diagonal”. Namely, in order to query $A$ at $v \in \{0,1\}^N$, we query $L$ at $D_v \in \{0,1\}^{N^2}$, where $D_v$ is the diagonal matrix that has $v$ as its diagonal.

We show that, since $L$ is accepted with probability at least $1 - \varepsilon$, $L$ must satisfy any individual constraint $P_j(w) = c_j$ with probability at least $1 - O(\varepsilon)$, and this directly implies that the linear local assignment generator $A$ has error $O(\varepsilon)$. (See Section 8.2 for details.)

The discussion so far already gives us a weak bound on the soundness error of the ALMSS verifier, namely $1 - O(\frac{1}{N\log N})$. Indeed, for $k = O(\log N)$ and $\varepsilon = O(\frac{1}{N\log N})$, we can apply the lemma above (in Section 2.3) to conclude that $C(x) = 1$.

However, our goal is to show that the ALMSS verifier (as is) has constant soundness error, and doing so requires more technical work, which we discuss next.

**Remark 2.1.** We stress that proving a soundness error of even $1 - O(\frac{1}{N\log N})$ is a non-trivial statement. This is in contrast to the classical setting, where if an assignment satisfies an $1 - \varepsilon$ fraction of the $M = N + 1$ constraints for $\varepsilon < 1/M$, then, trivially, all constraints are satisfied.

### 2.5 The ALMSS verifier has constant soundness error

Our goal is to prove that the ALMSS verifier has constant soundness error. In a first step (Section 2.5.1), we use the soundness error proved above (Section 2.4) to show that the $t$-repeated ALMSS verifier has soundness error $\gamma$ when $t = \Omega(\log N + \log \frac{1}{\varepsilon})$. In a second step (Section 2.5.2), we prove that the basic ALMSS verifier (no repetitions) has constant soundness error. The second step is generic and of independent interest: we prove that if a $t$-repeated verifier has soundness error $\exp(-t)$, then the corresponding basic verifier has soundness error $O(1)$.
2.5.1 The $t$-repeated ALMSS verifier has soundness error $\exp(-t)$

While in the classical setting reducing soundness error via simple repetition is straightforward ($t$-wise repetition reduces soundness error from $\delta$ to $\delta^t$), in the non-signaling setting simple repetition does not work. Indeed, consider the non-signaling function (in fact, distribution) that, with probability $1 - \varepsilon$, answers the verifier’s queries in an accepting way, and otherwise answers randomly. This non-signaling function is accepted by the $t$-repeated verifier with probability $\approx 1 - \varepsilon$, which is about the same as the probability that it is accepted by a single verifier.

However, this example provides intuition for how one circumvents this issue. Informally, we would like to extract the “$1 - \varepsilon$ good part” that satisfies the verifier, and drop the “$\varepsilon$ bad part”. We follow a technique used in [KRR14] and, instead of arguing about the probability that $L$ passes the $t$-repeated verifier, we argue that the non-signaling function $L$ conditioned on passing the $t$-repeated verifier passes the basic verifier with high probability. Indeed, in the aforementioned example, conditioning on at least one test passing removes the “$\varepsilon$ bad part” injected by the distribution, and intuitively extracts the part of $L$ that is passing the verifier. An interesting feature of our analysis of the verifier is that our conclusion is about the basic verifier, not the relaxed $t$-repeated verifier, which plays a major role in the analysis in [KRR14,7]. This is a qualitative difference in our analysis arising from our use of property testing (not present in [KRR14]), which also simplifies the analysis.

In more detail, let $L'$ denote the linear non-signaling function that equals $L$ when conditioned on passing the $t$-repeated verifier. Namely, if $E$ is the (random) event that $L$ passes the $t$-repeated verifier, then for any $S \subseteq \{0, 1\}^n$ (of some maximal size) and $\vec{b} \in \{0, 1\}^S$, we define

$$
\Pr[L'(S) = \vec{b}] := \Pr[L(S) = \vec{b} \mid E] = \frac{\Pr[L(S) = \vec{b} \wedge E]}{\Pr[E]}. \tag{1}
$$

We then prove that $L'$ passes the basic verifier with probability at least $1 - \frac{1/\Pr[E]}{\exp(t)}$.

The proof uses a generic lemma (Lemma 5.1) stating that, if we run $t + d$ independent tests, then the probability that at most $r$ out of the first $d$ tests pass and all of the last $t$ tests pass is at most $(\frac{d}{t+d})^r + 1$. A naive application of this lemma (with $r = 0$ and $d = 1$) shows that $L'$ passes the basic verifier with probability at least $1 - \frac{1/\Pr[E]}{t+1}$. This is not enough, because (using $\Pr[E] \geq \gamma$) we would require $t = \Omega(N \log N \cdot \frac{1}{\gamma})$ to prove soundness, which is again far too many repetitions.

However, we leverage the linearity of $L$ to deduce the stronger guarantee, as we now explain. We want to bound the probability that $L'$ does not pass the basic verifier, which means we need to bound the probability that $L$ fails exactly the first test of $t + 1$ independent tests. We do this by arguing this individually for each of the two types of tests made by the ALMSS verifier: the tensor test “$L(D_u)L(D_v) = L(u \otimes v)$” and the satisfiability test “$L(\sum_{j=1}^{M} s_j P_j) = \sum_{j=1}^{M} s_j c_j$”. We will explain our techniques in the case of the satisfiability test; the same techniques work for the tensor test, but the algebra is messier.

In the case of the satisfiability test, we split the “special” test (i.e., the first one) into $d$ pairs of tests, such that each individual test is random, but each pair is correlated so that if both tests in some pair pass, then the original test passes. Specifically, we draw $d$ random vectors $s^{(1)}, \ldots, s^{(d)} \in \{0, 1\}^M$, and then we split the test “$L(\sum_{j=1}^{M} s_j P_j) = \sum_{j=1}^{M} s_j c_j$” into the $d$ pairs of

---

6 Even if simple repetition were to reduce soundness error from $\delta$ to $\delta^t$, then to get $\delta^t = \gamma$ when $\delta = 1 - O(\frac{1}{N \log N})$ we would need to repeat $t = \Omega(N \log N + \log \frac{1}{\gamma})$ times, which requires too large of a locality $k$ for the analysis.

7 The relaxed $t$-repeated verifier runs $t$ tests and accepts if a large fraction of them pass.
tests

\[ \mathcal{L} \left( \sum_{j=1}^{M} (s_j + s_j^{(i)} P_j) \right) = \sum_{j=1}^{M} (s_j + s_j^{(i)} c_j) \quad \text{and} \quad \mathcal{L} \left( \sum_{j=1}^{M} s_j^{(i)} P_j \right) = \sum_{j=1}^{M} s_j^{(i)} c_j. \]

This allows us to apply the lemma with \( d = O(t) \), and \( r = O(t) \), which shows that \( \mathcal{L}' \) passes the basic verifier with probability at least \( 1 - \frac{1}{\operatorname{Pr}[E]} \exp(t) \), an exponential decrease in \( t \).

The above analysis shows soundness error of \( \gamma \) for the \( t \)-repeated verifier, for \( t = O(\log N + \log \frac{1}{\gamma}) \).

Indeed, by the above argument, the conditioned function \( \mathcal{L}' \) passes the basic verifier with probability \( 1 - \frac{1}{\gamma} \exp(-t) = 1 - O\left(\frac{1}{N \log N}\right) \), by choice of \( t \). The analysis in the previous section (Section 2.4) then implies that \( C(x) = 1 \), proving soundness of the \( t \)-repeated verifier.

Setting \( \gamma = \exp(-t) \), the discussion so far merely shows that the \( t \)-wise repetition of the ALMSS verifier, which makes \( 4t \) queries, has soundness error \( \exp(-t) \) when \( t = \Omega(\log N) \); moreover, we get no conclusions for \( t = o(\log N) \). But we still did not conclude anything about the soundness of a single invocation of the 4-query ALMSS verifier. We next discuss how to handle this case.

### 2.5.2 Back to the 4-query ALMSS verifier

We establish that the ALMSS verifier has constant soundness error by proving a generic lemma. The lemma states that, for \( \mathcal{V} \) the basic PCP verifier, if the \( t \)-repeated verifier \( \mathcal{V}' \) has soundness error \( \exp(-t) \), then \( \mathcal{V} \) has soundness error \( O(1) \). Since we have already argued that the \( t \)-repeated ALMSS verifier has soundness error \( \exp(-t) \) for \( t = \Omega(\log N) \), we can conclude that the basic ALMSS verifier has soundness error \( O(1) \), against \( O(\log N) \)-non-signaling linear functions.

In the classical case, the proof of this generic fact is trivial: a (classical) function passes a PCP verifier \( \mathcal{V} \) with probability \( \delta \) if and only if it passes the \( t \)-repeated verifier \( \mathcal{V}' \) with probability \( \delta^t \). However, in the non-signaling case, it is not clear what one can say because a non-signaling function can provide correlated answers across repetitions. Nevertheless, we are able to lower bound the probability that \( \mathcal{V}' \) accepts by a quantity that is almost \( \delta^t \) (which is, in particular, almost tight).

To our knowledge, we are the first to relate the soundness of \( \mathcal{V} \) to the soundness of \( \mathcal{V}' \). Generic statements in prior works (starting with [KRR13]) have related the soundness of \( \mathcal{V}' \) to the soundness of the \( t \)-repeated relaxed verifier (which accepts if a vast majority of the \( t \) tests pass), but did not provide conclusions about the basic verifier \( \mathcal{V} \).

### 2.6 Testing linearity with low error

Below we discuss linearity testing with low error (Theorem 3) in more detail.

**Warmup: distributions.** We have discussed (in Section 2.1) how to test linearity with low error in the classical setting. In order to illustrate some of the difficulties that arise in the non-signaling setting, we first discuss a special case of it: testing linearity against a distribution over functions.

First, suppose that \( \mathcal{D} \) is a distribution over functions \( f: \{0,1\}^n \rightarrow \{0,1\} \) that passes the linearity test with probability \( 1 - \varepsilon \). The self-correction \( \mathcal{D}' \) that on input \( x \in \{0,1\}^n \) samples a random \( z \in \{0,1\}^n \) and outputs \( \mathcal{D}'(x) = \mathcal{D}(z + x) - \mathcal{D}(z) \) is \( 2\varepsilon \)-close to a distribution over linear functions \( \mathcal{D}'^* \), namely, for every \( x \in \{0,1\}^n \) it holds that \( |\operatorname{Pr}[\mathcal{D}'(x) = 1] - \operatorname{Pr}[\mathcal{D}'^*(x) = 1]| \leq 2\varepsilon \). Indeed, consider the distribution \( \mathcal{D}'^* \) that samples \( f \leftarrow \mathcal{D} \) and outputs any linear function \( f^* \) closest to
Then, for every function \( f \) and \( x \in \{0, 1\}^n \), the probability over a random \( z \in \{0, 1\}^n \) that \( f^*(z) = f(z) \) and \( f^*(z + x) = f(z + x) \) is at least \( 1 - 2\varepsilon_f \), where \( \varepsilon_f := 1 - \Pr_{x,y}[f(x) + f(y) = f(x + y)] \). Denoting by \( d_f \) denotes the probability that \( D \) samples the function \( f \), we conclude that 
\[
\left| \Pr[D^*(x) = 1] - \Pr[D(x) = 1] \right| \leq \sum_f 2\varepsilon_f \cdot d_f = 2\varepsilon.
\]

Next, suppose that we seek a self-correction of \( D \) that is \( \delta \)-close to a distribution over linear functions, for \( \delta \ll 2\varepsilon \). One idea is to follow the same strategy as in the case of a single function: repeat the linearity test and then do self-correction. This idea, however, does not work now.

Consider the distribution \( D = (1 - \varepsilon) \cdot 0 + \varepsilon \cdot 1 \), i.e., the distribution that with probability \( 1 - \varepsilon \) answers according to the all-zeros function (a linear function), and with probability \( \varepsilon \) according to the all-ones function (a function maximally far from linear functions). While \( D \) passes the linearity test with probability \( 1 - \varepsilon \), \( D \) also passes the \( t \)-repeated linearity test with probability \( 1 - \varepsilon \). In other words, if \( D \) passes the \( t \)-repeated linearity test with probability \( 1 - \varepsilon \), we can still only conclude that \( D \) is \( 2\varepsilon \)-close to linear, independent of \( t \).

While repeating the test does not increase the rejection probability, it can still be used to improve the quality of self-correction, by considering a different notion of self-correction that penalizes functions in the support of \( D \) that are far from linear. Concretely, consider the distribution \( D_t \) that equals \( D \) when conditioned on the event that the \( t \)-repeated linearity test passes, and then define \( \hat{D}_t \) to be the self-correction of \( D_t \). That is, \( \hat{D}_t \) samples \( f \) from \( D_t \) and answers any query \( x \in \{0, 1\}^n \) by sampling \( z \in \{0, 1\}^n \) and returning \( f(z + x) - f(z) \). We claim that \( \hat{D}_t \) is very close to linear.

Indeed, suppose that \( D_t \) passes the \( t \)-repeated test with probability \( \gamma > 0 \), and let \( c > 1 \) be a parameter. A function \( f \) sampled from \( D_t \) is \( \frac{\ln \varepsilon}{t} \)-close to linear with probability at least 
\[
\frac{\gamma - 1/c}{\gamma} = 1 - \frac{\ln \varepsilon}{c \gamma t}.
\]
Setting \( c := t/\log t \), the probability that \( \hat{D}_t \) outputs a function \( f \) that is \( \frac{\ln \varepsilon}{t} \)-far from linear is at most \( \frac{\log t}{\gamma t} \). Therefore, by applying the argument from the beginning of this subsection, we conclude that \( \hat{D}_t \) is \( O(\frac{\log t}{t}) \)-close to a distribution over linear functions.

We can further reduce the distance to be exponentially small in \( t \) by performing self-correction \( t \) times: \( \hat{D}_t(x) \) now samples \( z_1, \ldots, z_t \in \{0, 1\}^n \), and outputs the majority of \( \{D(z_i + x) - D(z_i)\}_{i \in [t]} \) conditioned on the event that the \( t \)-repeated linearity test passes. By setting \( c := 2^{t/10} \) in the discussion above, we conclude that if we sample \( f \) from \( D_t \), then \( f \) is \( 0.1 \)-close to a linear function \( f^* \) with probability \( 1 - \frac{\ln \varepsilon}{2^{t/10}t} \). In particular, for every \( x \in \{0, 1\}^n \) it holds that \( \Pr_{z_i}[f^*(z_i) = f(z_i)] \wedge f^*(z_i + x) = f(z_i + x)] \geq 0.8 \), and so the probability that the majority value of \( \{D(z_i + x) - D(z_i)\}_{i \in [t]} \) is not equal to \( f^*(x) \) is \( 2^{-O(t)} \). In sum, the \( t \)-repeated self-correction conditioned on the event that the \( t \)-repeated linearity test passes yields us a distribution that is \( \frac{1}{t} 2^{t}^{-O(t)} \)-close to linear.

The non-signaling case. The case of non-signaling strategies is similar to the case of distributions in that the analysis of the self-correction involves conditioning over a certain event. Yet, the conclusions and steps of the proof are quite different. Informally, this is because non-signaling functions are quasi-distributions (probabilities can be negative), which prevents us from doing a straightforward analysis such as the one above. We now discuss how we address this.

Suppose that we have a \( k \)-non-signaling function \( \mathcal{F} : \{0, 1\}^n \rightarrow \{0, 1\} \) that passes the linearity test with probability \( 1 - \varepsilon \). The result of [CMS18] proves that the self-correction \( \hat{\mathcal{F}} \) defined as \( \hat{\mathcal{F}}(x) := \mathcal{F}(z + x) - \mathcal{F}(z) \) (where \( z \) is chosen randomly from \( \{0, 1\}^n \)) is \( 2^{O(k)} \varepsilon \)-close to linear. This is too large in our setting as we have \( k = O(\log N) \), and we would like the distance to be \( O(\varepsilon) \).

---

\(^8\)Recall that if \( f \) is \( 0.25 \)-close to linear functions then \( f^* \) is unique. We do not rely on uniqueness.

\(^9\)The \( t \)-repeated linearity test accepts a function \( f \) that is \( \frac{\ln \varepsilon}{t} \)-far from linear with probability at most \( (1 - \frac{\ln \varepsilon}{t})^t \leq \frac{1}{t} \).
Instead, we prove a slightly different guarantee from \cite{CMS18}. Namely, we show that there is a linear non-signaling function $\mathcal{L}$, such that on every set $S$, $\mathcal{F}$ is $O(|S|/\sqrt{\varepsilon})$-close to $\mathcal{L}$. Unlike in the result of \cite{CMS18}, our distance now decays with $|S|$, and is in particular independent of $k$. This is sufficient for our purposes, since we set $|S| = 4$, the number of queries made by the ALMSS verifier.

In our proof, we consider a different self-correction $\hat{\mathcal{F}}$ that, unlike $\hat{\mathcal{F}}$, is only used in the analysis and is not used by the compiler. First, we show that $\mathcal{F}_t$ passes the linearity test with probability $1 - \exp(-t)$, and that the result of \cite{CMS18} implies that $\mathcal{F}_t$ is very close to a linear non-signaling function. Then, we relate $\hat{\mathcal{F}}_t$ and $\hat{\mathcal{F}}$ to show that $\hat{\mathcal{F}}$ is $O(\sqrt{\varepsilon})$-close to a linear non-signaling function.

Informally, the self-correction $\mathcal{F}_t$ equals $\mathcal{F}$ with the standard self-correction procedure repeated $t$ times, conditioned on $\mathcal{F}$ passing $(1 - \sqrt{\varepsilon})t$ of $t$ repetitions of the linearity test. In more detail, given a subset $S \subseteq \{0, 1\}^n$, $\mathcal{F}_t(S)$ is the following distribution. For each $x \in S$, sample uniform and independent $z_x^{(1)}, \ldots, z_x^{(t)} \in \{0, 1\}^n$ conditioned on satisfying the same linear dependencies as in $S$; for instance, if $S = \{x, y, x + y\}$, then $z_x^{(i)} + z_y^{(i)} = z_{x+y}^{(i)}$ holds for all $i$. Then $\mathcal{F}_t$ assigns to each $x \in S$ the value $\max_{y \in \{0, 1\}^n} \{\mathcal{F}(z_x^{(i)} + x) - \mathcal{F}(z_x^{(i)})\}$ conditioned on the event that $\mathcal{F}$ passes at least $(1 - \sqrt{\varepsilon})t$ of $t$ repetitions of the basic linearity test. We note that if $\mathcal{F}$ is linear, then $\mathcal{F}_t \equiv \hat{\mathcal{F}} \equiv \mathcal{F}$.

The first part of the analysis uses Lemma 5.2 which informally states that by conditioning on $\mathcal{F}$ passing most of the $t$-repeated linearity tests, we force the conditioned $\mathcal{F}$ to behave “close” to linear. Specifically, letting $b_x = \mathcal{F}(z_x^{(i)} + x) - \mathcal{F}(z_x^{(i)})$, we get that with probability $1 - \exp(-t)$ there is a bit $b_x$ that equals $b_x^{(i)}$ for at least $\frac{3t}{4}$ of the $i$’s (so the majority is a vast majority), which implies that $\mathcal{F}_t(x) = b_x$, and analogously for $y$ and $x + y$. Then, via a similar argument, we show that with probability $1 - \exp(-t)$ for at least $\frac{3t}{4}$ of the $i$’s it holds that $y_x^{(i)} + b_y^{(i)} = b_y^{(i)}$. By union bound, these events hold simultaneously, and so we conclude that $\mathcal{F}_t$ satisfies $\mathcal{F}_t(x) + \mathcal{F}_t(y) = \mathcal{F}_t(x + y)$ with probability $1 - \exp(-t)$. We then invoke the result of \cite{CMS18} and conclude that $\hat{\mathcal{F}}$ is very close to some linear non-signaling function $\mathcal{L}$.

In the second step, we relate $\hat{\mathcal{F}}$ to $\mathcal{F}_t$ by claiming that if $\Pr[\mathcal{F}(x) + \mathcal{F}(y) = \mathcal{F}(x + y)] \geq 1 - \varepsilon$, then $\hat{\mathcal{F}}_t$ and $\mathcal{F}_t$ are close in some precise sense (see Section 9.3 for details). We first observe that if we run the $t$-repeate linearity test, i.e., choose $x^{(1)}, y^{(1)}, \ldots, x^{(t)}, y^{(t)}$ and check that $\mathcal{F}(x^{(i)}) + \mathcal{F}(y^{(i)}) = \mathcal{F}(x^{(i)} + y^{(i)})$ for every $i$, then a simple Markov argument (Lemma 5.5) shows that with high probability, most of the linearity tests are satisfied. For instance, with probability $1 - \sqrt{\varepsilon}$ at least $(1 - \sqrt{\varepsilon})t$ of the $i$’s satisfy the linear constraint. This means that the event conditioned on in the definition of $\mathcal{F}_t$ is a large event. We also know from the first part of the analysis that, with high probability, the conditioning causes most of the evaluations of $\mathcal{F}(z_x^{(i)} + x) - \mathcal{F}(z_x^{(i)})$ to output the same value. Intuitively, this implies that $\hat{\mathcal{F}}$ is close to $\mathcal{F}_t$, via the following reasoning. Since $\mathcal{F}_t$ conditions on a large event, it is close to the corresponding self-correction that does not condition at all. Since the majority taken over the evaluations of $\mathcal{F}(z_x^{(i)} + x) - \mathcal{F}(z_x^{(i)})$ when computing $\mathcal{F}_t$ is a vast majority, with high probability $\hat{\mathcal{F}}$ (which is a sample from one of the elements the majority is over) will agree with the vast majority. This allows us to conclude that for any set $S$, $\hat{\mathcal{F}}_t$ will be $O(|S|/\sqrt{\varepsilon})$-close to $\mathcal{F}_t$.

See Section 9 for details.
3 Definitions

We introduce the main definitions used throughout the paper.

3.1 Boolean circuits and the class DSIZE(S)

All boolean circuits in this paper are layered and have maximum fan-in 2. In more detail, a boolean circuit is a layered directed acyclic graph, in which the vertices are called “gates” and the edges are called “wires”. Furthermore, each gate is labeled as either AND, OR or NOT, and it holds that (1) every wire goes from a gate in layer \( h \) to a gate in layer \( h + 1 \), (2) each gate has at most 2 incoming wires, (3) each gate with 2 incoming wires is either an AND gate or an OR gate, and (4) each gate with 1 incoming wire is a NOT gate.

The class DSIZE(S) is the main complexity class that we study in the paper.

**Definition 3.1.** Given a computable function \( S : \mathbb{N} \rightarrow \mathbb{N} \), a language \( L \) is in DSIZE(S) if there exists a uniform family of boolean circuits \( C_n : \{0,1\}^n \rightarrow \{0,1\} \) such that \( C_n \) has at most \( S(n) \) gates (and is constructible in time \( O(S(n)) \)) and, for all \( x \in \{0,1\}^n \), \( x \in L \) if and only if \( C_n(x) = 1 \).

The reader may find it helpful to compare DSIZE(S) with NSIZE(S), where the second condition is replaced by “\( x \in L \) if and only if there exists \( w \) such that \( C_n(x,w) = 1 \)”.

3.2 Non-signaling functions

We define non-signaling functions and introduce useful notation for them. The text below is taken almost verbatim from [CMS18].

**Definition 3.2.** A \( k \)-non-signaling function \( \mathcal{F} : D \rightarrow \{0,1\} \) is a collection \( \mathcal{F} = \{\mathcal{F}_S\}_{S \subseteq D,|S| \leq k} \) where (i) each \( \mathcal{F}_S \) is a distribution over functions \( f : S \rightarrow \{0,1\} \), and (ii) for every two subsets \( S \) and \( T \) each of size at most \( k \), the restrictions of \( \mathcal{F}_S \) and \( \mathcal{F}_T \) to \( S \cap T \) are equal as distributions. (If \( S = \emptyset \) then \( \mathcal{F}_S \) always outputs the empty string.)

Note that any function \( f : D \rightarrow \{0,1\} \) induces a \( |D| \)-non-signaling function by setting \( \mathcal{F}_S \) to be the distribution that outputs \( f|_S \) with probability 1. More generally, any distribution \( \mathcal{D} \) over functions \( f : D \rightarrow \{0,1\} \) induces a corresponding \( |D| \)-non-signaling function by defining \( \mathcal{F}_S \) to be the distribution that samples \( f \leftarrow \mathcal{D} \) and outputs \( f|_S \).

Given a set \( S \subseteq D \) of size \( |S| \leq k \) and a string \( \vec{b} \in \{0,1\}^S \), we define

\[
\Pr\left[ \mathcal{F}(S) = \vec{b} \right] := \Pr_{f \leftarrow \mathcal{F}_S}\left[ f(S) = \vec{b} \right].
\]

The non-signaling property in this notation is the following: for every two subsets \( S, T \subseteq D \) of sizes \( |S|, |T| \leq k \) and every string \( \vec{b} \in \{0,1\}^{S \cap T} \), \( \Pr[\mathcal{F}(S)|_{S \cap T} = \vec{b}] = \Pr[\mathcal{F}(T)|_{S \cap T} = \vec{b}] \), where the probability is over the randomness of \( \mathcal{F} \).

Sometimes it is more convenient to consider a vector of inputs (rather than a set of inputs), and so we define notation for this case. Given a vector \( (x_1, \ldots, x_s) \) with entries in \( D \) and a vector \( (b_1, \ldots, b_s) \) with entries in \( \{0,1\} \) (with \( s \in \{1, \ldots, k\} \)), we define \( \Pr[\mathcal{F}((x_1, \ldots, x_s)) = (b_1, \ldots, b_s)] \) and \( \Pr[\mathcal{F}(x_1) = b_1, \ldots, \mathcal{F}(x_s) = b_s] \) to be the probability

\[
\Pr_{f \leftarrow \mathcal{F}(x_1, \ldots, x_s)}[f(x_1) = b_1, \ldots, f(x_s) = b_s].
\]
Note that \( \{x_1, \ldots, x_s\} \) is an unordered set and its size may be less than \( s \), because the entries of the vector \((x_1, \ldots, x_s)\) may not be distinct. We abuse notation and still use symbols such as \( S \) and \( \vec{b} \) to denote vectors as above. We stress that we use an ordering on \( S \) merely to match each element of \( S \) to the corresponding element in \( \vec{b} \); the event remains unchanged if one permutes the entries of \( S \) and \( \vec{b} \) according to the same permutation.

### 3.3 Quasi-distributions

A quasi-distribution extends the notion of a probability distribution by allowing negative probabilities. Quasi-distributions are equivalent to non-signaling functions [CMS18], and are a useful to view, and analyze, non-signaling functions. The text below is taken almost verbatim from [CMS18].

**Definition 3.3 (quasi-distributions).** Let \( D \) be a finite domain, and denote by \( U_D \) the set of all boolean functions of the form \( f : D \to \{0,1\} \). A quasi-distribution \( Q \) over a subset \( G \subseteq U_D \) is a set of real numbers \( \{q_f\}_{f \in U_D} \) such that \( \sum_{f \in U_D} q_f = 1 \) and \( q_f = 0 \) for every \( f \not\in G \).

**Definition 3.4 (quasi-probability).** Given a quasi-distribution \( Q = \{q_f\}_{f \in U_D}, \) a subset \( S \subseteq D, \) and a string \( \vec{b} \in \{0,1\}^S \), we define the quasi-probability of the event “\( Q(S) = \vec{b} \)” to be the following (possibly negative) real number

\[
\tilde{\Pr}[Q(S) = \vec{b}] := \sum_{f \in U_D \text{ s.t. } f(S) = \vec{b}} q_f .
\]

As in the case of non-signaling functions, it is sometimes more convenient to consider a vector of inputs rather than a set. Given a vector \((x_1, \ldots, x_s)\) with entries in \( D \) and a vector \((b_1, \ldots, b_s)\) with entries in \( \{0,1\} \), we define \( \Pr[Q((x_1, \ldots, x_s)) = (b_1, \ldots, b_s)] \) and \( \Pr[Q(x_1) = b_1, \ldots, Q(x_s) = b_s] \) to be the (possibly negative) real number \( \sum_{f \in U_D \text{ s.t. } \forall i, f(x_i) = b_i} q_f \). We abuse notation and still use symbols such as \( S \) and \( \vec{b} \) to denote vectors as above.

**Definition 3.5 (locality).** Let \( D \) be a finite domain of size \( N \) and let \( k \in \{1, \ldots, N\} \). A quasi-distribution \( Q \) over \( U_D \) is \( k \)-local if for every subset \( S \subseteq D \) of size \( |S| \leq k \) and string \( \vec{b} \in \{0,1\}^S \),

\[
\tilde{\Pr}[Q(S) = \vec{b}] \in [0,1] .
\]

For completeness, we also say that all quasi-distributions are 0-local.

If \( Q \) is \( k \)-local, then for every subset \( S \subseteq D \) of size \( |S| \leq k \), we may view \( Q(S) \) as a probability distribution over \( \{0,1\}^S \). If \( Q \) is \( k \)-local then it is \( s \)-local for every \( s \in \{0,1, \ldots, k\} \).

For \( Q \) to be \( k \)-local, it suffices for all relevant \( \tilde{\Pr}[Q(S) = \vec{b}] \) to be non-negative (as opposed to be in \([0,1]\)). This is because \( \sum_f q_f = 1 \), so that \( \sum_{\vec{b} \in \{0,1\}^S} \tilde{\Pr}[Q(S) = \vec{b}] = 1 \) and, if all terms in this sum are non-negative, then we can deduce that \( \tilde{\Pr}[Q(S) = \vec{b}] \leq 1 \) for every \( \vec{b} \).

**Definition 3.6 (statistical distance).** Given a finite domain \( D \) and an integer \( k \in \{1, \ldots, |D|\} \), the \( \Delta_k \)-distance between two quasi-distributions \( Q \) and \( Q' \) is

\[
\Delta_k(Q, Q') := \max_{S \subseteq D, |S| \leq k} \Delta(Q_S, Q'_S) ,
\]

where \( \Delta(Q_S, Q'_S) := \max_{E \subseteq \{0,1\}^S} \left| \tilde{\Pr}[Q(S) \in E] - \tilde{\Pr}[Q'(S) \in E] \right| . \)

We say that \( Q \) and \( Q' \) are \( \varepsilon \)-close in the \( \Delta_k \)-distance if \( \Delta_k(Q, Q') \leq \varepsilon \); else, they are \( \varepsilon \)-far.
**Remark 3.7** (distance for non-signaling functions). The definition of $\Delta_k$-distance naturally extends to defining distances between $k'$-non-signaling functions, as well as between quasi-distributions and $k'$-non-signaling functions, provided that $k \leq k'$.

The notion above generalizes the standard notion of statistical (total variation) distance: if $Q$ and $Q'$ are distributions then their $\Delta_{|D|}$-distance equals their statistical distance. Also note that if $Q$ and $Q'$ are $k$-local quasi-distributions then their $\Delta_k$-distance equals the maximum statistical distance, across all subsets $S \subseteq D$ with $|S| \leq k$, between the two distributions $Q_S$ and $Q'_S$ — in particular this means that any experiment that queries exactly one set of size at most $k$ cannot distinguish between the two quasi-distributions with probability greater than $\Delta_k(Q, Q')$.

We conclude by stating a theorem from [CMS18], which motivates the use of quasi-distributions. Informally, the theorem states that quasi-distributions and non-signaling functions are equivalent.

**Theorem 4** ([CMS18]). For any finite domain $D$, the following hold:

1. For every $k$-local quasi-distribution $Q$ over functions $f : D \rightarrow \{0,1\}$ there exists a $k$-non-signaling function $F : D \rightarrow \{0,1\}$ such that $\Delta_k(Q, F) = 0$.

2. For every $k$-non-signaling function $F : D \rightarrow \{0,1\}$ there exists a $k$-local quasi-distribution $Q$ over functions $f : D \rightarrow \{0,1\}$ such that $\Delta_k(Q, F) = 0$. The set of such quasi-distributions forms an affine subspace of $R^{2^{|D|}}$.

### 3.4 Input-oblivious queries

An oracle algorithm $A : \{0,1\}^* \rightarrow \{0,1\}$ has input-oblivious queries if, given an input $x \in \{0,1\}^n$, the queries to its oracle are determined non-adaptively based solely on $n$ and on its internal randomness, and only the final decision depends on $x$. In more detail, one can view $A$ as specified by a pair of algorithms $(Q, D)$ where: (a) $Q$ is a probabilistic algorithm known as the query sampler that, given $n$, outputs a set of queries for the oracle; (b) $D$ is a deterministic algorithm known as the decision predicate that, given $x$ (and $Q$’s randomness), outputs a decision bit. All verifiers that we consider in this paper have input-oblivious queries.

### 3.5 Non-signaling PCPs

A non-signaling PCP (nsPCP) [KRR14] is a PCP in which soundness is further required to hold against any (sufficiently local) non-signaling function.

In this setting, for a query alphabet $D$ and answer alphabet $\Sigma$, a nsPCP verifier is an algorithm $V$ that, given oracle access to a non-signaling function $F : D \rightarrow \Sigma$, queries $F$ at a single subset $S \subseteq D$ (on which $F$ is defined), receives answers $a \in \Sigma^S$, and outputs a decision bit.

**Definition 3.8.** Given $\epsilon \in [0,1]$ and $k \in \mathbb{N}$, a nsPCP verifier $V$ for a language $L$ has soundness error $\epsilon$ against $k$-non-signaling functions if the following holds for every $x \in \{0,1\}^*$:

1. If $x \in L$ then there exists a function $\pi : D \rightarrow \Sigma$ such that $\Pr[V^\pi(x) = 1] = 1$.
2. If $x \notin L$ then, for every $k$-non-signaling function $F : D \rightarrow \Sigma$, $\Pr[V^F(x) = 1] \leq \epsilon$.

**Remark 3.9.** In Item 2, if we consider only all functions $\pi : D \rightarrow \Sigma$, rather than all $k$-non-signaling functions $F : D \rightarrow \Sigma$, then we recover the usual PCP definition.
If the query alphabet $D$ has size at most $\ell$, the PCP verifier uses at most $r$ random bits to sample a subset $S$ of size at most $q$ in time $T_1$, and the PCP verifier decides to accept/reject in time $T_2$ based on $F$’s answers to $S$, then we write

\[
L \in \text{nsPCP} \begin{bmatrix}
\text{soundness error: } & \varepsilon \\
\text{randomness: } & r \\
\text{proof length: } & \ell \\
\text{query complexity: } & q \\
\text{locality: } & k \\
\text{query sampler time: } & T_1 \\
\text{decision predicate time: } & T_2
\end{bmatrix}.
\]

Note that we must require $q \leq k$, since querying $F$ on a set of size greater than $k$ is undefined.

### 3.6 Linear non-signaling PCPs

A linear non-signaling PCP (nsLPCP) is a PCP that is sound against (sufficiently local) non-signaling functions that are linear, that is, they are induced by quasi-distributions over linear functions. This definition is motivated by the equivalence between non-signaling functions and quasi-distributions, as well as the linearity testing result, proved in [CMS18].

**Definition 3.10.** A $k$-non-signaling function $L: \{0,1\}^\ell \to \{0,1\}$ is linear if there exists a $k$-local quasi-distribution $Q$ over linear functions $f: \{0,1\}^\ell \to \{0,1\}$ such that $\Delta_k(L,Q) = 0$.

**Definition 3.11.** Given $\varepsilon \in [0,1]$ and $k \in \mathbb{N}$, a nsLPCP verifier $V$ for a language $L$ has soundness error $\varepsilon$ against linear $k$-non-signaling functions if the following holds for every $x \in \{0,1\}^n$:

1. If $x \in L$ then there exists a linear function $\pi: \{0,1\}^\ell \to \{0,1\}$ such that $\Pr[V^\pi(x) = 1] = 1$.
2. If $x \notin L$ then, for every linear $k$-non-signaling function $L: \{0,1\}^\ell \to \{0,1\}$, $\Pr[V^L(x) = 1] \leq \varepsilon$.

**Remark 3.12.** In Item 2, if we only consider all linear functions $\pi: \{0,1\}^\ell \to \{0,1\}$, rather than all linear $k$-non-signaling functions $L: \{0,1\}^\ell \to \{0,1\}$, then we recover the usual LPCP definition.

If the LPCP verifier uses at most $r$ random bits to sample a subset $S$ of size at most $q$ in time $T_1$, queries $L: \{0,1\}^\ell \to \{0,1\}$ on $S$, and decides to accept/reject in time $T_2$ based on $L$’s answers to $S$, then we write

\[
L \in \text{nsLPCP} \begin{bmatrix}
\text{soundness error: } & \varepsilon \\
\text{randomness: } & r \\
\text{proof length: } & \ell \\
\text{query complexity: } & q \\
\text{locality: } & k \\
\text{query sampler time: } & T_1 \\
\text{decision predicate time: } & T_2
\end{bmatrix}.
\]

Like before, we must require $q \leq k$, since querying $L$ on a set of size greater than $k$ is undefined.
4 Formal statements of our results

Our main theorem (discussed in Section 1.2) is that the exponential-length constant-query PCP in [ALMSS98] (with no modifications) is sound against non-signaling functions. We obtain:

**Theorem 5** (formal statement of Theorem 1).

\[
\text{DSIZE}(S) \subseteq \text{nsPCP} \quad \left[ \begin{array}{|c|}
\hline
\text{soundness error:} & 1 - 1/10^7 \\
\text{randomness:} & O(S^2) \\
\text{proof length:} & 2O(S^2) \\
\text{query complexity:} & 11 \\
\text{locality:} & O(\log^2 S) \\
\text{query sampler time:} & O(S^2) \\
\text{decision predicate time:} & O(n) \\
\hline
\end{array} \right].
\]

Our proof relies on several results (discussed in Section 1.3), which we formally state below. At the end of this section we show how to combine these results to prove Theorem 5.

First, we prove that the linear PCP in [ALMSS98] is sound against linear non-signaling functions.

**Theorem 6** (formal statement of Theorem 2; proved in Section 8). The linear PCP in [ALMSS98] is sound against linear non-signaling functions, giving us the following class inclusion:

\[
\text{DSIZE}(S) \subseteq \text{nsLPCP} \quad \left[ \begin{array}{|c|}
\hline
\text{soundness error:} & 39/40 \\
\text{randomness:} & O(S) \\
\text{proof length:} & O(S^2) \\
\text{query complexity:} & 4 \\
\text{locality:} & O(\log S) \\
\text{query sampler time:} & O(S^2) \\
\text{decision predicate time:} & O(n) \\
\hline
\end{array} \right].
\]

Next, we provide a linearity test for non-signaling functions with low error. The statement involves the notion of self-correction for a non-signaling function, which we define first.

**Definition 4.1.** The self-correction of a \( k \)-non-signaling function \( F : \{0,1\}^n \to \{0,1\} \) is the \( \lfloor k/2 \rfloor \)-non-signaling function \( \hat{F} : \{0,1\}^n \to \{0,1\} \) defined as follows. Given a set \( \{x_1, \ldots, x_s\} \subseteq \{0,1\}^n \) with \( s \leq \lfloor k/2 \rfloor \), \( \hat{F}(x_1, \ldots, x_s) \) samples uniform and independent \( z_1, \ldots, z_s \in \{0,1\}^n \) and answers each \( x_i \) with \( F(z_i + x_i) - F(z_i) \).

**Theorem 7** (formal statement of Theorem 3; proved in Section 9). Let \( k, \tilde{k} \in \mathbb{N} \) and \( \varepsilon \in (0, 1/400] \) be such that \( k = \Omega(\tilde{k} \cdot (\tilde{k} + \log k)) \). Suppose that \( F : \{0,1\}^n \to \{0,1\} \) is a \( k \)-non-signaling function such that \( \Pr_{x,y,F}[F(x) + F(y) = F(x + y)] \geq 1 - \varepsilon \). Then there exists a linear \( \tilde{k} \)-non-signaling function \( L : \{0,1\}^n \to \{0,1\} \) such that for all query sets \( Q \subseteq \{0,1\}^n \) with size \( |Q| \leq \tilde{k} \) and for all events \( E \subseteq \{0,1\}^Q \) it holds that

\[
\left| \Pr[\hat{F}(Q) \in E] - \Pr[L(Q) \in E] \right| \leq (6|Q| + 3\sqrt{\varepsilon}).
\]

The above result enables us to compile any nsLPCP into a corresponding nsPCP, via the usual transformation that runs a linearity test and then uses self-correction to ask the nsLPCP queries.
Lemma 4.2 (formal statement of Lemma 1.3; proved in Section 6). The classical LPCP-to-PCP compiler works in the non-signaling setting, giving us that, for every $\varepsilon \in (0,1)$,

\[
\begin{align*}
\text{nsLPCP} & \quad \begin{bmatrix}
\text{soundness error:} & 1 - \varepsilon \\
\text{randomness:} & r \\
\text{proof length:} & \ell \\
\text{query complexity:} & q \\
\text{locality:} & k \\
\text{query sampler time:} & T_1 \\
\text{decision predicate time:} & T_2
\end{bmatrix} \\
\subseteq \text{nsPCP} & \quad \begin{bmatrix}
\text{soundness error:} & 1 - \min\left\{ \frac{1}{400}, \frac{\varepsilon^2}{(6q+4)\ell} \right\} \\
\text{randomness:} & r + (q + 2)\ell \\
\text{proof length:} & 2\ell \\
\text{query complexity:} & 2q + 3 \\
\text{locality:} & O\left(\frac{k}{\ell} (k + \log \frac{1}{\varepsilon})\right) \\
\text{query sampler time:} & T_1 + O(q\ell) \\
\text{decision predicate time:} & T_2 + O(q)
\end{bmatrix},
\end{align*}
\]

Proof of Theorem 5. By invoking the transformation from Lemma 4.2 on the ALMSS verifier from Theorem 6, we directly obtain the inclusion stated in Theorem 5.

The PCP produced by Lemma 4.2 can at best have soundness error $1 - \frac{1}{400}$, even if the soundness error of the underlying linear PCP is sub-constant. In Appendix B, we give a similar compiler where the resulting PCP has almost the same soundness error as that of the given linear PCP.
5 Repeated tests against non-signaling functions

We state and prove several generic facts about the acceptance probability of repeated tests against non-signaling functions. Using these statements, we then prove a generic lemma showing that, given a verifier V, if the t-repeated verifier V′ is exp(−t)−sound against non-signaling functions, then V is O(1)−sound against non-signaling functions. The fact that this implication holds generically is non-trivial, as the reverse implication, which is typically used in the classical setting to establish that V′ has a much smaller soundness error than V, does not hold in the non-signaling setting.

A q-query test over a domain D is a pair (R, Π), where R is a distribution over q-size subsets of D, and Π: {0, 1}^q → {0, 1} is a predicate; the test first samples (x_1, ..., x_q) according to R, makes the queries to F, receives answers (a_1, ..., a_q) ∈ {0, 1}^q, and outputs Π(a_1, ..., a_q).

Given a k-non-signaling function F: D → {0, 1} and a q-query test over D, for any t ≤ [k/q] we can perform the test t times in parallel and thereby define t boolean random variables X_1, ..., X_t by setting X_i = 1 if the i-th test passes. These boolean variables are identically distributed but need not be independent, and their joint distribution is permutation invariant (for every b ∈ {0, 1}^t and permutation π: [t] → [t] it holds that Pr[X_i = b_i ∀ i ∈ [t]] = Pr[X_{π(i)} = b_i ∀ i ∈ [t])].

The lemmas below hold for boolean variables whose joint distribution is permutation-invariant.

**Lemma 5.1** (adapted from [KRR14, Claim 6.2]). Let r, d, t ∈ N be such that 0 ≤ r < d. Let X_1, ..., X_{t+d} be boolean random variables with a permutation-invariant joint distribution. Define the following events:

- H^d_{d−r} is the event that X_i = 1 for at least d − r indices i ∈ {1, ..., d};
- T^t_i is the event that X_i = 1 for all i ∈ {d + 1, ..., d + t}.

Then

\[ \Pr_{X_1, ..., X_{t+d}} [−H^d_{d-r} \land T^t_i] \leq \left( \frac{t+d-r-1}{t+d} \right)^r \left( \frac{d}{t+d} \right)^{r+1}. \]

In particular, if \( \Pr [T^t_i] > 0 \) then

\[ \Pr_{X_1, ..., X_{t+d}} [−H^d_{d-r} \mid T^t_i] \leq \frac{1}{\Pr [T^t_i]} \left( \frac{d}{t+d} \right)^{r+1}. \]

**Proof.** Letting I := \{i ∈ [t + d] : X_i = 0\}, we view the event −H^d_{d-r} ∩ T^t_i as happening in two steps. First, we sample X_1, ..., X_{t+d}, which determines |I|. Second, we re-index the samples by selecting a uniformly random permutation π: [t + d] → [t + d] to determine I. This implies that

\[ \Pr_{X_1, ..., X_{t+d}} [−H^d_{d-r} \land T^t_i] = \sum_{j=0}^{t+d} \Pr_{X_1, ..., X_{t+d}} [−H^d_{d-r} \land T^t_i \mid |I| = j] \cdot \Pr_{X_1, ..., X_{t+d}} [|I| = j]. \]

The above step uses the fact that the joint distribution of X_1, ..., X_{t+d} is permutation-invariant. If j ≤ r, then \( \Pr_{π} [−H^d_{d-r} \land T^t_i \mid |I| = j] = 0 \), because in this case H^d_{d-r} always holds. If j ≥ r + 1, then \( \Pr_{π} [−H^d_{d-r} \land T^t_i \mid |I| = j] = (t−d) / (t−d) \), because the event holds if and only if \( \pi(d+1), ..., \pi(d+t) \) are all not in I, which happens with probability \( (t−d−j) / (t−d) \) when |I| = j. Hence, \( \Pr_{π} [−H^d_{d-r} \land T^t_i \mid |I| = j] \leq (t−d−r−1) / (t−d) \), for j ≥ r + 1. Therefore,

\[ \Pr_{X_1, ..., X_{t+d}} [−H^d_{d-r} \land T^t_i] \]

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Then, with a permutation-invariant joint distribution. Define the following events:

\[ \Pr \left[ \left. X_1, \ldots, X_{t+d} \right| |I| = j \right] \]

\[ \leq \sum_{j=r+1}^{t+d} \frac{(t+d-(r+1))}{t} \cdot \Pr_{X_1, \ldots, X_{t+d}} \left[ |I| = j \right] \]

\[ \leq \left( \frac{t+d-r-1}{t} \right) \cdot \left( \frac{t+d}{t} \right) \]

\[ \, \square \]

Next we prove a variant of Lemma 5.1, where we consider the event that \( X_i = 1 \) for most indices \( i \) in \( \{d + 1, \ldots, d + t\} \), as opposed to all indices \( i \) in that set.

**Lemma 5.2.** Let \( d, t \in \mathbb{N} \), and let \( 0 \leq \varepsilon < \delta \leq 1 \). Let \( X_1, \ldots, X_{t+d} \) be boolean random variables with a permutation-invariant joint distribution. Define the following events:

- \( H_d^{(1-\delta)d} \) is the event that \( X_i = 1 \) for at least \((1-\delta)d\) indices \( i \in \{1, \ldots, d\} \);
- \( T^t_{(1-\varepsilon)t} \) is the event that \( X_i = 1 \) for at least \((1-\varepsilon)t\) indices \( i \in \{d + 1, \ldots, d + t\} \).

Then,

\[ \Pr_{X_1, \ldots, X_{t+d}} \left[ \left. -H_d^{(1-\delta)d} \land T^t_{(1-\varepsilon)t} \right| |I| = j \right] \leq e^{-\frac{1}{8}(\delta-\varepsilon)^2 \min\{t,d\}}. \]

In particular, if \( \Pr \left[ T^t_{(1-\varepsilon)t} \right] > 0 \) then

\[ \Pr_{X_1, \ldots, X_{t+d}} \left[ \left. -H_d^{(1-\delta)d} \mid T^t_{(1-\varepsilon)t} \right| |I| = j \right] \leq \frac{1}{\Pr \left[ T^t_{(1-\varepsilon)t} \right]} \cdot e^{-\frac{1}{8}(\delta-\varepsilon)^2 \min\{t,d\}}. \]

**Proof.** The proof strategy follows that for Lemma 5.1. Letting \( I := \{i \in [t + d] : X_i = 0\} \), we view the event \(-H_d^{(1-\delta)d} \land T^t_{(1-\varepsilon)t} \) as happening in two steps. First, we sample \( X_1, \ldots, X_{t+d} \), which determines \( |I| \). Second, we re-index the samples by selecting a uniformly random permutation \( \pi : [t + d] \to [t + d] \) to determine \( I \). This implies that

\[ \Pr_{X_1, \ldots, X_{t+d}} \left[ \left. -H_d^{(1-\delta)d} \land T^t_{(1-\varepsilon)t} \right| |I| = j \right] = \sum_{j=0}^{t+d} \Pr_{X_1, \ldots, X_{t+d}} \left[ \left. -H_d^{(1-\delta)d} \land T^t_{(1-\varepsilon)t} \right| |I| = j \right] \cdot \Pr_{X_1, \ldots, X_{t+d}} \left[ |I| = j \right]. \]

The above step uses the fact that the joint distribution of \( X_1, \ldots, X_{t+d} \) is permutation-invariant. By standard martingale concentration bounds (see, e.g., [McD98]), for \( j < \frac{\varepsilon+\delta}{2}(t + d) \) it holds that

\[ \Pr_{\pi} \left[ \left. -H_d^{(1-\delta)d} \mid |I| = j \right] \leq e^{-\frac{1}{2}(\delta-\varepsilon)^2 d} \leq e^{-\frac{1}{8}(\delta-\varepsilon)^2 d}, \]

and for \( j \geq \frac{\varepsilon+\delta}{2}(t + d) \) it holds that

\[ \Pr_{\pi} \left[ \left. T^t_{(1-\varepsilon)t} \mid |I| = j \right] \leq e^{-\frac{1}{2}(\delta-\varepsilon)^2 t} \leq e^{-\frac{1}{8}(\delta-\varepsilon)^2 t}. \]

Therefore, all probabilities \( \Pr_{\pi} \left[ \left. -H_d^{(1-\delta)d} \land T^t_{(1-\varepsilon)t} \mid |I| = j \right] \) are upper bounded by the maximum among the two bounds, which completes the proof of the lemma. \( \square \)

We now state and prove two lemmas that lower bound the probability that \( X_1 = \cdots = X_t = 1 \) in terms of the probability that \( X_i = 1 \). Both lemmas play a key role in our results and, to the best of our knowledge, have not appeared in prior work studying non-signaling functions.

Suppose that \( \Pr[X_i = 1] = \delta \). If the \( X_i \)'s are independent, then \( \Pr[X_1 = \cdots = X_t = 1] = \delta^t \). The same implication is not true when the \( X_i \)'s are arbitrarily correlated. Yet, if their joint distribution is permutation-invariant, then we can conclude that \( \delta^t \) (approximately) lower bounds this probability.
Lemma 5.3. Let $X_1, \ldots, X_r$ be boolean random variables with a permutation-invariant joint distribution. Let $t \in \mathbb{N}$ and $0 < \tau < \delta < 1$ be such that $r \geq \frac{4t}{\tau}$ and $\Pr[X_i = 1] = \delta$ for every $i \in [r]$. Then,

$$\Pr_{X_1, \ldots, X_t}[X_1 = \cdots = X_t = 1] \geq \frac{3}{16} \cdot \tau \cdot (\delta - \tau)^t.$$  

Proof. Let $X = \sum_{i=1}^r X_i$. We have that $\mathbb{E}[X] = \sum_{i=1}^r \mathbb{E}[X_i] = r\delta$. Since the joint distribution of the $X_i$’s is permutation-invariant, the probability that $X_1, \ldots, X_t$ are all 1 can be viewed as the probability that $X_i = 1$ for all $i$ in a random subset of $t$ $X_i$’s. Therefore,

$$\Pr_{X_1, \ldots, X_t}[X_1 = \cdots = X_t = 1] = \sum_{i=0}^r \Pr[X = i] \cdot \binom{t}{i} = \mathbb{E} \left[ \binom{X}{t} \right].$$

Below we lower bound the foregoing expression.

Since $\mathbb{E}[X] = r\delta$ we have $\Pr[X \geq (\delta - \tau/2)r] \geq \tau/2$. Therefore,

$$\Pr_{X_1, \ldots, X_t}[X_1 = \cdots = X_t = 1] = \mathbb{E} \left[ \binom{X}{t} \right] \geq \tau/2 \cdot \frac{\binom{(\delta-\tau/2)r}{t}}{\binom{t}{t}}.$$

It remains to lower bound the above binomial coefficient. We do this using the following claim.

Claim 5.4. Let $\beta \in [0, 1]$ and $r, t \in \mathbb{N}$ such that $r/t \geq 2/\beta$. Then $\frac{\binom{\beta r}{t}}{\binom{t}{t}} \geq \frac{3}{8} \cdot (\beta - 2t/r)^t$.

Proof. By Stirling’s approximation we have

$$\frac{3}{8} \leq \frac{\binom{n}{i}}{\sqrt{i^{i-n} \cdot i^{n-i}}} \leq 1.$$ 

Therefore

$$\frac{\binom{\beta r}{t}}{\binom{t}{t}} \geq \frac{\sqrt{\frac{\beta r}{t(\beta r-t)}} \cdot \frac{\binom{\beta r}{t}}{t^{(\beta r-t)}}}{\sqrt{\frac{r}{t(r-t)}} \cdot \frac{\binom{t}{t}}{t^{t}}} = \frac{3}{8} \cdot \frac{(1 - \frac{t}{r})}{\frac{1}{\beta r}} \cdot \frac{(\beta r - t)^t}{(r-t)^t}.$$

Using the assumption that $0 < \frac{t}{r} < \frac{t}{\beta r} < 1$ and the fact that $\frac{1-x}{e} \leq (1-x)^{1/x} \leq \frac{1-x/2}{e}$ for all $x \in (0, 1]$, we get that

$$\frac{\binom{\beta r}{t}}{\binom{t}{t}} \geq \frac{3}{8} \cdot \left( \frac{1 - \frac{t}{r}}{e} \right)^t \cdot \left( \beta - \frac{t}{r} \right)^t \geq \frac{3}{8} \cdot \left( \frac{1 - \frac{t}{r}}{1 - \frac{t}{2\beta r}} \right)^t \cdot \left( \beta - \frac{t}{r} \right)^t \geq \frac{3}{8} \cdot \left( \beta - 2\frac{t}{r} \right)^t,$$

where the last inequality uses the assumption that $r/t \geq 2/\beta$. \hfill \Box

Using Claim 5.4, we get that

$$\Pr_{X_1, \ldots, X_t}[X_1 = \cdots = X_t = 1] \geq \tau/2 \cdot \frac{3}{8} \cdot (\delta - \tau/2 - 2t/r)^t \geq \frac{3}{16} \cdot \tau \cdot (\delta - \tau)^t,$$

where the last inequality uses the fact that $r \geq \frac{4t}{\tau}$. \hfill \Box

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Next we prove a variant of Lemma 5.3, where we consider the event that among $X_1, \ldots, X_t$ most of the $X_i$’s are 1, as opposed to all $X_i$’s. We note that in this lemma, we do not require that $X_1, \ldots, X_t$ have a permutation-invariant joint distribution.

**Lemma 5.5.** Let $X_1, \ldots, X_t$ be boolean random variables with an arbitrary joint distribution, and let $\rho, \delta \in [0, 1]$ be such that $\Pr[X_i = 1] \geq 1 - \delta \rho$ holds for all $i \in [n]$. Then

$$\Pr_{X_1, \ldots, X_t} \left[ X_i = 1 \text{ for at least } (1 - \rho) t \text{ indices } i \in [t] \right] \geq 1 - \delta.$$

**Proof.** Let $X = t - \sum_{i=1}^t X_i$. Observe that $\mathbb{E}[X] \leq \rho \delta t$, and so $\Pr[X \geq \rho t] \leq \mathbb{E}[X]/\rho t = \delta$. This implies that $\Pr[X < \rho t] \geq 1 - \delta$, and the lemma follows. $\Box$

Returning to the given PCP verifier $V$, we wish to relate the soundness error of $V$ and of the $t$-repeated verifier $V^t$. In the classical setting, one typically shows that $V$ has soundness error $O(1)$, and then it follows trivially that $V^t$ has soundness error $\exp(-t)$. However, this approach fails in the non-signaling setting because soundness error $O(1)$ for $V$ does not (in general) imply soundness error $\exp(-t)$ for $V^t$. Surprisingly, using Lemma 5.3 we can show that the other direction holds, namely that soundness error $\exp(-t)$ for $V^t$ implies soundness error $O(1)$ for $V$.

**Lemma 5.6.** Let $V$ be a $q$-query PCP verifier. Suppose that there exist $t \in \mathbb{N}$ and $u \in [0, 1)$ such that $V^t$ is $u^t$-sound against $k$-non-signaling functions. If $t > \frac{6 - \log_2(3 - 3u)}{\log_2(\frac{3u}{2})}$, then $V$ is $\frac{3 + u}{4}$-sound against $\max(k, \frac{16tq}{1-u})$-non-signaling functions. For example, if $V^t$ is $\left(\frac{9}{10}\right)^t$-sound against $k$-non-signaling functions for some $t \geq 100$, then $V$ is $39/40$-sound against $\max(k, 160tq)$-non-signaling functions.

One can prove variations on the above lemma, based on the settings of $\delta$ and $\tau$ in Lemma 5.3. We also note that the same lemma holds when $V$ is a linear PCP verifier and all non-signaling functions considered are linear.

**Proof.** Let $k' := \max(k, \frac{16tq}{1-u})$, and let $F$ be a $k'$-non-signaling function. Since $k' \geq k$, $\Pr[V^t \text{ accepts } F] \leq u^t$. Hence, it suffices to show that if $\Pr[V \text{ accepts } F] > \frac{3 + u}{4}$, then $\Pr[V^t \text{ accepts } F] > u^t$.

Indeed, let $\delta := \frac{3 + u}{4}$, and $\tau := \frac{1-u}{4}$, and suppose that $F$ is a $k'$-non-signaling function such that $\Pr[V \text{ accepts } F] > \frac{3 + u}{4} = \delta$. Since $k' \geq \frac{16tq}{1-u}$, we can apply Lemma 5.3 with the above choices of $\delta$ and $\tau$ to get that $\Pr[V^t \text{ accepts } F] \geq \frac{3}{16} \cdot \frac{1-u}{4} \cdot (\frac{3+u}{2})^{3t}$, which is strictly greater than $u^t$ since $t > \frac{6 - \log_2(3 - 3u)}{\log_2(\frac{3u}{2})}$, as required. $\Box$

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6 From linear PCPs to standard PCPs

In this section we use Theorem 7 (testing linearity with low error) to prove Lemma 4.2, which shows that any non-signaling linear PCP can be compiled into a corresponding non-signaling PCP. This is the non-signaling analogue of the classical compiler from linear PCPs to PCPs [ALMSS98].

6.1 The PCP verifier

Let \( V_{LPCP} = (Q_{LPCP}, D_{LPCP}) \) be a linear PCP verifier that has oracle access to a linear \( k \)-non-signaling function \( \mathcal{L} : \{0,1\}^\ell \rightarrow \{0,1\} \). We need to construct a PCP verifier \( V_{PCP} = (Q_{PCP}, D_{PCP}) \) that has oracle access to a \( k_{PCP} \)-non-signaling function \( \mathcal{F} : \{0,1\}^\ell \rightarrow \{0,1\} \).

The construction of \( V_{PCP} \) from \( V_{LPCP} \) is exactly the same as in the classical case. That is, \( V_{PCP} \) runs the [BLR93] linearity test, and then runs \( V_{LPCP} \), simulating its oracle by querying the proof in a self-correcting manner. More specifically, for every query \( a \in \{0,1\}^\ell \) that \( V_{LPCP} \) makes to its oracle, \( V_{PCP} \) samples a random \( z_a \in \{0,1\}^\ell \), queries \( \mathcal{F}(z_a) \), \( \mathcal{F}(z_a + a) \), and answers \( V_{LPCP} \) with \( \mathcal{F}(z_a + a) - \pi(z_a) \). Overall we can express \( V_{PCP} = (Q_{PCP}, D_{PCP}) \) in the following way:

The query sampler \( Q_{PCP} \) works as follows:
1. Sample uniformly random \( x, y \in \{0,1\}^\ell \).
2. Run \( Q_{LPCP} \) to obtain a set of queries \( S_{LPCP} \subseteq \{0,1\}^\ell \) for the linear PCP.
3. For each \( a \in S_{LPCP} \) sample \( z_a \in \{0,1\}^\ell \) independently and uniformly at random.
4. Output the query set \( S_{PCP} = \{x, y, x + y\} \cup \{z_a, z_a + a : a \in S_{LPCP}\} \).

Given the answers \( \mathcal{F}(S_{PCP}) \) from the proof, define \( \hat{\mathcal{F}}(S_{LPCP}) \) by letting \( \hat{\mathcal{F}}(a) = \mathcal{F}(z_a + a) - \mathcal{F}(a) \) for all \( a \in S_{LPCP} \). The decision predicate \( D_{PCP} \) accepts if and only if \( \mathcal{F}(x) + \mathcal{F}(y) = \hat{\mathcal{F}}(S_{LPCP}) \) and \( D_{LPCP}(\hat{\mathcal{F}}(S_{LPCP})) = 1 \).

6.2 Parameters

The parameters of the new verifier \( V_{PCP} \) follow from those of \( V_{LPCP} \) and Theorem 7

Randomness: \( V_{PCP} \) uses \( 2\ell \) random bits for the linearity test, plus additional \( r + q\ell \) random bits to simulate \( V_{LPCP} \). The total number of random bits is \( r + (q+2)\ell \).

Length: The length is \( 2^\ell \), because the honest proof in the completeness case is a (classical) function \( \pi : \{0,1\}^\ell \rightarrow \{0,1\} \), which takes \( 2^\ell \) bits to write down.

Queries: \( V_{PCP} \) makes 2 queries for every query made by \( V_{LPCP} \), plus an additional 3 queries for the linearity test. The total number of queries is \( 2q + 3 \).

Locality: We need \( k_{PCP} = O((k_{LPCP}/\epsilon) \cdot (k_{LPCP} + \log \frac{1}{\delta})) \) so that we can apply Theorem 7.

Input obliviousness: The linearity test clearly does not depend on the input (or even \( L \)), and the self-correction only depends on its input set \( S_{LPCP} \). Hence, if \( V_{LPCP} \) has input-oblivious queries, then so does \( V_{PCP} \).

Query sampler time: The query sampler of \( V_{PCP} \) samples uniformly random \( x, y \in \{0,1\}^\ell \). Then, it runs the query sampler of \( V_{LPCP} \) to generate \( q \) queries, and finally generates the self-corrected queries. This takes a total of \( T_1 + O(q\ell) \) time.

Decision predicate time: The decision predicate of \( V_{PCP} \) first runs the linearity test predicate. Then, it computes the self-corrected answers to the queries of \( V_{LPCP} \), and finally runs the decision predicate of \( V_{LPCP} \). This takes a total of \( T_2 + O(q) \) time.
6.3 Completeness

Completeness is straightforward. Suppose that there exists a linear function $F : \{0, 1\}^\ell \rightarrow \{0, 1\}$ such that $V^F_{LPCP}$ always accepts. We claim that $F$ is also accepted by $V_{PCP}$ with probability 1. Indeed, if $F$ is linear then $\hat{F} \equiv F$, and hence $\Pr[D_{LPCP}(\hat{F}(Q_{LPCP})) = 1] = \Pr[D_{LPCP}(F(Q_{LPCP})) = 1]$. Therefore,

$$\Pr[V^F_{PCP} = 1] = \Pr[F(x) + F(y) = F(x+y) \land D_{LPCP}(\hat{F}(Q_{LPCP}))] = \Pr[D_{LPCP}(F(Q_{LPCP})) = 1] = 1.$$

6.4 Soundness

Suppose that there exists a $k_{PCP}$-non-signaling function $F : \{0, 1\}^\ell \rightarrow \{0, 1\}$ that is accepted by $V_{PCP}$ with probability greater than $1 - \gamma$ for some $\gamma \leq \min\{\frac{1}{400}, \frac{\epsilon^2}{(6q+1)^2}\}$. In particular, this implies that $\Pr[F(x) + F(y) = F(x+y)] \geq 1 - \gamma \geq 1 - \frac{1}{400}$. Since $k_{PCP} = O(\frac{k_{LPCP}}{\epsilon} \cdot (k_{LPCP} + \log \frac{1}{\epsilon}))$, by Theorem 7 there exists a linear $k_{LPCP}$-non-signaling function $L$ such that $\Delta_q(L, \hat{F}) \leq (6q + 3)\sqrt{\gamma}$. Then,

$$\Pr[V^L_{LPCP} = 1] \geq \Pr[V^F_{LPCP} = 1] - (6q + 3)\sqrt{\gamma} > 1 - \gamma - (6q + 3)\sqrt{\gamma} \geq 1 - \epsilon,$$

since $\gamma \leq \frac{\epsilon^2}{(6q+1)^2}$. This is a contradiction, and so the soundness of $V_{LPCP}$ implies soundness of $V_{PCP}$.
7 Linear local assignment generators

In this section we define the notion of a linear local assignment generator, and prove that, in order to establish soundness of an LPCP, it suffices to construct a linear local assignment generator with sufficiently small error. This is one of the steps in the proof of Theorem 6 in Section 8, and we state it separately because it works for any nsLPCP (not just the ALMSS verifier analyzed in Section 8).

Informally, a linear local assignment generator for a circuit $C$ on input $x$ is a linear non-signaling function that satisfies each local constraint in the computation of $C$ on $x$ with high probability. This definition augments earlier ones in [KRR14; PR17] with the linearity condition (see Section 2.3).

**Definition 7.1.** Let $C: \{0,1\}^n \rightarrow \{0,1\}$ be a boolean circuit with $N$ wires, and let $x \in \{0,1\}^n$ be an input for it. A linear $k$-local assignment generator for $(C, x)$ with error $\varepsilon$ is a linear $k$-non-signaling function $A: \{0,1\}^N \rightarrow \{0,1\}$ that satisfies the following.

1. The input wires $w_1, \ldots, w_n$ are consistent with $x$: for every $i \in [n]$, $\Pr[A(e_i) = x_i] \geq 1 - \varepsilon$.
2. The output wire $w_N$ is 1: $\Pr[A(e_N) = 1] \geq 1 - \varepsilon$.
3. AND/OR gates compute correctly: for every binary gate $g$ with inputs $w_{i_1}, w_{i_2}$ and output $w_j$,
   \[
   \Pr[g(A(e_{i_1}), A(e_{i_2})) = A(e_j)] \geq 1 - \varepsilon.
   \]
4. NOT gates compute correctly: for every unary gate $g$ with input $w_i$ and output $w_j$,
   \[
   \Pr[g(A(e_i)) = A(e_j)] \geq 1 - \varepsilon.
   \]

We prove that the existence of a linear $k$-local assignment generator for $(C, x)$, with sufficiently large locality $k$ and sufficiently small error $\varepsilon$, implies that $C(x) = 1$. This tells us that, in order to prove soundness for a nsLPCP, it suffices to show that the existence of a linear non-signaling function accepted by the nsLPCP verifier with good probability implies the existence of a linear local assignment generator with low error.

**Lemma 7.2.** Let $C: \{0,1\}^n \rightarrow \{0,1\}$ be a boolean circuit with $N$ wires, and let $x \in \{0,1\}^n$ be an input for it. If there exists a linear $k$-local assignment generator $A$ for $(C, x)$ with error $\varepsilon$, where $k \geq 2 \log_2(1/\varepsilon)$ and $\varepsilon \log_2(1/\varepsilon) \leq \frac{1}{2N}$, then $C(x) = 1$.

The bound on the error is almost tight, regardless of the locality $k$.

**Lemma 7.3.** For sufficiently large $N$ and any $n$, there exist a boolean circuit $C: \{0,1\}^n \rightarrow \{0,1\}$ with $N$ wires and an input $x \in \{0,1\}^n$ such that $C(x) = 0$ and yet exists a linear $2^N$-local assignment generator $A: \{0,1\}^N \rightarrow \{0,1\}$ for $(C, x)$ with error $\varepsilon \geq \frac{\ln N}{N}$.

We first prove Lemma 7.2 in Section 7.1, and then prove Lemma 7.3 in Section 7.2.

### 7.1 Proof of Lemma 7.2

Fix a boolean circuit $C$ (see Section 3.1) with $W$ wires, and an input $x \in \{0,1\}^n$ for $C$; denote by $H$ the number of layers in $C$. For $h \in [H]$, let $W_h$ be the set of wires in layer $h$, and let $D_h = \operatorname{span}\{e_i : w_i \in W_h\}$. In other words, $D_h$ is the subset of $\{0,1\}^N$ containing all vectors $v$ such that $v_i = 0$ for all indices $i \in [N]$ where the wire $w_i$ is not in layer $h$. There is a unique correct transcript of the wires of the circuit for the computation of $C$ on $x$; we let $\operatorname{Tr}_{(C, x)}: \{0,1\}^N \rightarrow \{0,1\}$ denote the linear extension of this transcript.
For a set \( S \subseteq \{0,1\}^N \) let \( E(S) \) be the event that \( \mathcal{A}(S) \) returns values that are consistent with the correct transcript. That is, \( E(S) \) occurs if and only if for each \( \alpha \in S \) it holds that \( \mathcal{A}(\alpha) = \langle \text{Tr}_{C,x} \rangle(\alpha) \).

When the set \( S = \{\alpha_1, \ldots, \alpha_t\} \) is specified explicitly by its elements, we allow ourselves to write \( E(\alpha_1, \ldots, \alpha_t) \) to refer to \( E(\{\alpha_1, \ldots, \alpha_t\}) \).

Let \( t = \lceil \log_2(1/\varepsilon) \rceil \) be a parameter. By our choice of parameters, \( k \geq 2t \geq t + 2 \). For each layer \( h \in [H] \) of the circuit, let \( z_{h,1}, \ldots, z_{h,t} \in D_h \) be chosen uniformly at random. That is, each \( z_{h,i} \) is a uniformly random linear combination of the wires in layer \( h \) of \( C \). For a layer \( h \in [H] \) of the circuit \( C \), let \( E_h \) denote the event that \( E(z_{h,1}, \ldots, z_{h,t}) \) holds, i.e., all values output by \( \mathcal{A}(z_{h,1}, \ldots, z_{h,t}) \) are consistent with the correct transcript of \( C \) on input \( x \). Note that the randomness of \( E_h \) is over both \( z_{h,1}, \ldots, z_{h,t} \) and over \( \mathcal{A} \); \( E_h \) will be the event that we will condition on in layer \( h \) (see Section 2.3).

In the proof we use the following definition.

**Definition 7.4.** A set \( S \subseteq \{0,1\}^N \) (of linear combinations of wires) is \( p \)-good for layer \( h \) if \( \Pr[E(S)|E_h] \geq p. \) When \( S = \{\alpha\} \) is a singleton, we say that \( \alpha \) is \( p \)-good for layer \( h \).

Note that in the foregoing definition we do not insist that the elements of \( S \) belong to \( D_h \), and it is well-defined to ask whether a linear combination of wires outside \( D_h \) is good for layer \( h \).

We first argue that if \( E_h \) occurs with probability at least 0.5, then any linear combination of wires in layer \( h \) is \( p \)-good for \( p \) close to 1.

**Claim 7.5.** If \( \Pr[E_h] \geq 0.5 \), then every \( \alpha_h \in D_h \) is \( (1 - 2^{-t+1}) \)-good for layer \( h \).

**Proof.** Fix \( \alpha_h \in D_h \) and let \( \beta_1, \ldots, \beta_t \in D_h \) be chosen independently and uniformly at random. For \( i = 1, \ldots, t \), pick \( z_{h,i} \in \{\beta_i, \alpha_h + \beta_i\} \) uniformly; note that the \( z_{h,i} \)’s are uniform in \( D_h \).

We claim that \( \Pr[E(z_{h,1}, \ldots, z_{h,t}) | \neg E(\alpha_h)] \leq 2^{-t} \). We first observe that if \( E(\alpha_h) \) does not hold then \( \mathcal{A}(\alpha_h) \) is not consistent with the transcript \( \text{Tr}_{C,x} \), i.e., \( \mathcal{A}(\alpha_h) \neq \langle \text{Tr}_{C,x} \rangle(\alpha_h) \). Therefore, if \( k \geq t + 1 \), then by the linearity of \( \mathcal{A} \) for each \( i \in [t] \) it holds that either \( \mathcal{A}(\beta_i) \neq \langle \text{Tr}_{C,x} \rangle(\beta_i) \) or \( \mathcal{A}(\alpha_h + \beta_i) \neq \langle \text{Tr}_{C,x} \rangle(\alpha_h + \beta_i) \). In other words, if \( E(\alpha_h) \) does not hold, then either \( E(\beta_i) \) does not hold or \( E(\alpha_h + \beta_i) \) does not hold for all \( i \in [t] \) independently of each other.\(^{10}\) Hence, \( \Pr[E(z_{h,1}, \ldots, z_{h,t}) | \neg E(\alpha_h)] \leq 2^{-t} \).

Therefore, since \( k \geq t + 1 \), we have

\[
\Pr[\neg E(\alpha_h) | E_h] = \frac{\Pr[\neg E(\alpha_h) \land E_h]}{\Pr[E_h]} \leq 2 \Pr[\neg E(\alpha_h) \land E(z_{h,1}, \ldots, z_{h,t})] \quad \text{(since } \Pr[E_h] \geq 0.5 \text{)}
\]
\[
\leq 2 \Pr[E(z_{h,1}, \ldots, z_{h,t}) | \neg E(\alpha_h)] \leq 2^{-t+1}.
\]

By definition, this means that \( \alpha_h \) is \( (1 - 2^{-t+1}) \)-good for layer \( h \), as claimed. \( \square \)

We now use the above claim, and the condition that \( \mathcal{A} \) satisfies the constraints representing the computation of \( C \) on \( x \) with high probability, to show that if \( E_h \) occurs with probability at least 0.5, then any linear combination of wires in the next layer is good with high probability.

**Claim 7.6.** If \( \Pr[E_h] \geq 0.5 \), then every \( \alpha_{h+1} \in D_{h+1} \) is \( (1 - (2^{-t+2} + 2\varepsilon) \cdot N_{h+1}) \)-good for layer \( h \).

\(^{10}\)This is the only place in the proof where we use the fact that the \( z_{h,i} \)'s are random.
We now provide the details of the induction and corresponding conclusion.

Proof. Fix $\alpha_{h+1} \in D_{h+1}$. If $\alpha_{h+1} = e_i$ is a standard basis vector with $1$ in a coordinate corresponding to some wire $w_i$ in $W_{h+1}$, then there is some gate $g$ in $C$ taking two inputs $w_{i_1}, w_{i_2}$ in layer $h$ and outputting $w_i$ in layer $h+1$. By Claim 7.5 and union bound, we have

$$\Pr[E(\{e_{i_1}, e_{i_2}\}) | E_h] \geq 1 - 2^{-t+2}.$$ 

By the assumption that $\mathcal{A}$ is a local assignment generator for $(C, x)$ with error $\varepsilon$ it holds that

$$\Pr[g(\mathcal{A}(e_{i_1}), \mathcal{A}(e_{i_2})) = \mathcal{A}(e_i)] \geq 1 - \varepsilon,$$

which implies that

$$\Pr[-E(e_i) \wedge E(\{e_{i_1}, e_{i_2}\})] \leq \Pr[-E(e_i) | E(\{e_{i_1}, e_{i_2}\})] \leq \varepsilon.$$

Hence,

$$\Pr[-E(e_i) | E_h] = \Pr[-E(e_i) \wedge E(\{e_{i_1}, e_{i_2}\}) | E_h] + \Pr[-E(e_i) \wedge E(\{e_{i_1}, e_{i_2}\}) | E_h]$$

\leq $\Pr[-E(\{e_{i_1}, e_{i_2}\}) | E_h] + \Pr[-E(e_i) \wedge E(\{e_{i_1}, e_{i_2}\}) | E_h]$

\leq $2^{-t+2} + \frac{\Pr[-E(e_i) \wedge E(\{e_{i_1}, e_{i_2}\})]}{\Pr[E_h]}$

\leq $2^{-t+2} + 2\varepsilon$ (since $\Pr[E_h] \geq 0.5$).

Therefore, for $\alpha_{h+1} = e_i$ it holds that $\mathcal{A}(\alpha_{h+1})$ conditioned on $E_h$ is consistent with $\text{Tr}_{(C, x)}$ with probability at least $1 - (2^{-t+2} + 2\varepsilon)$, i.e., such $\alpha_{h+1}$ is $(1 - (2^{-t+2} + 2\varepsilon))$-good for layer $h$.

If $\alpha_{h+1}$ is not a standard basis vector, then it is a linear combination of at most $N_{h+1}$ basis vectors, and since $\mathcal{A}$ is a linear non-signaling function we get by union bound that $\mathcal{A}(\alpha_{h+1})$ is consistent with $\text{Tr}_{(C, x)}$ with probability at least $1 - (2^{-t+2} + 2\varepsilon) \cdot N_{h+1}$, as required.

We now use the foregoing claims to prove Lemma 7.2.

We use induction on the depth of $C$ to show that $E_h$ holds with high probability for every layer $h$. Specifically, we show that, for every layer $h$, $\Pr[E_h] \geq 1 - (2^{-t+2} + 2\varepsilon) \cdot \left(\sum_{i=1}^{N_i} N_i\right) \cdot t$. In particular, this will imply that $\Pr[E_H] \geq 0.5$, which will allow us to conclude that $C(x) = 1$.

Note that by the choice of the parameters, the obtained bound on $\Pr[E_h]$ will be greater than $0.5$ for every layer $h \in [H]$, which will allow us to apply Claim 7.6 in the induction step. Indeed, for $k \geq 2 \log_2(1/\varepsilon) \cdot \varepsilon \log_2(1/\varepsilon) \leq \frac{1}{2 \varepsilon^2}$, and $t = \lceil \log_2(1/\varepsilon) \rceil \geq \log_2(1/\varepsilon) - 1$, we have

$$1 - (2^{-t+2} + 2\varepsilon) \cdot \left(\sum_{i=1}^{N_i} N_i\right) \cdot t \geq 1 - (2^{-t+2} + 2\varepsilon) \cdot N \cdot t \geq 1 - 10\varepsilon \cdot N \log_2(1/\varepsilon) \geq 0.5.$$

We now provide the details of the induction and corresponding conclusion.

Base case. For the base case, we show that $E_1$ holds with probability $1 - (2^{-t+2} + 2\varepsilon) \cdot N_1 \cdot t$. Fix $i \in [n]$ so that $w_i$ is the $i$-th input wire in $C$. Then, by the assumption that $\mathcal{A}$ is a linear local assignment generator for $(C, x)$ with error $\varepsilon$ we have $\Pr[\mathcal{A}(e_i) = x_i] \geq 1 - \varepsilon$. By the linearity of $\mathcal{A}$ and union bound, if $\alpha \in D_1$ then

$$\Pr[\mathcal{A}(\alpha) = \langle \text{Tr}_{(C, x)}, \alpha \rangle] \geq 1 - \varepsilon n = 1 - \varepsilon N_1 \geq 1 - (2^{-t+2} + 2\varepsilon) \cdot N_1 \cdot t,$$
which proves the base case of the induction.

**Induction step.** For the induction step we assume that \( \Pr[E_h] \geq 1 - (2^{-t+2} + 2\varepsilon)(\sum_{i=1}^{h} N_i) t \), and prove the corresponding statement for \( E_{h+1} \). By the earlier discussion we have \( \Pr[E_h] \geq 0.5 \), and thus by applying Claim 7.6 we get that for all \( z_{h+1} \in D_{h+1} \) it holds that \( \Pr[E(z_{h+1}) \mid E_h] \geq 1 - (2^{-t+2} + 2\varepsilon) \cdot N_{h+1} \). Therefore, for any \( z_{h+1,1}, \ldots, z_{h+1,t} \in D_{h+1} \), by union bound we have

\[
\Pr[E_{h+1} \mid E_h] = \Pr[E(\{z_{h+1,1}, \ldots, z_{h+1,t}\}) \mid E_h] \geq 1 - (2^{-t+2} + 2\varepsilon) \cdot N_{h+1} \cdot t .
\]

This implies that

\[
\Pr[E_{h+1}] \geq \Pr[E_{h+1} \mid E_h] \cdot \Pr[E_h] = (1 - (2^{-t+2} + 2\varepsilon) \cdot N_{h+1} \cdot t) \cdot \Pr[E_h] \\
\geq (1 - (2^{-t+2} + 2\varepsilon) \cdot N_{h+1} \cdot t) \left(1 - (2^{-t+2} + 2\varepsilon) \cdot \sum_{i=1}^{h} N_i \cdot t \right) \\
\geq 1 - (2^{-t+2} + 2\varepsilon) \cdot \sum_{i=1}^{h+1} N_i \cdot t ,
\]

which proves the induction step.

**Proving that \( C(x) = 1 \).** It remains to show that \( \Pr[E_H] \geq 0.5 \) implies that \( C(x) = 1 \). Let \( w_N \) be the output wire of \( C \), and let \( e_N \) be the corresponding basis vector. By Claim 7.5 we have \( \Pr[A(e_N) = C(x) \mid E_H] \geq 1 - 2^{-t+1} \). This implies that

\[
\Pr[A(e_N) = C(x)] \geq \Pr[A(e_N) = C(x) \wedge E_H] = \Pr[A(e_N) = C(x) \mid E_H] \cdot \Pr[E_H] \\
\geq (1 - 2^{-t+1})/2 \\
\geq 1/4 .
\]

On the other hand, by the assumption that \( A \) is a linear local assignment generator for \( (C, x) \) with error \( \varepsilon \), we have \( \Pr[A(e_N) \neq 1] \leq \varepsilon \). Therefore,

\[
\Pr[C(x) \neq 1] = \Pr[C(x) \neq A(e_N) \wedge A(e_N) = 1] + \Pr[C(x) = A(e_N) \wedge A(e_N) \neq 1] \\
\leq \Pr[C(x) \neq A(e_N)] + \Pr[A(e_N) \neq 1] \\
\leq 3/4 + \varepsilon \\
< 1 .
\]

Since \( C(x) \) is a deterministic value, we conclude that \( C(x) = 1 \), thus proving Lemma 7.2.

### 7.2 Proof of Lemma 7.3

Let \( C : \{0,1\}^n \rightarrow \{0,1\} \) be a boolean circuit with only OR gates and let the input \( x \) be \( 0^n \); note that \( C(x) = 0 \). Fix \( \varepsilon > 0 \). We first define a distribution \( D \) over functions \( f : [N] \rightarrow \{0,1\} \).

To sample a function \( f \) from \( D \), we assign values to \( f \) one layer at a time. For each wire \( w_i \) in layer 0 (the input layer), \( f(i) = x_i \) with probability \( 1 - \varepsilon \). Then, for each \( w_i \) in layer \( h + 1 \), letting \( w_{i_1} \) and \( w_{i_2} \) be the wires in layer \( h \) that are inputs to the gate \( g \) that computes \( w_i \), \( f(i) \) is sampled
as follows. If either \( f(i_1) \) or \( f(i_2) \) is 1, then \( f(i) = 1 \) with probability 1. Otherwise, \( f(i) = 0 \) with probability \( 1 - \varepsilon \), and 1 with probability \( \varepsilon \). In other words, if at least one of the input wires is 1, then the gate is computed correctly, otherwise it is computed correctly with probability \( 1 - \varepsilon \).

Given the distribution \( \mathcal{D} \), we define \( \mathcal{A} : \{0,1\}^N \rightarrow \{0,1\} \) to be the linear extension of \( \mathcal{D} \) to the domain \( \{0,1\}^N \). Namely, while \( \mathcal{D} \) is a distribution over functions \( f : [N] \rightarrow \{0,1\} \), \( \mathcal{A} \) is the corresponding distribution over linear functions \( h_f : \{0,1\}^N \rightarrow \{0,1\} \) such that \( h_f \) is the linear function defined by \( h_f(e_i) = f(i) \) for all \( i \in [N] \).

Note that \( \mathcal{A} \) is a \( 2^N \)-non-signaling function (since it is a distribution), and is linear since it is a distribution over linear functions. Moreover, \( \mathcal{A} \) satisfies all gate constraints locally with probability at least \( 1 - \varepsilon \), since \( \mathcal{D} \) either computes the gate correctly with probability \( 1 - \varepsilon \), or it computes the gate correctly (but with incorrect inputs) with probability 1. Also, \( \mathcal{A} \) assigns correct values to the input wires with probability \( 1 - \varepsilon \), since \( \mathcal{D} \) does. Hence, to show that \( \mathcal{A} \) is a linear local assignment generator with error \( \varepsilon \), it remains to show that \( \Pr[\mathcal{A}(e_N) = 1] \geq 1 - \varepsilon \).

We have that \( \Pr[\mathcal{A}(e_N) = 1] = 1 - (1 - \varepsilon)^N \), since \( \mathcal{A} \) assigns 1 to the output gate if at least one wire has the wrong value, and the probability that all wires are correct is \( (1 - \varepsilon)^N \). For \( \varepsilon \geq \frac{\ln N}{N} \) and \( N \) sufficiently large, \( 1 - (1 - \varepsilon)^N \geq 1 - \varepsilon \). This is because \( \frac{(1 - \varepsilon)^N}{\varepsilon} \leq e^{-\varepsilon N - \ln \varepsilon} \leq e^{-\ln N - \ln \ln N} = e^{-\ln N} < 1 \) for \( N \) sufficiently large.

In sum, \( \mathcal{A} \) is a \( 2^N \)-local assignment generator for \( C \) on input \( x \) with error \( \varepsilon \), but \( C(x) = 0 \).
8 The linear PCP verifier of ALMSS

In this section we prove Theorem 6.

Throughout, we fix a language \( L \in \text{DSIZE}(S) \) and an instance \( x \in \{0,1\}^n \). We denote by \( C := C_n \) the boolean circuit with \( N := S(n) \) wires such that \( x \in L \) if and only if \( C(x) = 1 \).

8.1 The ALMSS verifier

We analyze the LPCP verifier of [ALMSS98] (the “ALMSS verifier”), which we now recall.

The computation of \( C \) on \( x \) is viewed as a system of \( M := N + 1 \) equations \( \{P_j(w) = c_j\}_{j \in [M]} \) over \( N \) boolean variables \( w = (w_1, \ldots, w_N) \in \{0,1\}^N \), where \( P_1, \ldots, P_M : \{0,1\}^N \to \{0,1\} \) are quadratic polynomials (each involving at most three variables in \( w \)) and \( c_1, \ldots, c_M \) are boolean constants. Each variable represents the value of one of the wires of \( C \) on input \( x \); \( w_1, \ldots, w_n \) are the \( n \) input wires and \( w_N \) is the output wire. There are three types of constraints:

Input consistency: For every \( j \in \{1, \ldots, n\} \), \( P_j(w) := w_i \) and \( c_j := x_j \).

Gate consistency: For every \( j \in \{n+1, \ldots, N\} \),

- If the wire represented by the variable \( w_j \) is an output of an AND gate \( g \) where the inputs to \( g \) are given by \( w_{j_1}, w_{j_2} \), then \( P_j(w) := w_j - w_{j_1} \cdot w_{j_2} \) and \( c_j := 0 \).
- If the wire represented by the variable \( w_j \) is an output of an OR gate \( g \) where the inputs to \( g \) are given by \( w_{j_1}, w_{j_2} \), then \( P_j(w) := w_j - w_{j_1} - w_{j_2} + w_{j_1} \cdot w_{j_2} \) and \( c_j := 0 \).
- If the wire represented by the variable \( w_j \) is an output of a NOT gate \( g \) where the input to \( g \) is given by \( w_{j_1} \), then \( P_j(w) := w_j - w_{j_1} \) and \( c_j := 1 \).

Accepting output: \( P_M(w) := w_N \) and \( c_M := 1 \).

We overload notation and use \( P_j \) to also denote the upper triangular matrix in \( \{0,1\}^{N^2} \) such that \( P_j(w) = \langle P_j, w \otimes w \rangle \); that is, if \( P_j(w) = \sum a_i w_i + \sum_{1 \leq i < j \leq N} a_i' w_i w_j \), then \( P_j \) has \( a_i \) in the diagonal entry \((i,i)\) and \( a_i' \) in the entry \((i,i')\), for \( 1 \leq i < i' \leq N \). Also, for \( a \in \{0,1\}^N \), \( D_a \) is the diagonal matrix in \( \{0,1\}^{N^2} \) whose diagonal is \( a \). The LPCP verifier of [ALMSS98] is as follows.

The ALMSS verifier, given input \( x \in \{0,1\}^n \) and oracle access to a linear non-signaling function \( L : \{0,1\}^{N^2} \to \{0,1\} \), works as follows:

1. Use the circuit \( C \) and input \( x \) to construct the matrices \( P_1, \ldots, P_M \in \{0,1\}^{N^2} \) and constants \( c_1, \ldots, c_M \in \{0,1\} \), which represent the computation of \( C \) on \( x \).
2. Sample \( u, v \in \{0,1\}^N \) and \( s \in \{0,1\}^M \) uniformly and independently at random.
3. Query the oracle \( L \) on the 4-element set \( S = \{D_u, D_v, u \otimes v, \sum_{j=1}^M s_j P_j\} \).
4. Check that \( L(D_u)L(D_v) = L(u \otimes v) \) and \( L(\sum_{j=1}^M s_j P_j) = \sum_{j=1}^M s_j c_j \).

We also analyze the \( t \)-wise repetition of the ALMSS verifier (the “\( t \)-repeated ALMSS verifier”).

Parameters. The parameters of the ALMSS verifier are the usual ones, which we review.

The ALMSS verifier uses \( 2N + M = O(S) \) random bits and makes 4 queries to the proof. The proof is expected to be a linear function \( \pi : \{0,1\}^{N^2} \to \{0,1\} \), which can be represented via a vector of \( N^2 = O(S^2) \) bits. The queries are input oblivious because they depend only on the language \( L \) and input size \( n \), which in turn determine the circuit \( C \).
The time complexity of the query sampler is dominated by the cost of constructing the matrices $P_1, \ldots, P_M \in \{0,1\}^{N^2}$ and constants $c_1, \ldots, c_M \in \{0,1\}$, which takes time $O(S^2)$.

The time complexity of the decision predicate is $O(n)$, because it takes $O(n)$ time to compute $\sum_{j=1}^n s_j c_j$, and the rest of the computation takes constant time. Note that the query sampler precomputes $\sum_{j=n+1}^M s_j c_j$, since it does not depend on $x$ (only $c_1, \ldots, c_n$ depend on $x$).

**Completeness.** Completeness of the ALMSS verifier is as in the classical setting. Namely, suppose that $C(x) = 1$, and let $w = (w_1, \ldots, w_N) \in \{0,1\}^N$ be such that $w_i$ is the value of the $i$-th wire when $C$ computes on $x$. Then, the linear function defined as $\pi(Z) := \langle w \otimes w, Z \rangle = \sum_{i,i' \in [N]} w_i w_{i'} \cdot Z_{i,i'}$ for all $Z \in \{0,1\}^{N^2}$ is accepted with probability 1. Indeed,

1. for every $u, v \in \{0,1\}^N$ it holds that
   \[
   \pi(D_u)\pi(D_v) = \left( \sum_{i \in [N]} w_i u_i \right) \left( \sum_{i' \in [N]} w_{i'} v_{i'} \right) = \sum_{i,i' \in [N]} w_i w_{i'} u_i v_{i'} = \pi(u \otimes v) ;
   \]

2. for every $s \in \{0,1\}^M$ it holds that
   \[
   \pi\left( \sum_{j=1}^M s_j P_j \right) = \sum_{j=1}^M s_j \pi(P_j) = \sum_{j=1}^M s_j c_j ,
   \]
   because $\pi(P_j) = \langle w \otimes w, P_j \rangle = P_j(w) = c_j$ for all $j \in [M]$.

**Soundness.** Soundness against non-signaling functions is quite unlike the classical case, and is discussed in the next few sub-sections. We shall prove that if the ALMSS verifier accepts a given linear $k$ non-signaling function $L: \{0,1\}^{N^2} \rightarrow \{0,1\}$ with sufficiently large constant probability, then $C(x) = 1$, and thus $x \in L$.

### 8.2 Step 1: Constructing the linear local assignment generator

We have proved in Section 7 that, to conclude that $C(x) = 1$, it suffices to construct a linear local assignment generator with sufficiently small error and sufficiently large locality. Here we show how to construct a linear local assignment generator from any linear $k'$-non-signaling function $L'$ (possibly different from $L$) that satisfies every test of the ALMSS verifier with large probability.

**Lemma 8.1.** Let $k' \geq 4$. Suppose that $L': \{0,1\}^{N^2} \rightarrow \{0,1\}$ is a linear $k'$-non-signaling function for which the following holds.

1. Every tensor test is satisfied with probability $1 - \varepsilon_1$
   \[\forall x, y \in \{0,1\}^N, \Pr_{L'}[L'(D_x)L'(D_y)] = L'(x \otimes y)] \geq 1 - \varepsilon_1 .\]

2. Every satisfiability test is satisfied with probability $1 - \varepsilon_2$
   \[\forall s \in \{0,1\}^M, \Pr_{L'}[L'\left( \sum_{j=1}^M s_j P_j \right)] = \sum_{j=1}^M s_j c_j ] \geq 1 - \varepsilon_2 .\]
Then there exists a linear \( k' \)-local assignment generator \( A \) for \((C, x)\) with error \( \varepsilon' \leq \varepsilon_1 + \varepsilon_2 \).

**Proof.** Being a linear \( k' \)-non-signaling function, \( \mathcal{L}' \) can be represented uniquely as a quasi-distribution over linear functions \( f: \{0, 1\}^N \to \{0, 1\} \). For each \( B \in \{0, 1\}^N \), let \( \ell_B \in \mathbb{R} \) be the weight that \( \mathcal{L}' \) assigns to the linear function \((B, \cdot)\). Define \( A: \{0, 1\}^N \to \{0, 1\} \) to be the quasi-distribution over linear functions \( f: \{0, 1\}^N \to \{0, 1\} \) that assigns to each linear function \((\alpha, \cdot)\) the weight \( p_\alpha \), where

\[
p_\alpha := \sum_{B \in \{0, 1\}^N \atop \text{s.t. diag}(B)=\alpha} \ell_B ,
\]

where \( \text{diag}(B) \in \{0, 1\}^N \) is the vector on the diagonal of the matrix \( B \). We now prove that \( A \) is a \( k' \)-local assignment generator for \((C, x)\) with error \( \varepsilon' \leq \varepsilon_1 + \varepsilon_2 \).

**Locality.** The quasi-distribution \( A \) is \( k' \)-local because, for every subset \( S \subseteq \{0, 1\}^N \) of size at most \( k' \) and every string \( \vec{b} \in \{0, 1\}^S \),

\[
\Pr \left[ A(S) = \vec{b} \right] = \sum_{\alpha \in \{0, 1\}^N \atop \forall x \in S \ (\alpha, x) = \vec{b}_x} p_\alpha
\]

\[
= \sum_{\alpha \in \{0, 1\}^N \atop \forall x \in S \ (\alpha, x) = \vec{b}_x} \left( \sum_{B \in \{0, 1\}^N \atop \text{s.t. diag}(B)=\alpha} \ell_B \right)
\]

\[
= \sum_{B \in \{0, 1\}^N \atop \text{s.t. } \forall x \in S \ (\alpha, x) = \vec{b}_x} \ell_B
\]

\[
= \Pr \left[ \forall x \in S \ (B, D_x) = \vec{b}_x \right] \geq 0 . \quad \text{(as } \mathcal{L}' \text{ is } k' \text{-local)}
\]

**Satisfying the constraints.** The quasi-distribution \( A \) satisfies all \( M \) constraints if and only if for for every constraint \( P_j \) it holds that \( P_j(A(S_j)) = c_j \), where the set \( S_j \) contains the indices for the variables influencing \( P_j \), and \( P_j \) is interpreted as the polynomial that accepts only these variables (i.e., ignoring the variables in the coordinates \([N] \setminus S_j\) ). Recall that each \( P_j \) depends on at most 3 variables. Therefore, we may query \( A \) on such sets \( S_j \), provided that \( A \) is \( k' \)-local for \( k' \geq 3 \).

Below we show that for each \( j \in [M] \) it holds that

\[
\Pr[P_j(A(S_j)) = c_j] \geq 1 - (\varepsilon_1 + \varepsilon_2) . \tag{2}
\]

Suppose first that \( P_j(\cdot) = c_j \) is a constraint that does not involve a quadratic term, so that it is either an input consistency constraint or a NOT gate constraint. We will prove the statement only for the case of an input consistency constraint; the case of a NOT gate is analogous, and we omit its proof. An input consistency constraint has the form \( w_i = c_j \) for some \( i \in [N] \) and \( c_j \in \{0, 1\} \). Therefore,

\[
\Pr[P_j(A(S_j)) = c_j] = \Pr[A(e_i) = c_j] = \Pr[\mathcal{L}'(D_{e_i}) = c_j] = \Pr[\mathcal{L}'(P_j) = c_j] \geq 1 - \varepsilon_2 .
\]

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Otherwise, $P_j(\cdot) = c_j$ involves a quadratic term. More specifically, $P_j$ either corresponds to an AND gate, and hence is of the form $w_i - w_{i_1} \cdot w_{i_2} = c_j$, or it corresponds to an OR gate, and is of the form $w_i - w_{i_1} - w_{i_2} + w_{i_1} \cdot w_{i_2} = c_j$. Below we prove that $\Pr[P_j(A(S_j)) = c_j] \geq 1 - (\varepsilon_1 + \varepsilon_2)$ for the case of an AND gate; the case of the OR gate is analogous, and we omit its proof.

If $P_j$ corresponds to an AND gate, then

$$\Pr[P_j(A(S_j)) = c_j] = \Pr[A(e_i) - A[e_{i_1}] \cdot A(e_{i_2}) = c_j] = \Pr[\mathcal{L}'(D_{e_i}) - \mathcal{L}'(D_{e_{i_1}}) \cdot \mathcal{L}'(D_{e_{i_2}}) = c_j]$$

Since $\mathcal{L}'$ individually satisfies each tensor test constraints with probability $1 - \varepsilon_1$ and $k' \geq 4$,

$$\Pr[\mathcal{L}'(D_{e_i}) - \mathcal{L}'(D_{e_{i_1}}) \mathcal{L}'(D_{e_{i_2}}) = c_j] \geq \Pr[\mathcal{L}'(D_{e_i}) - \mathcal{L}'(e_{i_1} \otimes e_{i_2}) = c_j] - \varepsilon_1 .$$

We conclude that

$$\Pr[P_j(A(S_j)) = c_j] \geq \Pr[\mathcal{L}'(D_{e_i}) - \mathcal{L}'(e_{i_1} \otimes e_{i_2}) = c_j] - \varepsilon_1$$

$$= \Pr[\mathcal{L}'(D_{e_i} - e_{i_1} \otimes e_{i_2}) = c_j] - \varepsilon_1$$

$$= \Pr[\mathcal{L}'(P_j) = c_j] - \varepsilon_1$$

$$\geq 1 - \varepsilon_1 - \varepsilon_2 .$$

\[\square\]

### 8.3 Step 2: Reducing the soundness error

We now argue that if $\mathcal{L}$ passes the $t$-repeated ALMSS verifier with probability at least $\gamma$ and $k = \Omega(t)$, then there exists a (possibly different) linear $(k - 4t)$-non-signaling function $\mathcal{L}'$ that satisfies every constraint of the (non-repeated) ALMSS verifier with high probability $(1 - 2^{-O(t^2)})$. Namely, if $\mathcal{L}$ satisfies the $t$-repeated ALMSS verifier on average, then $\mathcal{L}'$ satisfies the ALMSS verifier with high probability in the worst case. Informally, $\mathcal{L}'$ equals $\mathcal{L}$ conditioned on the $t$-repeated verifier passing.

**Definition 8.2.** Given a linear $k$-non-signaling function $\mathcal{L} : \{0, 1\}^{N^2} \rightarrow \{0, 1\}$ and a parameter $t$, we let $\text{accept}_t$ denote the event that the $t$-repeated ALMSS verifier accepts $\mathcal{L}$, namely, that for uniformly random $u^{(1)}, v^{(1)}, \ldots, u^{(t)}, v^{(t)} \in \{0, 1\}^N$ and $s^{(1)}, \ldots, s^{(t)} \in \{0, 1\}^M$ it holds that

$$\forall i \in [t], \quad \mathcal{L}(D_{u^{(i)}}) \mathcal{L}(D_{v^{(i)}}) = \mathcal{L}(u^{(i)} \otimes v^{(i)}) \quad \text{and} \quad \mathcal{L}\left(\sum_{j=1}^M s_j^{(i)} P_j\right) = \sum_{j=1}^M s_j^{(i)} c_j .$$

For $k' := k - 4t$, we define $\mathcal{L}' : \{0, 1\}^{N^2} \rightarrow \{0, 1\}$ be the linear $k'$-non-signaling function defined as

$$\Pr_{\mathcal{L}'}[\mathcal{L}'(S) = \vec{b}] := \Pr_{\mathcal{L}}[\mathcal{L}(S) = \vec{b} | \text{accept}_t] = \frac{\Pr[\mathcal{L}(S) = \vec{b} \land \text{accept}_t]}{\Pr[\text{accept}_t]} .$$

**Lemma 8.3.** For every $k, t \in \mathbb{N}$ such that $k \geq 51t$ and linear $k$-non-signaling function $\mathcal{L} : \{0, 1\}^{N^2} \rightarrow \{0, 1\}$, if $\Pr[\text{accept}_t] \geq \gamma$ then

$$\forall x, y \in \{0, 1\}^N, \quad \Pr_{\mathcal{L}'}[\mathcal{L}'(D_x) \mathcal{L}'(D_y) = \mathcal{L}'(x \otimes y)] \geq 1 - \varepsilon_1$$

$$\forall s \in \{0, 1\}^M, \quad \Pr_{\mathcal{L}'}[\mathcal{L}'\left(\sum_{j=1}^M s_j P_j\right) = \sum_{j=1}^M s_j c_j] \geq 1 - \varepsilon_2 ,$$

where $\varepsilon_1, \varepsilon_2 \leq \frac{5}{\gamma} \left(\frac{4}{5}\right)^t$.  

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Proof. The proof relies on the following two claims.

**Claim 8.4.** For every $d \in \mathbb{N}$ such that $k \geq \max\{12d + 3t, 4t + 3\}$,

$$\forall x, y \in \{0, 1\}^N, \quad \Pr_{L'}[L'(D_x)L'(D_y) \neq L'(x \otimes y)] \leq \frac{4}{\gamma} \left( \frac{d}{t + d} \right)^{\frac{d+1}{2}}.$$

**Claim 8.5.** For every $d \in \mathbb{N}$ such that $k \geq \max\{t + 2d, 4t + 1\}$,

$$\forall s \in \{0, 1\}^M, \quad \Pr_{L'} \left[ L' \left( \sum_{j=1}^M s_j P_j \right) \neq \sum_{j=1}^M s_j c_j \right] \leq \frac{2}{\gamma} \left( \frac{d}{t + d} \right)^{\lfloor d/2 \rfloor}.$$

We prove the claims further below, and for now show how they imply the lemma. By the assumption that $k \geq 51t$, we can set $d := 4t$ in order to get

$$\Pr_{L'}[L'(D_x)L'(D_y) \neq L'(x \otimes y)] \leq \frac{4}{\gamma} \left( \frac{4t}{5t} \right)^{\lfloor t \frac{1}{2} \rfloor} \leq \frac{5}{\gamma} \left( \frac{4}{5} \right)^t,$$

$$\Pr_{L'} \left[ L' \left( \sum_{j=1}^M s_j P_j \right) \neq \sum_{j=1}^M s_j c_j \right] \leq \frac{2}{\gamma} \left( \frac{4t}{5t} \right)^{2t-1} \leq \frac{5}{\gamma} \left( \frac{4}{5} \right)^t. \quad \Box$$

We now prove the two claims.

**Proof of Claim 8.4.** Denote by $\text{tensor}_i$ the event “$L(D_{u^{(i)}})L(D_{v^{(i)}}) = L(u^{(i)} \otimes v^{(i)})$ for all $i \in [t]$". If we assume that $x$ and $y$ are chosen randomly, then a naive application of Lemma 5.1 implies that

$$\Pr_{L'}[L'(D_x)L'(D_y) \neq L'(x \otimes y)] = \frac{\Pr[\text{accept}_i | L(D_x)L(D_y) \neq L(x \otimes y) \wedge \text{accept}_i]}{\Pr[\text{accept}_i]} \leq \frac{\Pr[\text{accept}_i]}{\Pr[\text{accept}_i]} \leq \frac{1}{\gamma(t+1)}.$$

We will instead invoke the lemma more carefully and use the assumption that the locality of $L$ is sufficiently large to prove a stronger bound for fixed $x$ and $y$. Specifically, we use the assumption that $k \geq \max\{12d + 3t, 4t + 3\}$, together with the linearity of $L$ to prove that for all $x, y \in \{0, 1\}^N$, it holds that

$$\Pr[\text{accept}_i | L(D_x)L(D_y) \neq L(x \otimes y) \wedge \text{tensor}_i] \leq 4 \left( \frac{d}{t + d} \right)^{\frac{d+1}{2}}. \quad (3)$$

Indeed, this implies that

$$\Pr_{L'}[L'(D_x)L'(D_y) \neq L'(x \otimes y)] = \frac{\Pr[\text{accept}_i | L(D_x)L(D_y) \neq L(x \otimes y) \wedge \text{accept}_i]}{\Pr[\text{accept}_i]} \leq \frac{\Pr[\text{accept}_i | L(D_x)L(D_y) \neq L(x \otimes y) \wedge \text{tensor}_i]}{\Pr[\text{accept}_i]} \leq \frac{4}{\gamma} \left( \frac{d}{t + d} \right)^{\frac{d+1}{2}},$$

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which is what we wanted to prove.

It remains to argue that Eq. (3) holds. Let \( w_x^{(1)}, \ldots, w_x^{(d)}; w_y^{(1)}, \ldots, w_y^{(d)} \in \{0, 1\}^n \) be 2d vectors chosen uniformly at random. By linearity of \( \mathcal{L} \), for all \( i \in [d] \) we have

\[
\mathcal{L}(D_x)\mathcal{L}(D_y) = \mathcal{L}(D_{w_x^{(i)}})\mathcal{L}(D_{w_y^{(i)}}) + \mathcal{L}(D_{w_x^{(i)}})\mathcal{L}(D_{y+w_y^{(i)}}) + \mathcal{L}(D_{x+w_x^{(i)}})\mathcal{L}(D_{y+w_y^{(i)}}) \tag{4}
\]

For a positive integer \( r < d/4 \), Lemma 5.1 implies that

\[
\Pr \left[ \text{for at least } r \text{ indices } i \in [d] \right. \\
\mathcal{L}(D_{w_x^{(i)}})\mathcal{L}(D_{w_y^{(i)}}) \neq \mathcal{L}(w_x^{(i)} \otimes w_y^{(i)}) \land \text{tensor}_t \left. \right] \leq \left( \frac{d}{t+d} \right)^r,
\]

\[
\Pr \left[ \text{for at least } r \text{ indices } i \in [d] \right. \\
\mathcal{L}(D_{w_x^{(i)}})\mathcal{L}(D_{y+w_y^{(i)}}) \neq \mathcal{L}(x \otimes (y+w_y^{(i)})) \land \text{tensor}_t \left. \right] \leq \left( \frac{d}{t+d} \right)^r,
\]

\[
\Pr \left[ \text{for at least } r \text{ indices } i \in [d] \right. \\
\mathcal{L}(D_{x+w_x^{(i)}})\mathcal{L}(D_y) \neq \mathcal{L}((x+w_x^{(i)}) \otimes y) \land \text{tensor}_t \left. \right] \leq \left( \frac{d}{t+d} \right)^r.
\]

Denote by neq the union of the foregoing events. By union bound, \( \Pr[\text{neq}] \leq 4\left( \frac{d}{t+d} \right)^r \), and hence

\[
\Pr[\mathcal{L}(D_x)\mathcal{L}(D_y) \neq \mathcal{L}(x \otimes y) \land \text{tensor}_t] \leq \Pr[\mathcal{L}(D_x)\mathcal{L}(D_y) \neq \mathcal{L}(x \otimes y) \land \text{tensor}_t] - \Pr[\text{neq}] + 4\left( \frac{d}{t+d} \right)^r.
\]

We show next that \( \Pr[\mathcal{L}(D_x)\mathcal{L}(D_y) \neq \mathcal{L}(x \otimes y) \land \text{tensor}_t] - \Pr[\text{neq}] = 0 \). Indeed, if \( \text{neq} \) holds and \( 4r < d \), then there exists \( i \in [d] \) such that

\begin{itemize}
  \item \( \Pr[\mathcal{L}(D_{w_x^{(i)}})\mathcal{L}(D_{w_y^{(i)}}) = \mathcal{L}(w_x^{(i)} \otimes w_y^{(i)})] \),
  \item \( \Pr[\mathcal{L}(D_{w_x^{(i)}})\mathcal{L}(D_{y+w_y^{(i)}}) = \mathcal{L}(x \otimes (y+w_y^{(i)}))] \),
  \item \( \mathcal{L}(D_{x+w_x^{(i)}})\mathcal{L}(D_y) = \mathcal{L}((x+w_x^{(i)}) \otimes y) \), and
  \item \( \mathcal{L}(D_{x+w_x^{(i)}})\mathcal{L}(D_{y+w_y^{(i)}}) = \mathcal{L}((x+w_x^{(i)}) \otimes (y+w_y^{(i)})).
\end{itemize}

Therefore, by Eqs. (4) and (5), if \( \text{neq} \) holds and \( 4r < d \) then \( \mathcal{L}(x \otimes y) = \mathcal{L}(D_x)\mathcal{L}(D_y) \). This implies that \( \Pr[\mathcal{L}(D_x)\mathcal{L}(D_y) \neq \mathcal{L}(x \otimes y) \land \text{tensor}_t] - \Pr[\text{neq}] = 0 \), and so for \( r = \left\lfloor \frac{d-1}{4} \right\rfloor \) we get that

\[
\Pr[\mathcal{L}(D_x)\mathcal{L}(D_y) \neq \mathcal{L}(x \otimes y) \land \text{tensor}_t] \leq 4\left( \frac{d}{t+d} \right)^r = 4\left( \frac{d}{t+d} \right)^{\left\lfloor \frac{d-1}{4} \right\rfloor},
\]

which proves Eq. (3), and thus completes the proof of Claim 8.4. \( \square \)

**Proof of Claim 8.5.** Let \( s^{(1)}, \ldots, s^{(i)} \in \{0, 1\}^M \) be chosen uniformly at random, and denote by sat\(_t\) the event that

\[
\forall i \in [t] \quad \mathcal{L} \left( \sum_{j=1}^M s_j^{(i)} P_j \right) = \sum_{j=1}^M s_j^{(i)} c_j.
\]

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We use the assumption that \( k \geq \max\{t + 2d, 4t + 1\} \) to prove that for all \( s \in \{0, 1\}^M \) it holds that
\[
\Pr \left[ \mathcal{L} \left( \sum_{j=1}^{M} s_j P_j \right) \neq \sum_{j=1}^{M} s_j c_j \land \text{sat}_t \right] \leq 2 \left( \frac{d}{t + d} \right)^{d/2}.
\] (6)

Indeed, this implies that
\[
\Pr \left[ \mathcal{L}' \left( \sum_{j=1}^{M} s_j P_j \right) \neq \sum_{j=1}^{M} s_j c_j \right] = \Pr \left[ \mathcal{L} \left( \sum_{j=1}^{M} s_j P_j \right) \neq \sum_{j=1}^{M} s_j c_j \land \text{accept}_t \right] \leq \Pr \left[ \mathcal{L} \left( \sum_{j=1}^{M} s_j P_j \right) \neq \sum_{j=1}^{M} s_j c_j \land \text{sat}_t \right] \leq 2 \left( \frac{d}{t + d} \right)^{d/2},
\]
which is what we wanted to prove.

It remains to argue that Eq. (6) holds. Let \( y^{(1)}, \ldots, y^{(d)} \in \{0, 1\}^M \) be chosen uniformly at random. By linearity of \( \mathcal{L} \) and Lemma 5.1, we have
\[
\Pr \left[ \mathcal{L} \left( \sum_{j=1}^{M} s_j P_j \right) \neq \sum_{j=1}^{M} s_j c_j \land \text{sat}_t \right] = \sum_{s^{(1)}, \ldots, s^{(d)} \in \{0, 1\}^M} \Pr_{y^{(1)}, \ldots, y^{(d)} \in \{0, 1\}^M} \left[ \mathcal{L} \left( \sum_{j=1}^{M} (s_j + y_j^{(i)}) P_j \right) \neq \sum_{j=1}^{M} (s_j + y_j^{(i)}) c_j \land \text{sat}_t \right]
\]

for every index \( i \in [d] \).

By Lemma 8.3, there exists a linear \( k' \)-non-signaling function \( \mathcal{L}' \) with \( k' = k - 4t = 47t \) such that
\[
\forall x, y \in \{0, 1\}^N, \quad \Pr_{\mathcal{L}' \mathcal{L}'} \left[ \mathcal{L}'(D_x)\mathcal{L}'(D_y) = \mathcal{L}'(x \otimes y) \right] \geq 1 - \varepsilon_1
\]

\[\square\]

8.4 Step 3: Soundness of the ALMSS verifier

Fix \( N \geq 16 \), let \( t := 25 \log_2 N \), and suppose that \( \mathcal{L} \) is a linear \( k \)-non-signaling function \( \mathcal{L} \) that is accepted by the \( t \)-repeated ALMSS verifier with probability at least \( \gamma = \left( \frac{a}{10} \right)^t \), where \( k = 51t = O(\log N) \). We need to conclude that \( C(x) = 1 \).

By Lemma 8.3, there exists a linear \( k' \)-non-signaling function \( \mathcal{L}' \) with \( k' = k - 4t = 47t \) such that
\[
\forall x, y \in \{0, 1\}^N, \quad \Pr_{\mathcal{L}' \mathcal{L}'} \left[ \mathcal{L}'(D_x)\mathcal{L}'(D_y) = \mathcal{L}'(x \otimes y) \right] \geq 1 - \varepsilon_1
\]
∀ s ∈ \{0, 1\}^M, \quad \Pr_{\mathcal{C}'} \left[ \mathcal{L}' \left( \sum_{j=1}^{M} s_j P_j \right) = \sum_{j=1}^{M} s_j c_j \right] \geq 1 - \varepsilon_2 ,

with \( \varepsilon_1 + \varepsilon_2 \leq \frac{10}{\gamma} \left( \frac{4}{9} \right)^t = 10 \cdot \left( \frac{4}{9} \right)^t \).

By applying Lemma 8.1 to \( \mathcal{L}' \) we obtain a linear \( k' \)-local assignment generator \( \mathcal{A} : \{0, 1\}^N \rightarrow \{0, 1\} \) for \( (C, x) \) with error \( \varepsilon_1 + \varepsilon_2 = 10 \cdot \left( \frac{4}{9} \right)^t \).

By the choice of the parameters we have \( N \geq 16 \) and \( t = 25 \log_2 N \) so that \( 10 \cdot \left( \frac{4}{9} \right)^t < \frac{1}{40N^2} \) and \( t > 2 \log_2 (40N^2) \).

Therefore, \( \mathcal{A} \) is a linear \( k' \)-local assignment generator for \( (C, x) \) with error \( \varepsilon' = \frac{1}{40N^2} \), and locality \( k' \geq t = 25 \log_2 N \geq 2 \log_2 (40N^2) \geq 2 \log_2 \frac{1}{\varepsilon'} \). Moreover, \( \varepsilon' \log_2 (1/\varepsilon') \leq \frac{5+2 \log_2 N}{40N^2} < \frac{1}{20N} \).

Hence, by applying Lemma 7.2 to \( \mathcal{A} \), we get that \( C(x) = 1 \).

Finally, applying Lemma 5.6, we conclude that if a linear \( O(\log N) \)-non-signaling function \( \mathcal{L} \) is accepted by the basic ALMSS verifier with probability greater than \( \frac{39}{40} \), then \( C(x) = 1 \).
9 Testing linearity with low error

In this section we prove Theorem 7. Toward proving the theorem we define a repeated self-correction of \( \mathcal{F} \), which we denote by \( \mathcal{F}_{t,d} \), and prove that if \( \mathcal{F} \) passes the [BLR93] linearity test with probability \( 1 - \varepsilon \), then (i) \( \mathcal{F}_{t,d} \) is close to a linear \( k \)-non-signaling function \( \mathcal{L} \), and (ii) for all for all query sets \( S \subseteq \{0,1\}^n \) of size at most \( \bar{k} \) it holds that the statistical distance between \( \mathcal{F}(S) \) and \( \mathcal{F}_{t,d}(S) \) is upper bounded by \( O(|S| \sqrt{\varepsilon}) \). We proceed with the details below.

9.1 Defining \( \mathcal{F}_{t,d} \)

The definition of the repeated self-correction \( \mathcal{F}_{t,d} \) relies on components that we now introduce. First, we define a repeated and relaxed version of the linearity test, parameterized by \( t \in \mathbb{N} \) and \( \rho \in [0,1] \).

- **Query sampler:** \( Q_t \) samples \( 2t \) uniformly random vectors \( x^{(1)}, y^{(1)}, \ldots, x^{(t)}, y^{(t)} \in \{0,1\}^n \) and outputs the \( 3t \) queries \( (x^{(i)}, y^{(i)}, x^{(i)} + y^{(i)})_{i \in [t]} \). This is the basic linearity test repeated \( t \) times.

- **Decision predicate:** \( D_{t,\rho} \) receives \( 3t \) answers \( (a_{x^{(i)}}, a_{y^{(i)}}, a_{x^{(i)} + y^{(i)}})_{i \in [t]} \) to the queries, and accepts if and only if \( a_{x^{(i)}} + a_{y^{(i)}} = a_{x^{(i)} + y^{(i)}} \) for at least \((1 - \rho)t\) indices \( i \in [t] \). That is, the test accepts if at least a \((1 - \rho)\)-fraction of the repetitions pass.

Next, we define encoding and decoding procedures used to query \( \mathcal{F} \); they are parametrized by \( d \in \mathbb{N} \).

- **Encoding procedure:** given a set \( S = (x_1, \ldots, x_s) \subseteq \{0,1\}^n \), \( \text{Enc}_d \) outputs a set of vector pairs \( \{(z_j^{(i)}, z_j^{(i)} + x_j)\}_{i \in [d]} \) for each \( x_j \in S \) according to the following distribution: for each \( i \in [d] \), \( \text{Enc}_d \) samples \( n \) vectors \( y_1^{(i)}, \ldots, y_n^{(i)} \) from \( \{0,1\}^n \) independently, and sets \( z_j^{(i)} \) to be \( Mx_j \), where \( M \) is the \( n \times n \) matrix with \( y_l^{(i)} \) as the \( l \)-th column. Explicitly, \( z_j^{(i)} = \sum_{l=1}^n (x_j)_l \cdot y_l^{(i)} \), where we note that \( (x_j)_l \) is a bit, while \( y_l^{(i)} \) is a vector.

- **Decoding procedure:** given \( \mathcal{F}(\text{Enc}_d(S)) \) (the answers of \( \mathcal{F} \) to the encoding of \( S \)), \( \text{Dec}_d \) outputs the vector \( \bar{b} = (b_1, \ldots, b_s) \in \{0,1\}^S \) of answers defined as
  \[
  b_j = \max_{i \in [d]} \left\{ \mathcal{F}(z_j^{(i)} + x_j) - \mathcal{F}(z_j^{(i)}) \right\} .
  \]

Given the foregoing components, the **conditioned self-correction** \( \mathcal{F}_{t,d} \) of \( \mathcal{F} \) is defined as
\[
\Pr[\mathcal{F}_{t,d}(S) = \bar{b}] := \Pr[\text{Dec}_d(\mathcal{F}(\text{Enc}_d(S))) = \bar{b} \mid D_{t,\rho}(Q_t) = 1]
\]
for every \( S \subseteq \{0,1\}^n \) of some maximal size and \( \bar{b} \in \{0,1\}^S \). That is, on an input \( x \), the self-correction \( \mathcal{F}_{t,d} \) outputs the majority over \( d \) corrections, conditioned on the event that \((1 - \rho)\)-fraction of the \( t \) repetitions of the basic linearity test pass.

Before proceeding, we prove that \( \mathcal{F}_{t,d} \) is indeed a non-signaling function. While this may seem obvious, we note that there are choices of \( \text{Enc}_d \) and \( \text{Dec}_d \) for which \( \mathcal{F}_{t,d} \) will not be non-signaling.

**Lemma 9.1.** If \( \mathcal{F} \) is \( k \)-non-signaling and \( k \geq 2dk + 3t \), then \( \mathcal{F}_{t,d} \) is a \( \bar{k} \)-non-signaling.
Proof. We will prove that for all \( T \supseteq S \) and \( \vec{b} \in \{0,1\}^S \) it holds that \( \Pr[\mathcal{F}_{t,d}(T)\mid S = \vec{b}] = \Pr[\mathcal{F}_{t,d}(S) = \vec{b}] \). Given \( T \supseteq S \) and \( \vec{b} \in \{0,1\}^S \), we have
\[
\Pr[\mathcal{F}_{t,d}(T)\mid S = \vec{b}] = \Pr[\text{Dec}_d(\mathcal{F}(\text{Enc}_d(T)))\mid S = \vec{b} \land \text{Dec}_t(\mathcal{F}(Q_t)) = 1] .
\]
The specification of \( \text{Enc}_d, \text{Dec}_d \) and the non-signaling property of \( \mathcal{F} \) imply that \( \Pr[\text{Dec}_d(\mathcal{F}(\text{Enc}_d(T)))\mid S = \text{Dec}_d(\mathcal{F}(\text{Enc}_d(S)))] = 1 \). Indeed, by definition of \( \text{Enc}_d \) we have that \( \text{Enc}_d(S) \) is equal to \( \text{Enc}_d(T) \) when we ignore the encodings of the elements in \( T \setminus S \). Since \( \mathcal{F} \) is non-signaling, it follows that \( \mathcal{F}(\text{Enc}_d(S)) \) is equal to \( \mathcal{F}(\text{Enc}_d(T)) \) when in the latter distribution we ignore the answers provided by \( \mathcal{F} \) to the encodings of elements in \( T \setminus S \). Finally, since \( \text{Dec}_d \) decodes each element of \( S \) (or \( T \) separately), it follows that \( \text{Dec}_d(\mathcal{F}(\text{Enc}_d(S))) \) and \( \text{Dec}_d(\mathcal{F}(\text{Enc}_d(T))) \mid S \) are equal. Thus,
\[
\Pr[\text{Dec}_d(\mathcal{F}(\text{Enc}_d(T)))\mid S = \vec{b} \land \text{Dec}_t(\mathcal{F}(Q_t)) = 1] = \Pr[\text{Dec}_d(\mathcal{F}(\text{Enc}_d(S)))\mid S = \vec{b} \land \text{Dec}_t(\mathcal{F}(Q_t)) = 1] .
\]

Therefore,
\[
\Pr[\mathcal{F}_{t,d}(T)\mid S = \vec{b}] = \Pr[\text{Dec}_d(\mathcal{F}(\text{Enc}_d(T)))\mid S = \vec{b} \land \text{Dec}_t(\mathcal{F}(Q_t)) = 1] = \Pr[\text{Dec}_d(\mathcal{F}(\text{Enc}_d(S)))\mid S = \vec{b} \land \text{Dec}_t(\mathcal{F}(Q_t)) = 1] = \Pr[\mathcal{F}_{t,d}(S) = \vec{b}] ,
\]
and so \( \mathcal{F}_{t,d} \) is non-signaling. Note that all the probabilities are well-defined provided that \( k \geq 2d |S| + 3t \) (as \( \text{Enc}_d(S) \subseteq \text{Enc}_d(T) \)) which holds for \( |S| \leq \bar{k} \). This completes the proof of Lemma 9.1. \( \square \)

### 9.2 \( \mathcal{F}_{t,d} \) is close to linear

In this section we show that if \( \Pr[\text{Dec}_t(\mathcal{F}(Q_t)) = 1] \) is bounded away from zero, then \( \mathcal{F}_{t,d} \) is close to linear assuming that \( t \) and \( d \) are sufficiently large. Specifically, we prove the following lemma.

**Lemma 9.2.** For any \( \varepsilon > 0 \) and \( \gamma > 0 \), if \( t, d = \Omega(\bar{k} + \log \frac{1}{\varepsilon} + \log \frac{1}{\gamma}) \) and \( \Pr[\text{Dec}_t(\mathcal{F}(Q_t)) = 1] \geq \gamma \) and \( \rho \leq 1/20 \), then there exists a linear \( k \)-non-signaling function \( \mathcal{L} \) such that \( \Delta_k(\mathcal{L}, \mathcal{F}_{t,d}) < \varepsilon \).

**Proof.** The proof invokes the following theorem on linearity testing, due to [CMS18].

**Theorem 8 ([CMS18]).** Let \( \hat{\mathcal{F}} \) be a \( \bar{k} \)-non-signaling function such that for all \( x, y \in \{0,1\}^n \), \( \Pr[\hat{\mathcal{F}}(x) + \hat{\mathcal{F}}(y) = \hat{\mathcal{F}}(x + y)] \geq 1 - \varepsilon \) and \( \Pr[\hat{\mathcal{F}}(0^n) = 0] = 1 \). Then there exists a linear \( k \)-non-signaling function \( \mathcal{L} \) such that \( \Delta_k(\hat{\mathcal{F}}, \mathcal{L}) \leq 2^k \varepsilon \).

The key step of the proof of Lemma 9.2 will be to prove the following lemma.

---

\(^{11}\)Note that \( \text{Enc}_d \) is a randomized procedure. We assume both calls to \( \text{Enc}_d \) use the same randomness.
Lemma 9.3. Let $\overline{F}_{t,d}$ be the correction of $F$, and suppose that $\Pr[D_{t,\rho}(F(Q_t)) = 1] \geq \gamma$. If $\rho < 1/16$, then for all $x, y \in \{0,1\}^n$ it holds that

$$\Pr[\overline{F}_{t,d}(x) + \overline{F}_{t,d}(y) = F_{t,d}(x + y)] \geq 1 - \hat{\varepsilon}$$

for $\hat{\varepsilon} = \frac{6}{\gamma} \cdot e^{-\frac{(1/8 - \rho)^2 \min(t,d)}{8}} + c' \cdot d \cdot 2^{-4/16} + \frac{2}{\gamma} \cdot e^{-\frac{(1/8 - \rho)^2 \min(t,d)}{8}}$, where $c'$ is some absolute constant.

Indeed, by Theorem 8 this implies that there exists a linear $\bar{k}$-non-signaling function $L: \{0,1\}^n \rightarrow \{0,1\}$ such that

$$\Delta_k(\overline{F}_{t,d}, L) \leq 2^{2k} \cdot \hat{\varepsilon}.$$ 

Therefore, for $\rho \leq 1/20$ and $t, d = \Omega(\bar{k} + \log(1/\hat{\varepsilon}) + \log(1/\gamma))$ we get that

$$\Delta_k(\overline{F}_{t,d}, L) \leq 2^{2k} \cdot \left(\frac{6}{\gamma} \cdot e^{-O(t)} + c' \cdot d \cdot 2^{-O(t)} + \frac{2}{\gamma} \cdot e^{-O(t)}\right) < \varepsilon,$$

which proves Lemma 9.2.

We now prove Lemma 9.3. Recall that $\overline{F}_{t,d}$ is defined as $\text{Dec}_d \circ F \circ \text{Enc}_d$ conditioned on $D_{t,\rho}(F(Q_t)) = 1$. For the query set $S = \{x, y, x + y\}$, the encoding $\text{Enc}_d(S)$ can be split into three sets: $\{(z_x^{(i)}, z_x^{(i)} + x) : i \in [d]\}$, $\{(z_y^{(i)}, z_y^{(i)} + y) : i \in [d]\}$, and $\{(z_{x+y}^{(i)}, z_{x+y}^{(i)} + x + y) : i \in [d]\}$, where $z_x^{(1)}, \ldots, z_x^{(d)}; z_y^{(1)}, \ldots, z_y^{(d)} \in \{0,1\}^n$ are independent and uniform in $\{0,1\}^n$, and $z_{x+y}^{(i)} = z_x^{(i)} + z_y^{(i)}$ for all $i \in [d]$.

It will be convenient to define $b_x^{(i)} = F(z_x^{(i)} + x) - F(z_x^{(i)})$ for all $i \in [d]$. Similarly, define $b_y^{(i)} = F(z_y^{(i)} + y) - F(z_y^{(i)})$ and $b_{x+y}^{(i)} = F(z_{x+y}^{(i)} + x + y) - F(z_{x+y}^{(i)})$ for all $i \in [d]$. Using this notation, the decoding $\text{Dec}_d$ returns the answer $\text{maj}_{i \in [d]} \{b_x^{(i)}\}$ for $x$. Analogously, the decoding $\text{Dec}_d$ returns $\text{maj}_{i \in [d]} \{b_y^{(i)}\}$ for $y$, and $\text{Dec}_d(x + y) = \text{maj}_{i \in [d]} \{b_{x+y}^{(i)}\}$ for $x + y$.

The proof of Lemma 9.3 consists of the following 3 steps.

1. We first show that, if we condition on $D_{t,\rho}(F(Q_t)) = 1$, then with high probability there is some $b_x \in \{0,1\}$ such that $b_x^{(i)} = b_x$ for most indices $i \in [d]$. By symmetry, this holds for $y$ and $x + y$.

2. We then show that, if we condition on $D_{t,\rho}(F(Q_t)) = 1$, then with high probability $b_x^{(i)} + b_y^{(i)} = b_{x+y}^{(i)}$ for most indices $i \in [d]$.

3. Finally, by union bound all these events happen simultaneously with high probability, implying that $\overline{F}_{t,d}(x) + \overline{F}_{t,d}(y) = F_{t,d}(x + y)$ with high probability.

9.2.1 Step 1: Majority is a vast majority

Lemma 9.4. Let $\rho \in [0,1]$ and $t, d, r \in \mathbb{N}$ be parameters such that $d$ is even, $r < d/2$, and $\rho < r/4d$. Let $F$ be a $k$-non-signaling function with $k \geq 6d + 3t$. Fix $x \in \{0,1\}^n$, and let $x^{(1)}, \ldots, x^{(d)} \in \{0,1\}^n$ be independent uniformly random vectors. Denote by $\text{Agr}_d(x)$ the event that $\exists b \in \{0,1\}$ such that $F(x^{(b)} + x) - F(x^{(b)}) = b$ for at least $q$ indices $j \in [d]$. Then

$$\Pr[\text{Agr}_{d-q}(x) \mid D_{t,\rho}(F(Q_t)) = 1] \geq 1 - \frac{2}{\gamma} \cdot e^{-\frac{(r/4d - \rho)^2 \min(t,d)}{\gamma}} - \frac{c \cdot d \cdot 2^{-r/4}}{\gamma},$$

where $c$ is some absolute constant.
Proof. Instead of arguing about Agr_q(x) directly, we first define a related event PairAgr_{q'}(x) for some q' that depends on q, and show that \( \Pr\{\text{PairAgr}_{q'}(x) \mid D_{t, \rho}(Q_t) = 1\} \) is high. We then relate PairAgr_{q'}(x) to Agr_q(x).

Denote by PairAgr_{q'}(x) the event that \( \mathcal{F}(x^{(2j-1)} + x) - \mathcal{F}(x^{(2j-1)}) = \mathcal{F}(x^{(2j)}) - \mathcal{F}(x^{(2j)}) \) for at least \( q \) indices \( j \in [d/2] \). Note that if \( \mathcal{F}(x^{(2j-1)} + x) - \mathcal{F}(x^{(2j-1)}) \neq \mathcal{F}(x^{(2j)}) + \mathcal{F}(x^{(2j)}), \) then either \( \mathcal{F}(x^{(2j-1)} + x) + \mathcal{F}(x^{(2j)}) \neq \mathcal{F}(x^{(2j-1)} + x^{(2j)} + x) \) or \( \mathcal{F}(x^{(2j)}) + \mathcal{F}(x^{(2j-1)}) \neq \mathcal{F}(x^{(2j-1)} + x^{(2j)} + x) \). Therefore,

\[
\Pr\left[ \neg \text{PairAgr}_{q' \cdot \frac{1}{4}}(x) \mid D_{t, \rho}(Q_t) = 1 \right] \\
\leq \Pr\left[ \mathcal{F}(x^{(2j-1)} + x) + \mathcal{F}(x^{(2j)}) \neq \mathcal{F}(x^{(2j-1)} + x^{(2j)} + x) \mid D_{t, \rho}(Q_t) = 1 \right] \\
+ \Pr\left[ \mathcal{F}(x^{(2j)}) + \mathcal{F}(x^{(2j-1)}) \neq \mathcal{F}(x^{(2j-1)} + x^{(2j)} + x) \mid D_{t, \rho}(Q_t) = 1 \right].
\]

Since \( x^{(2j-1)} + x \) and \( x^{(2j)} \) are uniformly random and independent, by Lemma 5.2 each of the two terms is upper bounded by \( \frac{1}{4} \cdot e^{-\frac{(r/4d-\rho)}{s} \min(t, d)} \). Therefore,

\[
\Pr\left[ \neg \text{PairAgr}_{q' \cdot \frac{1}{4}}(x) \mid D_{t, \rho}(Q_t) = 1 \right] \leq \frac{2}{2} \cdot e^{-\frac{(r/4d-\rho)}{s} \min(t, d)}. \]

In the next claim we relate PairAgr_{q'}(x) to Agr_q(x).

Claim 9.5. Let \( d, r \in \mathbb{N} \) be parameters such that \( d \) is even and \( r < d/2 \). There exists an absolute constant \( c > 0 \) (independent of \( d \) or \( r \)) such that \( \Pr\{\text{PairAgr}_{q' \cdot \frac{1}{4}}(x) \mid \neg \text{Agr}_{d-r}(x)\} \leq c \cdot d \cdot 2^{-r/4} \).

Proof. For each \( i \in [d] \) let \( b_i = \mathcal{F}(x^{(i)} + x) - \mathcal{F}(x^{(i)}) \) \( \in \{0, 1\} \), let \( b^* = \text{maj}\{b_i\} \) and let \( q = \left| \{i \in [d] : b_i = b^*\} \right| \). Then Agr_{d-r}(x) holds if and only if \( q \geq d - r \). We stress that Agr_{d-r}(x) is independent of the order of the indices \( i \in [d] \). Next, we use the following lemma to relate PairAgr_{q'}(x) to Agr_q(x).

Lemma 9.6. Let \( d \) be an even integer, and let \( b_1, \ldots, b_d \in \{0, 1\} \) be such that not all \( b_i \)'s are equal. Let \( q = \left| \{i \in [d] : b_i = 0\} \right| \), and let \( \alpha = \min\left(\frac{d}{2}, 1 - \frac{1}{d}\right) \) for some \( \alpha \in (0, 0.5) \). If \( \sigma : [d] \rightarrow [d] \) is a uniformly random permutation of \([d]\), then

\[
\Pr[\text{there are at least } \alpha d/4 \text{ indices } i \in [d/2] \text{ such that } b_{\sigma(2i-1)} \neq b_{\sigma(2i)}] \geq 1 - c \cdot d \cdot 2^{-\alpha d/4}
\]

for some absolute constant \( c > 0 \).

We postpone the proof of lemma to Appendix A, and show now how it implies the claim. Indeed, by conditioning on \( \neg \text{Agr}_{d-r}(x) \), or equivalently on \( d/2 \leq q < d - r \), and choosing a random permutation of the indices \( i \in [d] \), by Lemma 9.6 we have

\[
\Pr[\exists \text{ at least } r/4 \text{ indices } i \in [d/2] \text{ such that } b_{\sigma(2i-1)} \neq b_{\sigma(2i)} \mid \neg \text{Agr}_{d-r}(x)] \geq 1 - c \cdot d \cdot 2^{-r/4},
\]

for some absolute constant \( c > 0 \), independent of \( d \) or \( r \). Therefore, \( \Pr[\text{PairAgr}_{q' \cdot \frac{1}{4}}(x) \mid \neg \text{Agr}_{d-r}(x)] \leq c \cdot d \cdot 2^{-r/4} \). \( \square \)
We use Claim 9.5 to complete the proof of Lemma 9.4. We have that
\[
1 - 2 \cdot e^{-(r/4d - \rho)^2 \min(t,d) \over s} \leq \Pr[\text{PairAgr}_{x,t} \cdot \zeta(x) \mid D_{t,\rho}(F(Q_t)) = 1] \\
\leq \Pr[\text{PairAgr}_{x,t} \cdot \zeta(x) \wedge \text{Agr}_{d-r}(x) \mid D_{t,\rho}(F(Q_t)) = 1] \\
+ \Pr[\text{PairAgr}_{x,t} \cdot \zeta(x) \wedge \neg \text{Agr}_{d-r}(x) \mid D_{t,\rho}(F(Q_t)) = 1] \\
\leq \Pr[\text{Agr}_{d-r}(x) \mid D_{t,\rho}(F(Q_t)) = 1] + \frac{\Pr[\text{PairAgr}_{x,t} \cdot \zeta(x) \mid \neg \text{Agr}_{d-r}(x)]}{\Pr[D_{t,\rho}(F(Q_t)) = 1]}.
\]
Therefore,
\[
\Pr[\text{Agr}_{d-r}(x) \mid D_{t,\rho}(F(Q_t)) = 1] \geq 1 - 2 \cdot e^{-(r/4d - \rho)^2 \min(t,d) \over s} \leq 2 \cdot e^{-(r/4d - \rho)^2 \min(t,d) \over s}
\]
for some absolute constant \(c \in \mathbb{R}\). This concludes the proof of Lemma 9.4.

9.2.2 Step 2: Individual indices behave linearly

We argue that, with high probability, \(b_x^{(i)} + b_y^{(i)} = b_{x+y}^{(i)}\) for most indices \(i \in [d]\).

**Lemma 9.7.** Let \(t, d, r \in \mathbb{N}\) and \(\rho \in [0, 1]\) be parameters such that \(1 \leq r \leq d/2\) and \(\rho < r/d\). Let \(F\) be a \(k\)-non-signaling function with \(k \geq 6d + 3t\). Fix \(x, y \in \{0, 1\}^n\), and let \(z_x^{(1)}, \ldots, z_x^{(d)}; z_y^{(1)}, \ldots, z_y^{(d)} \in \{0, 1\}^n\) be independent and uniformly distributed in \(\{0, 1\}^n\). For all \(i \in [d]\) define \(b_x^{(i)} = F(z_x^{(i)} + x) - F(z_x^{(i)})\), \(b_y^{(i)} = F(z_y^{(i)} + y) - F(z_y^{(i)})\), and \(b_{x+y}^{(i)} = F(z_{x+y}^{(i)} + x + y) - F(z_{x+y}^{(i)})\). Then
\[
\Pr[\exists \text{ more than } d - 2r \text{ indices } i \in [d] \text{ s.t. } b_x^{(i)} + b_y^{(i)} = b_{x+y}^{(i)} \mid D_{t,\rho}(F(Q_t)) = 1] \geq 1 - 2 \cdot e^{-(r/4d - \rho)^2 \min(t,d) \over s},
\]
where the probability is over the random \(z_x^{(i)}\)’s and \(z_y^{(i)}\)’s and the randomness of \(F\).

**Proof.** By Lemma 5.2 for uniformly random \(z_x^{(i)}; z_y^{(i)}\), and \(z_{x+y}^{(i)} = z_x^{(i)} + z_y^{(i)}\) it holds that
\[
\Pr[\text{for at least } d - r + 1 \text{ indices } i \in [d] \mid F(z_x^{(i)}) + F(z_y^{(i)}) = F(z_{x+y}^{(i)}) \mid D_{t,\rho}(F(Q_t)) = 1] \geq 1 - \frac{1}{\gamma} e^{-(r/4d - \rho)^2 \min(t,d) \over s}.
\]
Similarly, it holds that
\[
\Pr[\text{for at least } d - r + 1 \text{ indices } i \in [d] \mid F(z_x^{(i)} + x) + F(z_y^{(i)} + y) = F(z_{x+y}^{(i)} + x + y) \mid D_{t,\rho}(F(Q_t)) = 1] \geq 1 - \frac{1}{\gamma} e^{-(r/4d - \rho)^2 \min(t,d) \over s}.
\]
Lemma 9.7 follows by union bound over these two events.

9.2.3 Step 3: Putting it together

Finally, we combine Lemma 9.4 and Lemma 9.7 in order to prove Lemma 9.3.
We argue that Therefore, by the union bound we have This completes the proof of Lemma 9.3.

Suppose that Claim 9.8.

Proof. We show that for any event $E \subseteq \{0,1\}^S$ it holds that $\Pr{\mathcal{F}(S) \in E} \geq \Pr{\hat{\mathcal{F}}(S) \in E} - (5\rho + 3) |S| - \varepsilon$. This suffices because the same bound also holds for the complement event.
Consider the sampling procedure of $\mathcal{F}_{t,d}$ on the input $S$. For each $w \in S$ and $i \in [d]$ define $b^{(i)}_w = \mathcal{F}(z^{(i)}_w + w) - \mathcal{F}(z^{(i)}_w)$, where $z^{(i)}_w$’s are chosen according to the definition of $Enc_d$, and let $b_w = \max\{b^{(i)}_w : i \in [d]\}$.

By Lemma 9.4 for all $w \in S$ we have

$$\Pr[\text{Agr}_{d-r}(w) | D_{t,\rho}(\mathcal{F}(Q_t)) = 1] \geq 1 - \frac{2}{1 - \rho} \cdot e^{-\left(\frac{\epsilon^2}{4(d - \rho)^2} \right) \min\{t,d\}} - \frac{\epsilon \cdot d \cdot 2^{-r/4}}{1 - \rho}.$$  

Hence, by letting $r = 5pd$ and taking the union bound we get

$$\Pr[\text{Agr}_{(1-5\rho)d}(w) \forall w \in S | D_{t,\rho}(\mathcal{F}(Q_t)) = 1] \geq 1 - \epsilon'$$

for $\epsilon' = \frac{2|S|}{1-\rho} \cdot e^{-O(\rho^2 \min\{t,d\})} + \frac{c'd|S|2^{-5pd/4}}{1-\rho}$. In particular, this implies that

$$\Pr[\text{Agr}_{(1-5\rho)d}(w) \forall w \in S] \geq \Pr[\text{Agr}_{(1-5\rho)d}(w) \forall w \in S \land D_{t,\rho}(\mathcal{F}(Q_t)) = 1] \geq (1-\epsilon')(1-\rho) \geq 1-\epsilon'-\rho.$$  

Therefore,

$$\Pr[\mathcal{F}_{t,d}(S) \in E] = \Pr[(b_w : w \in S) \in E | D_{t,\rho}(\mathcal{F}(Q_t)) = 1]$$

$$\geq \Pr[(b_w : w \in S) \in E \land \text{Agr}_{(1-5\rho)d}(w) \forall w \in S | D_{t,\rho}(\mathcal{F}(Q_t)) = 1]$$

$$\geq \Pr[(b_w : w \in S) \in E \land \text{Agr}_{(1-5\rho)d}(w) \forall w \in S] - \rho.$$  

Observe $\tilde{\mathcal{F}}(S)$ can be sampled by choosing distinct $i_w \in [t]$ for each $w \in S$, and outputting $(b^{(i_w)}_w : w \in S)$. Note that if $S$ was a linearly independent set, then we do not need to require $i_w$’s to be distinct. However, since $i_w$’s are all distinct with probability $1 - \frac{|S|^2}{t}$, this is essentially immaterial. Therefore,

$$\Pr[\mathcal{F}_{t,d}(S) \in E] \geq \Pr[(b_w : w \in S) \in E \land \text{Agr}_{(1-5\rho)d}(w) \forall w \in S] - \rho$$

$$\geq \Pr[\tilde{\mathcal{F}}(S) \in E \land \text{Agr}_{(1-5\rho)d}(w) \forall w \in S] - 5\rho |S| - \frac{|S|^2}{t} - \epsilon' - 2\rho.$$  

and the claim follows for $\epsilon = \epsilon' + \frac{|S|^2}{t}$. \hfill $\Box$

### 9.4 Proof of Theorem 7

Let $k$ and $\bar{k} \in \mathbb{N}$, $\varepsilon \leq 1/400$, and let $\rho = \sqrt{\varepsilon} \leq 1/20$. Let $t = d = O\left(\frac{1}{\rho^2}(\bar{k} + \log \frac{1}{\rho})\right)$, and suppose that $k \geq 2d\bar{k} + 3t$, i.e., that $k = O\left(\frac{\bar{k}}{\varepsilon} \cdot (\bar{k} + \log \frac{1}{\varepsilon})\right)$. Let $\mathcal{F}$ be a $k$-non-signaling function, and suppose that $\Pr_{x,y,\mathcal{F}}[\mathcal{F}(x) + \mathcal{F}(y) = \mathcal{F}(x + y)] \geq 1 - \varepsilon$.

By Lemma 5.5 we deduce that $\Pr[|D_{t,\rho}(\mathcal{F}(Q_t)) = 1] \geq 1 - \rho$.

Since $k \geq 2d\bar{k} + 3t$, it follows by Lemma 9.1 that $\mathcal{F}_{t,d}$ is a $k$-non-signaling function.

Since $\rho \leq 1/20$ and $t = d = O\left(\frac{1}{\rho^2}(\bar{k} + \log \frac{1}{\rho})\right) \geq O\left(\bar{k} + \log \frac{1}{\rho} + \log \frac{1}{\rho^2}\right)$, and $\Pr[|D_{t,\rho}(\mathcal{F}(Q_t)) = 1] \geq 1 - \rho$, by Lemma 9.2 there exists a linear $k$-non-signaling function $\mathcal{L}$ such that $\Delta_k(\mathcal{L}, \mathcal{F}_{t,d}) < \rho/2$.  

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Since \( t = d = O(\frac{1}{\rho^n}(k + \log \frac{1}{\rho})) \), by Claim 9.8 for all sets \( S \subseteq \{0,1\}^n \) of size at most \( \bar{k} \) it holds that
\[
\left| \Pr[\mathcal{F}_{t,d}(S) \in E] - \Pr[\mathcal{F}(S) \in E] \right| \leq (6|S| + 2.5)\rho .
\]
Hence for all sets \( S \subseteq \{0,1\}^n \) of size at most \( \bar{k} \) and events \( E \subseteq \{0,1\}^S \) it holds that
\[
\left| \Pr[\mathcal{F}(S) \in E] - \Pr[\mathcal{L}(S) \in E] \right| \leq (6|S| + 3)\rho ,
\]
as required.
A Proof of Lemma 9.6

In the proof we will use Stirling’s formula to approximate the binomial coefficients involved.

**Proposition A.1** (Stirling’s approximation formula). There exists an absolute constant \( c_0 = \frac{3}{8} > 0 \) such that the following holds.

1. For all integers \( n > i \geq 1 \) it holds that \( c_0 \leq \frac{\binom{n}{i}}{\sqrt{i^{n-i} n^n}} \leq 1 \).

2. For all integers \( n, i, j \geq 1 \) such that \( n > i + j \) it holds that \( c_0 \leq \frac{\binom{n}{i,j,n-i-j}}{\sqrt{i,j,n-i-j} n^n} \leq 1 \).

We now prove Lemma 9.6. Let \( d \) be an even integer, and let \( b_1, \ldots, b_d \in \{0, 1\} \) be such that not all \( b_i \)'s are equal, and let \( \alpha = \min\{\frac{d}{4}, 1 - \frac{d}{4}\} \) for \( q = |\{i \in [d] : b_i = 0\}| \). For an integer \( 0 < r \leq (1 - \alpha)d \) denote by \( N_{d,q}(r) \) the number of permutations \( \sigma : [d] \rightarrow [d] \) such that there are exactly \( r \) indices \( i \in [d/2] \) such that \( b_{\alpha(d-i-1)} \neq b_{\alpha(2i)} \). Clearly, if \( r \neq q \mod 2 \), then \( N_{d,q}(r) = 0 \). If \( r \equiv q \mod 2 \), then

\[
N_{d,q}(r) = \binom{d/2}{r, q-r, d-q-r} \cdot 2^r \cdot q!(d-q)!
\]

where (i) the binomial term reflects the choice of which of the pairs of indices \( (2j - 1, 2j) \) will have only zeros, only ones, or both, (ii) the term \( 2^r \) accounts for the transpositions within the pairs with both a zero and a one, and (iii) \( q!(d-q)! \) counts separately the permutations of the zeros and the permutations of the ones in the positions specified by the previous two items.

For an integer \( 0 \leq r < \alpha d \) denote by \( \text{neq}_{d,q}(r) \) the event that there are exactly \( r \) indices \( i \in [d/2] \) such that \( b_{\alpha(2i-1)} \neq b_{\alpha(2i)} \), and let \( \beta = \frac{\alpha}{d} \in [0, \beta) \). Then

\[
\Pr[\text{neq}_{d,q}(r)] = 2^r \cdot \frac{N_{d,q}(r)}{d!} = 2^r \cdot \binom{d/2}{r, q-r, d-q-r} \binom{d}{q}
\]  

Next, we show that for all \( 0 \leq r < \frac{\alpha d}{4} \) it holds that \( \Pr[\text{neq}_{d,q}(r)] < c 2^{-\alpha d/4} \) for some absolute constant \( c > 0 \). The claim follows immediately as

\[
\Pr[\text{there are less than } \alpha d/4 \text{ indices } i \in [d/2] \text{ such that } b_{\alpha(2i-1)} \neq b_{\alpha(2i)}] = \sum_{0 \leq r < \frac{\alpha d}{4}} \Pr[\text{neq}_{d,q}(r)] \leq c d \cdot 2^{-\alpha d/2}
\]

as required.

**The case of \( r = 0 \):** For \( r = 0 \) and even \( q \) the expression in Eq. (9) reduces to

\[
\Pr[\text{neq}_{d,q}(r)] = \frac{q/2}{d} \left(\frac{q}{2}\right)^{d/2} \frac{2^{d/2} \cdot q^{d/2}}{q^{d/2}} = \frac{\sqrt{2}}{c_0} \cdot q^{d/2} (d-q)^{d/2}/d^{d/2}
\]

By substituting \( q = (1 - \alpha)d \) into the foregoing expression we get

\[
\Pr[\text{neq}_{d,q}(r)] \leq \frac{\sqrt{2}}{c_0} \cdot \left[\alpha^\alpha \cdot (1 - \alpha)^{1-\alpha}\right]^{d/2} \leq c \cdot 2^{-\alpha d/2}
\]
for the constant \( c = \frac{\sqrt{2}}{c_0} \).

**The case of \( r > 0 \):** By Stirling’s formula, the nominator in the expression in Eq. (9) is upper bounded by

\[
\left( r, \frac{d/2}{2}, \frac{d-q-r}{2} \right) \leq \sqrt{\frac{d/2}{\beta d \cdot \frac{\alpha - \beta}{2} \cdot \frac{1 - \alpha - \beta}{2}}} \cdot \left( \frac{(d/2)^{d/2}}{(\beta d)^{\beta d} \cdot \left( \frac{\alpha - \beta}{2} \cdot \frac{1 - \alpha - \beta}{2} \right)^{\frac{\alpha - \beta}{2} \cdot \frac{1 - \alpha - \beta}{2}}} \right)
\]

The denominator is lower bounded by

\[
\left( \frac{d}{q} \right) \geq c_0 \cdot \sqrt{\frac{d}{(1 - \alpha)d \cdot \alpha d}} \cdot \left( \frac{d^d}{((1 - \alpha)d)^{(1 - \alpha)d} \cdot (\alpha d)^{\alpha d}} \right)
\]

Therefore, we get the following bound on \( \Pr[\text{neq}_{d,q}(r)] \).

\[
\Pr[\text{neq}_{d,q}(r)] \leq \frac{1}{c_0} \cdot \sqrt{\frac{2\alpha(1 - \alpha)}{\beta(\alpha - \beta)(1 - \alpha - \beta)d^2} \cdot \left[ \frac{\alpha^\alpha \cdot (1 - \alpha)^{1 - \alpha}}{(2\beta)^{\beta \cdot (\alpha - \beta) \cdot \frac{\alpha - \beta}{2} \cdot (1 - \alpha - \beta) \cdot \frac{1 - \alpha - \beta}{2}}} \right]^d}
\]

(10)

Note that since \( r \) is an integer, we have \( \beta \in [1/d, \alpha/4] \), and the expression under the square root is upper bounded by

\[
\frac{2\alpha(1 - \alpha)}{(\beta d)^{\frac{3\alpha}{4} \cdot (1 - \frac{3\alpha}{4})}} = \frac{8}{3} \cdot \frac{1 - \alpha}{4 - 3\alpha} \leq \frac{8}{3} .
\]

Next, we upper bound the base of the exponent in last term. We first rewrite it follows.

\[
\frac{(1 - \alpha)^{1 - \alpha} \cdot \alpha^\alpha}{(2\beta)^{\beta \cdot (\alpha - \beta) \cdot \frac{\alpha - \beta}{2} \cdot (1 - \alpha - \beta) \cdot \frac{1 - \alpha - \beta}{2}}} = \left( \frac{\sqrt{(\alpha - \beta)(1 - \alpha - \beta)}}{2\beta} \right)^\beta \cdot \left( \frac{\alpha}{\sqrt{\alpha - \beta}} \right)^\alpha \cdot \left( \frac{1 - \alpha}{\sqrt{1 - \alpha - \beta}} \right)^{1 - \alpha} .
\]

For the first term note that \( \sqrt{(\alpha - \beta)(1 - \alpha - \beta)} \leq \sqrt{\alpha(1 - \alpha)} \leq 1/2 \), and hence

\[
\left( \frac{\sqrt{(\alpha - \beta)(1 - \alpha - \beta)}}{2\beta} \right)^\beta \leq \left( \frac{\sqrt{\alpha(1 - \alpha)}}{2\beta} \right)^\beta \leq \left( \frac{1}{4\beta} \right)^\beta \leq \left( \frac{1}{\alpha} \right)^{\alpha/4} ,
\]

where the last inequality follows from the fact that the function \( f(\beta) = (1/4\beta)^\beta \) is monotonically increasing in the interval \([0, 1/4e]\), and hence, for \( \alpha < 0.5 \) and \( \beta \in (0, \alpha/4) \) the maximum is obtained for \( \beta_{\text{max}} = \alpha/4 \), where \( f(\beta_{\text{max}}) = \alpha^{-\alpha/4} \).

For the second term, we have

\[
\left( \frac{\alpha}{\sqrt{\alpha - \beta}} \right)^\alpha \leq \alpha^{\alpha/4} .
\]
Indeed, the bound is equivalent to the inequality $\alpha^{3/4} < \sqrt{\alpha - \beta}$, which holds for all $\alpha < 0.5$ and $\beta \leq \alpha/4$.

For the third term, for all $\alpha \leq 0.5$ and $\beta \in [0, \alpha/4]$ we have

$$
\left( \frac{1 - \alpha}{\sqrt{1 - \alpha - \beta}} \right)^{1-\alpha} \leq \left( \frac{1 - \alpha}{\sqrt{1 - 5\alpha/4}} \right)^{1-\alpha} \leq 2^{-\alpha/4}.
$$

Multiplying the three terms, we get that the last term in Eq. (10) is upper bounded by $2^{-\alpha d/4}$. Therefore, for the constant $c = \frac{8}{3c_0}$ the expression in Eq. (10) satisfies

$$
\Pr[\text{neq}_{d,q}(r)] < c 2^{-\alpha d/4}
$$

for all $r < \alpha d/4$, and the claim follows.
B Compiling linear PCPs with sub-constant soundness error

In Section 8.4 we showed that the $t$-repeated ALMSS linear PCP verifier has soundness error $(\frac{n}{10})^t$ against linear $O(t)$-non-signaling functions when $t = \Omega(\log S)$. This, however, does not imply that the $t$-repeated ALMSS PCP verifier also has $\exp(-t)$ soundness error, because the compiler in Lemma 4.2 can at best achieve soundness error $1 - \frac{1}{400}$. We show a different way to compile linear PCPs into PCPs in a way that approximately preserves soundness error. This compilation, applied to the $t$-repeated ALMSS linear PCP, produces a PCP that still has $\exp(-t)$ soundness error.

**Lemma B.1.** For every $\varepsilon > 0$ and $d = O(k + \log \frac{1}{\varepsilon} + \log \frac{1}{\gamma + \varepsilon})$, it holds that

\[
\begin{align*}
\text{nsLPCP} & \quad \subseteq \text{nsPCP} \\
\text{query complexity: } & \quad q \\
\text{locality: } & \quad k \\
\text{decision predicate time: } & \quad T_1 \\
\text{query sampler time: } & \quad T_2
\end{align*}
\]

Applying the above lemma to the $t$-repeated ALMSS linear PCP, with $t = \Omega(\log S)$, we get

\[
\begin{align*}
\text{DSIZE}(S) & \subseteq \text{nsPCP} \\
\text{soundness error: } & \quad 2^{-O(t)} \\
\text{randomness: } & \quad O(t^2 S^2) \\
\text{proof length: } & \quad 2^{O(S^2)} \\
\text{query complexity: } & \quad O(t^2) \\
\text{locality: } & \quad O(t^2) \\
\text{query sampler time: } & \quad O(t^3 S^2) \\
\text{decision predicate time: } & \quad O(tn) + O(t^2)
\end{align*}
\]

**Proof of Lemma B.1.** The compiler uses $\mathcal{F}_{d,d}$ (see Section 9.1) rather than $\mathcal{F}$. Given $V_{\text{LPCP}} = (Q_{\text{LPCP}}, D_{\text{LPCP}})$, we define $V_{\text{PCP}} = (Q_{\text{PCP}}, D_{\text{PCP}})$ as follows. The query sampler $Q_{\text{PCP}}$ runs $Q_{\text{LPCP}}$ to obtain a set of queries $S_{\text{LPCP}} \subseteq \{0,1\}^\ell$ for the linear PCP, and outputs $\text{Enc}_d(S_{\text{LPCP}}) \cup Q_d$. Given the answers $\mathcal{F}(\text{Enc}_d(S_{\text{LPCP}})) \cup Q_d$, the decision predicate $D_{\text{PCP}}$ accepts if and only if $D_{d,0}(\mathcal{F}(Q_d)) = 1$ and $D_{\text{LPCP}}(\text{Dec}_d(\mathcal{F}(\text{Enc}_d(S_{\text{LPCP}})))) = 1$.

We only show soundness of $V_{\text{PCP}}$; the rest of the proof is straightforward.

Suppose that there exists a $(2k+3)d$-non-signaling function $\mathcal{F} : \{0,1\}^\ell \rightarrow \{0,1\}$ that is accepted by $V_{\text{PCP}}$ with probability greater than $\gamma + \varepsilon$, namely, that $\Pr[\mathcal{F} = 1] > \gamma + \varepsilon$. Observe that

\[
\begin{align*}
\Pr[V_{\text{PCP}} = 1] & = \Pr[D_{\text{LPCP}}(\text{Dec}_d(\mathcal{F}(\text{Enc}_d(Q_{\text{LPCP}})))) = 1 \land D_{d,0}(\mathcal{F}(Q_d)) = 1] \\
& = \Pr[D_{\text{LPCP}}(\text{Dec}_d(\mathcal{F}(\text{Enc}_d(Q_{\text{LPCP}})))) = 1 \mid D_{d,0}(\mathcal{F}(Q_d)) = 1 \mid \Pr[D_{d,0}(\mathcal{F}(Q_d)) = 1] \\
& = \Pr[D_{\text{LPCP}}(\mathcal{F}_{d,d}(Q_{\text{LPCP}})) = 1 \mid \Pr[D_{d,0}(\mathcal{F}(Q_d)) = 1] = \Pr[V_{\text{LPCP}}^{\mathcal{F}_{d,d}} = 1] \Pr[D_{d,0}(\mathcal{F}(Q_d)) = 1].
\end{align*}
\]

In particular, this implies that $\Pr[D_{d,0}(\mathcal{F}(Q_d)) = 1] \geq \gamma + \varepsilon$. Since $d = O(k + \log \frac{1}{\varepsilon} + \log \frac{1}{\gamma + \varepsilon})$, we can apply Lemma 9.2 to get a linear $k$-non-signaling function $L$ such that $\Delta_k(L, \mathcal{F}_{d,d}) \leq \varepsilon$. Therefore,

\[
\Pr[V_{\text{LPCP}}^L = 1] \geq \Pr[V_{\text{LPCP}}^{\mathcal{F}_{d,d}} = 1] - \varepsilon > (\gamma + \varepsilon) - \varepsilon = \gamma,
\]

which is a contradiction, and so the soundness of $V_{\text{LPCP}}$ implies soundness of $V_{\text{PCP}}$.  

\[
\text{Pr}[\mathcal{F} = 1] = \text{Pr}[V_{\text{PCP}} = 1] - \varepsilon < (\gamma + \varepsilon) - \varepsilon = \gamma.
\]


\[
\text{Pr}[\mathcal{F} = 1] = \text{Pr}[V_{\text{PCP}} = 1] - \varepsilon < (\gamma + \varepsilon) - \varepsilon = \gamma.
\]
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References


