# COMPRESSED SENSING WITH SEQUENTIAL OBSERVATIONS 

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#### Abstract

Compressed sensing allows perfect recovery of sparse signals (or signals sparse in some basis) using only a small number of measurements. The results in the literature have focused on the asymptotics of how many samples are required and the probability of making an error for a fixed batch of samples. We investigate an alternative scenario where observations are available in sequence and can be stopped as soon as there is reasonable certainty of correct reconstruction. This approach does not require knowing how sparse is the signal, and allows reconstruction using the smallest number of samples.

Central to our sequential approach is the stopping rule. For the random Gaussian ensemble we show that a simple stopping rule gives the absolute minimum number of observations required for exact recovery, with probability one. However, for other ensembles like Bernoulli or Fourier, this is no longer true, and the rule is modified to trade off delay in stopping and probability of error. We also consider near-sparse signals and describe how to estimate the reconstruction error from the sequence of solutions. This enables stopping once the error falls below a desired tolerance. Our sequential approach to compressed sensing involves a sequence of linear programs, and we outline how such a sequence can be solved efficiently.


Index Terms- Sequential compressed sensing

## 1. INTRODUCTION

In compressed sensing (CS) a few random linear measurements of a signal are taken, and the signal is recovered using sparse representation algorithms such as the $\ell_{1}$-based basis pursuit [1]. This is most useful when the cost of taking measurements is much larger than the computational overhead of recovering the signal, hence minimizing the number of required measurements is a primary concern. Existing analytical results provide guidelines on how many measurements are needed to ensure exact recovery with high probability, but these are often seen to be pessimistic $[1,2]$ and rely on $a$ priori knowledge about the sparsity of the unknown signal.

[^0]We consider an alternative scenario where one is able to get observations in sequence, and perform computations in between observations to decide whether enough samples have been obtained. Exact recovery is now possible from the smallest possible number of observations, and without any $a$-priori knowledge about how sparse the underlying signal is.

This, however, requires a computationally efficient decoder which can detect exactly when enough samples have been received. We first consider the case when noiseless measurements are taken using the random Gaussian ensemble, and we show that simply checking for one-step agreement yields such a decoder. After $M$ samples the decoder ${ }^{1}$ solves

$$
\begin{equation*}
\hat{\mathbf{x}}_{M}=\arg \min \|\mathbf{x}\|_{1} \quad \text { s.t. } \quad \mathbf{a}_{i}^{\prime} \mathbf{x}=y_{i}, \quad i=1, . ., M \tag{1}
\end{equation*}
$$

In case of one-step agreement, i.e. $\hat{\mathbf{x}}_{M+1}=\hat{\mathbf{x}}_{M}$, the decoder declares $\hat{\mathbf{x}}_{M}$ to be the reconstruction and stops requesting new measurements. In Section 2 we show that this decoder gives exact reconstruction with probability one.

For some other measurement ensembles, such as random Bernoulli and the ensemble of random rows from a Fourier basis, the one-step agreement stopping rule no longer has zero probability of error. We modify the rule to wait until $T$ subsequent solutions $\hat{\mathbf{x}}^{M}, \ldots, \hat{\mathbf{x}}^{M+T}$ all agree. In Section 3 we show that in the Bernoulli case the probability of making an error using this stopping rule decays exponentially with $T$, allowing trade-off of error probability and delay.

In Section 4 we show how the error in reconstruction can be estimated from the sequence of recovered solutions. This enables the decoder to stop once the error is below a required tolerance - even for signals that are not exactly sparse, but in which the energy is largely concentrated in a few components.

Finally, we propose an efficient way to solve the sequential problem in Section 5. Rather than re-solving the linear program from scratch after an additional measurement is received, we use an augmented linear program that uses the solution at step $M$ to guide its search for the new solution. We show empirically that this approach significantly reduces computational complexity.

[^1]

Fig. 1. A new constraint is added: $\mathbf{a}_{M+1}^{\prime} \mathbf{x}=y_{M+1}$. Probability that this hyperplane passing through $\mathrm{x}^{*}$ also passes through $\hat{\mathbf{x}}^{M}$ is zero.

## 2. STOPPING RULE IN THE GAUSSIAN CASE

We now analyze the sequential CS approach with the Gaussian measurement ensemble. Suppose that the underlying sparse signal $\mathbf{x}^{*} \in \mathbb{R}^{N}$ has $K$ non-zero components (we denote the number of non-zero entries in $\mathbf{x}$ by $\|\mathrm{x}\|_{0}$ ). We sequentially receive random measurements $y_{i}=\mathbf{a}_{i}^{\prime} \mathbf{x}^{*}$, where $\mathbf{a}_{i} \sim \mathcal{N}(0, I)$ is a $N$-vector of i.i.d. Gaussian samples. At step $M$ we solve the basis-pursuit problem in (1) using all the received data. Results in compressed sensing [1, 2] indicate that after receiving around $M \propto K \log (N)$ measurements, solving (1) recovers the signal $\mathbf{x}^{*}$ with high probability. This requires the knowledge of $K$, which may not be available, and only rough bounds on the scaling constants are known. Our approach is different - we compare the solutions at step $M$ and $M+1$, and if they agree, we declare correct recovery.

Proposition 1 If $\hat{\mathbf{x}}^{M+1}=\hat{\mathbf{x}}^{M}$ in the Gaussian measurement ensemble, then $\hat{\mathbf{x}}^{M}=\mathbf{x}^{*}$, with probability 1 .

Proof. Let $\mathbf{y}_{1: M} \triangleq\left[y_{1}, \ldots, y_{M}\right]^{\prime}$, and $A^{M} \triangleq\left[\mathbf{a}_{1}, \ldots, \mathbf{a}_{M}\right]^{\prime}$. Suppose that $\hat{\mathbf{x}}^{M} \neq \mathbf{x}^{*}$. We have that $\mathbf{y}_{1: M}=A^{M} \hat{\mathbf{x}}^{M}$ and $\mathbf{y}_{1: M}=A^{M} \mathbf{x}^{*}$ : both $\mathbf{x}^{*}$ and $\hat{\mathbf{x}}^{M}$ belong to the $(N-M)$ dimensional affine space $\left\{\mathbf{x} \mid \mathbf{y}_{1: M}=A^{M} \mathbf{x}\right\}$. The next measurement passes a random hyperplane $y_{M+1}=\mathbf{a}_{M+1}^{\prime} \mathbf{x}^{*}$ through $x^{*}$ and reduces the dimension of the affine subspace of feasible solutions by 1 . In order for $\hat{\mathbf{x}}^{M}$ to remain feasible at step $M+1$, it must hold that $\mathbf{y}_{M+1}=\mathbf{a}_{M+1}^{\prime} \hat{\mathbf{x}}^{M}$. Since we also have $\mathbf{y}_{M+1}=\mathbf{a}_{M+1}^{\prime} \mathbf{x}^{*}$, then $\hat{\mathbf{x}}^{M}$ remains feasible only if $\left(\hat{\mathbf{x}}^{M}-\mathbf{x}^{*}\right)^{\prime} \mathbf{a}_{M+1}=0$, i.e. if $\mathbf{a}_{M+1}$ falls in the $N-1$ dimensional subspace of $\mathbb{R}^{N}$ corresponding to $\operatorname{Null}\left(\left(\hat{\mathbf{x}}^{M}-\mathbf{x}^{*}\right)^{\prime}\right)$. As $\mathbf{a}_{M+1}$ is random and independent of $\hat{\mathbf{x}}^{M}$ and of the previous samples $\mathbf{a}_{1}, \ldots, \mathbf{a}_{M}$, the probability that this happens is 0 (event with measure zero). See Figure 1 for illustration.

Clearly, if we obtain $\hat{\mathbf{x}}^{M}=\mathrm{x}^{*}$, then the solution will not change with additional samples: $\mathbf{x}^{*}$ is always in the feasible set, and the feasible set is shrinking with each new sample. In the Gaussian case the stopping rule can be simplified further: if $\hat{\mathbf{x}}^{M}$ has fewer than $M$ non-zero entries, then $\hat{\mathbf{x}}^{M}=\mathbf{x}^{*}$.


Fig. 2. Gaussian ensemble example: $N=100$, and $K=10$. (Top): $\left\|\hat{\mathbf{x}}^{M}\right\|_{0}$. (Middle): $\left\|\hat{\mathbf{x}}^{M}\right\|_{1}$. (Bottom): $\left\|\mathbf{x}^{*}-\hat{\mathbf{x}}^{M}\right\|_{2}$.

Proposition 2 If $\left\|\hat{\mathbf{x}}^{M}\right\|_{0}<M$, then $\hat{\mathbf{x}}^{M}=\mathbf{x}^{*}$ with prob. $1^{2}$.
Proof. Let $A=A^{M}$ to simplify notation. Then $A$ is $M \times N$, with $M<N$. The key fact about random Gaussian matrices is that any $M \times M$ submatrix of $A$ is non-singular with probability $1^{3}$. Let $\mathcal{I}$ be the support of $\mathbf{x}^{*}$, i.e. $\mathcal{I}=$ $\left\{i \mid x_{i}^{*} \neq 0\right\}$. Suppose that there is another sparse feasible vector $\hat{\mathbf{x}} \neq \mathbf{x}^{*}$ with support $\mathcal{J}$, such that $|\mathcal{J}|<M$. There are two possibilities: $\mathcal{I} \subset \mathcal{J}$ and $\mathcal{I} \neq \mathcal{I} \cap \mathcal{J}$. We show that in both cases $\hat{\mathbf{x}} \neq \mathbf{x}^{*}$ can occur only with probability zero.

First suppose $\mathcal{I} \subset \mathcal{J}$. Then $\hat{\mathbf{x}}-\mathbf{x}^{*} \in \operatorname{Null}(A)$, and support of $\hat{\mathbf{x}}-\mathbf{x}^{*}$ is a subset of $\mathcal{J}$, hence it is smaller than $M$. But that means that fewer than $M$ columns of $A$ are linearly dependent, which only happens with probability zero.

Now consider the case $\mathcal{I} \neq \mathcal{I} \cap \mathcal{J}$. First fix some such set $\mathcal{J}$. We use the notation $\mathcal{I} \backslash \mathcal{J}=\{i \in I \mid i \notin J\}$. The probability that the vector $\tilde{\mathbf{y}}=A_{\mathcal{I} \backslash \mathcal{J}} \mathbf{x}_{\mathcal{I} \backslash \mathcal{J}}^{*}$ falls into $\operatorname{span}\left(A_{\mathcal{J}}\right)$ is zero, as $|J|<M$ and the elements of $A_{\mathcal{I} \backslash \mathcal{J}}$ are independent of $A_{\mathcal{J}}$ and $\mathbf{x}^{*}$. The number of such possible subsets $\mathcal{J}$ is finite (albeit large), so the event that $\tilde{\mathbf{y}}$ falls into span of any such $A_{\mathcal{J}}$ still has probability zero. Hence, with probability 1 there is only one solution with $\|\mathbf{x}\|_{0}<M$, namely $\mathbf{x}^{*}$.

Consider an example in Figure 2 with $N=100$, and $K=$ 10. We keep solving (1) until agreement, $\hat{\mathbf{x}}^{M}=\hat{\mathbf{x}}^{M+1}$. The top plot shows that $\left\|\hat{\mathbf{x}}^{M}\right\|_{0}$ increases linearly with $M$ until $M=35$, at which point it drops to $K=10$ and we have $\hat{\mathbf{x}}^{M}=\mathbf{x}^{*}$. The middle plot shows the monotonic increase in $\left\|\hat{\mathbf{x}}^{M}\right\|_{1}$ (as the feasible set is shrinking with $M$ ). The bottom plot shows the error-norm of the solution, $\left\|\hat{\mathbf{x}}^{M}-\mathbf{x}^{*}\right\|_{2}$. On average it tends to go down with more observations, but nonmonotonically. After $M=35$ the error becomes zero.

[^2]
## 3. STOPPING RULE IN THE BERNOULLI CASE

Now suppose that the measurement vectors $\mathbf{a}_{i}$ have equiprobable i.i.d. Bernoulli entries $\pm 1$. A difference emerges from the Gaussian case: the probability that all $M \times M$ submatrices of $A^{M}$ are non-singular is no longer 0 . This makes it possible (with non-zero probability) for $\hat{\mathbf{x}}^{M+1}$ to agree with $\hat{\mathbf{x}}^{M}$ when $\hat{\mathbf{x}}^{M} \neq \mathbf{x}^{*}$, and for erroneous solutions $\hat{\mathbf{x}}^{M}$ to have cardinality less than $M$. We modify the stopping rule to require agreement for several steps - success is declared only when last $T$ solutions all agree. We show in proposition 3 that the probability of error decays exponentially with $T$. We use the following Lemma from [3]:
Lemma 1 Let a be an i.i.d. equiprobable Bernoulli vector with $\mathbf{a} \in\{-1,1\}^{N}$. Let $W$ be a deterministic d-dimensional subspace of $\mathbb{R}^{N}, 0 \leq d<N$. Then $P(\mathbf{a} \in W) \leq 2^{d-N}$.

We are now ready to establish the following claim:
Proposition 3 Consider the Bernoulli measurement case. If $\hat{\mathbf{x}}^{M}=\hat{\mathbf{x}}^{M+1}=\ldots=\hat{\mathbf{x}}^{M+T}$, then $\hat{\mathbf{x}}^{M}=\mathbf{x}^{*}$ with probability greater than or equal to $1-2^{-T}$.

Proof. Suppose that $\hat{\mathbf{x}}^{M} \neq \mathbf{x}^{*}$. Denote the support of $\mathbf{x}^{*}$ by $\mathcal{I}$ and the support of $\hat{\mathbf{x}}^{M}$ by $\mathcal{J}$. At step $M$ we have $A^{M} \mathbf{x}^{*}=A^{M} \hat{\mathbf{x}}^{M}$. Let $W=\left\{\mathbf{a} \mid\left(\hat{\mathbf{x}}^{M}-\mathbf{x}^{*}\right)^{\prime} \mathbf{a}=0\right\}$, i.e. the nullspace of $\left(\hat{\mathbf{x}}^{M}-\mathbf{x}^{*}\right)^{\prime}$. Then $W$ is an $(N-1)$-dimensional subspace of $\mathbb{R}^{N}$.

Given a new random Bernoulli sample $\mathbf{a}_{M+1}$, the vector $\hat{\mathbf{x}}^{M}$ can remain feasible at step $M+1$ only if ( $\hat{\mathbf{x}}^{M}-$ $\left.\mathbf{x}^{*}\right)^{\prime} \mathbf{a}_{M+1}=0$, i.e. if $\mathbf{a}_{M+1}$ falls into $W$. By Lemma 1, the probability that $\mathbf{a}_{M+1} \in W$ is a most $1 / 2$. The same argument applies to all subsequent samples of $\mathbf{a}_{M+i}$ for $i=$ $1, . ., T$, so the probability of having $T$-step agreement with an incorrect solution is bounded above by $2^{-T}$.

We now pursue an alternative heuristic analysis, more akin to Proposition 2. For the Bernoulli case, $\left\|\hat{\mathbf{x}}^{M}\right\|_{0}<M$ does not imply $\hat{\mathbf{x}}^{M}=\mathbf{x}^{*}$. However, we believe that if $N^{2} 2^{1-M}$ $\ll 1$, then $\hat{\mathbf{x}}^{M}=\mathbf{x}^{*}$ with high probability. Since the elements of $\mathbf{a}_{i}$ belong to finite set $\{-1,1\}$, an $M \times M$ submatrix of $A^{M}$ can be singular with non-zero probability. Surprisingly, characterizing this probability is a very hard question. It is conjectured [3] that the dominant source of singularity is the event that two columns or two rows are equal or opposite in sign. This leads to the following estimate (here $X_{M}$ is $M \times M$ ): ${ }^{4}$

$$
\begin{equation*}
P\left(\operatorname{det} X_{M}=0\right)=(1+o(1)) M^{2} 2^{1-M} \tag{2}
\end{equation*}
$$

However the very recent best provable bound on this probability is still rather far: $P\left(\operatorname{det} X_{M}=0\right)=\left(\left(\frac{3}{4}+o(1)\right)^{M}\right)$ [3]. If we assume that the simple estimate based on pairs of columns is accurate, similar analysis shows that the probability that a random $\pm 1 M \times N$ matrix with $M \ll N$ having all $M \times M$ submatrices non-singular is $(1+o(1)) N^{2} 2^{1-M}$.

[^3]
## 4. NEAR-SPARSE SIGNALS

In practical settings, e.g. when taking Fourier and wavelet transforms of smooth signals, we may only have approximate sparseness: a few values are large, and most are very small. The stopping rule from Section 2 is vacuous for near-sparse signals, as $\left\|\mathbf{x}^{*}\right\|_{0}=N$ and all samples are needed for perfect recovery. We consider Gaussian measurements $A^{M}$ for this section, and show that the reconstruction error satisfies

$$
\begin{equation*}
\left\|\mathbf{x}^{*}-\hat{\mathbf{x}}^{M}\right\|_{2} \leq C_{T}\left\|\hat{\mathbf{x}}^{M}-\hat{\mathbf{x}}^{M+T}\right\|_{2} \tag{3}
\end{equation*}
$$

where $C_{T}$ is a random variable. We characterize $E\left[C_{T}\right]$ and $\operatorname{Var}\left[C_{T}\right]$ - this gives us a confidence interval on the reconstruction error using the observed change $\left\|\hat{\mathbf{x}}^{M}-\hat{\mathbf{x}}^{M+T}\right\|_{2}$. We can now stop taking new measurements once the error falls below a desired tolerance. Note that our analysis does not assume a model of decay, and bounds the reconstruction error by comparing subsequent solutions. In contrast, related results in CS literature assume a power-law decay of entries of $\mathbf{x}^{*}$ (upon sorting) and show that with roughly $O(K \log N)$ samples, $\hat{\mathbf{x}}^{M}$ in (1) will have similar error to that of keeping the $K$ largest entries in $\mathrm{x}^{*}[1]$.

We now outline the analysis in (3). Consider Figure 1 again. Solution $\hat{\mathbf{x}}^{M+T}$ lies on the hyperplane $H_{M+T} \triangleq\left\{\mathbf{x} \mid y_{i}=\right.$ $\left.\mathbf{a}_{i}^{\prime} \mathbf{x}, i=1, . ., M+T\right\}$. Let $\theta_{T}$ be the angle between the line connecting $\mathbf{x}^{*}$ with $\hat{\mathbf{x}}^{M}$, and $H_{M+T}$. Denote distance to $H_{M+T}$ by $d\left(\hat{\mathbf{x}}^{M}, H_{M+T}\right)$. Then

$$
\begin{equation*}
d\left(\mathbf{x}^{*}, \hat{\mathbf{x}}^{M}\right)=\frac{d\left(\hat{\mathbf{x}}^{M}, H_{M+T}\right)}{\sin \left(\theta_{T}\right)} \leq \frac{d\left(\hat{\mathbf{x}}^{M+T}, \hat{\mathbf{x}}^{M}\right)}{\sin \left(\theta_{T}\right)} \tag{4}
\end{equation*}
$$

$\theta_{T}$ is a random variable - the angle between a fixed vector in $\mathbb{R}^{N-M}$ and a random $N-(M+T)$ dimensional hyperplane. The constant $C_{T}$ in (3) is equal to $\frac{1}{\sin \left(\theta_{T}\right)}$. We next analyze the distribution of $\theta_{T}$ and hence of $C_{T}$.

Let $L=N-M$. In the Gaussian case (due to invariance to orthogonal transformations) it is equivalent to consider the angle $\theta$ between a fixed $(L-T)$-dimensional subspace $H$ and a random vector $\mathbf{h}$ in $\mathbb{R}^{L}$. Let $H$ be the span of the last $L-T$ coordinate vectors, and $\mathbf{h}$ be i.i.d. Gaussian. Then: $C_{T}=\frac{1}{\sin (\theta)}=\sqrt{\sum_{i=1}^{L} x_{i}^{2}} / \sqrt{\sum_{i=1}^{T} x_{i}^{2}}$.

Using the properties of $\chi_{L}, \chi_{L}^{2}$, and inverse- $\chi_{L}^{2}$ distributions [4] and Jensen's inequality, we have $E\left[C_{T}\right] \approx \sqrt{\frac{L}{T}}$, and an upper bound on the variance: ${ }^{5}$

$$
\begin{equation*}
\operatorname{Var}\left[C_{T}\right] \leq \frac{L-2}{T-2}-\frac{L}{T} \tag{5}
\end{equation*}
$$

In Figure 3 (left) we plot the mean estimate and the standard deviation bound for $L=100$ and a range of $T$. We compare them to sample mean and standard deviation of $C_{T}$ based

$$
\begin{aligned}
& { }^{5} \text { Consider } E\left[\sin (\theta)^{2}\right]=\left(\sum_{i=1}^{T} x_{i}^{2}\right) /\|\mathbf{x}\|_{2}^{2} \text {. We have } E\left[\frac{x_{i}^{2}}{\|\mathbf{x}\|_{2}^{2}}\right]=\frac{1}{L} \\
& \text { (Dirichlet dist.), so } E\left[\sin (\theta)^{2}\right]=\frac{T}{L} \text {. Using Jensen's ineq. with } \sqrt{1 / x} \text {, } \\
& E[1 / \sin (\theta)] \geq \sqrt{\frac{L}{T}} \text {. Finally, } E\left[\frac{1}{\sin (\theta)^{2}}\right]=\frac{L-2}{T-2} \text { (for } T>2 \text { ). }
\end{aligned}
$$



Fig. 3. (Top left) Mean estimate and (bottom left) standard deviation bound of $C_{T}$ vs. averages over 2500 samples. $L=100$. (Top right): Estimated and actual errors: power-law decay, and (bottom right) blocky signals. $N=80, T=3$.
on 2500 samples. They give very good approximation for most of the range of $T>2$. Standard deviation quickly falls off with $T$, giving tight confidence intervals (by Chebyshev ineq. $\left.p\left(|a-E[a]| \geq k \sigma_{a}\right) \leq \frac{1}{k^{2}}\right)$. On the right plots we show that predicted errors (Chebyshev-bounds) follow closely the actual errors for certain near-sparse signals.

## 5. EFFICIENT SEQUENTIAL SOLUTION

The main motivation for the sequential approach is to reduce the number of measurements to as few as possible. Yet, we would also like to keep the computational complexity of the sequential approach low. Instead of re-solving the linear program (1) after each new sample, we would like to use the solution to the previous problem to guide the current problem. We now investigate a linear programming approach to accomplish this. In related work, [5] proposed to use Row-action methods for compressed sensing, which rely on a quadratic programming formulation equivalent to (1) and can take advantage of sequential measurements.

We can not use the solution $\hat{\mathbf{x}}^{M}$ directly as a starting point for the new problem at step $M+1$, because in general it will not be feasible. In the Gaussian measurement case, unless $\hat{\mathbf{x}}^{M}=\mathbf{x}^{*}$, the new constraint $\mathbf{a}_{M+1}^{\prime} \hat{\mathbf{x}}^{M}=y_{M+1}$ will be violated. One way to handle this is through a dual formulation, but we instead use an augmented primal formulation [6].

First, to model (1) as a linear program we use the standard trick: define $x_{i}^{+}=\max \left(x_{i}, 0\right), x_{i}^{-}=\max \left(-x_{i}, 0\right)$, and $\mathbf{x}=$ $\mathbf{x}^{+}-\mathbf{x}^{-}$. This gives a linear program in standard form:

$$
\begin{gather*}
\min \mathbf{1}^{\prime} \mathbf{x}^{+}+\mathbf{1}^{\prime} \mathbf{x}^{-}  \tag{6}\\
\mathbf{y}_{1: M}=\left[A^{M}-A^{M}\right]\left[\begin{array}{l}
\mathbf{x}^{+} \\
\mathbf{x}^{-}
\end{array}\right], \quad \text { and } \mathbf{x}^{+}, \mathbf{x}^{-} \geq 0
\end{gather*}
$$

Next we need to add an extra constraint $y_{M+1}=\mathbf{a}_{M+1}^{\prime} \mathbf{x}^{+}-$ $\mathbf{a}_{M+1}^{\prime} \mathbf{x}^{-}$. Suppose that $\mathbf{a}_{M+1}^{\prime} \hat{\mathbf{x}}^{M}>y_{M+1}$. We add an extra slack variable $z$ to the linear program, and a high positive cost


Fig. 4. A comparison of the number of simplex iterations when solving (1) from scratch (LP1) and using the solution at step $M-1$ (LP2). Plot shows \# iter. vs. $M$, over 100 trials.
$Q$ on $z$. This gives the following linear program:

$$
\begin{gather*}
\min \mathbf{1}^{\prime} \mathbf{x}^{+}+\mathbf{1}^{\prime} \mathbf{x}^{-}+Q z  \tag{7}\\
\mathbf{y}_{1: M}=\left[A^{M}-A^{M}\right]\left[\begin{array}{l}
\mathbf{x}^{+} \\
\mathbf{x}^{-}
\end{array}\right], \quad \text { and } \mathbf{x}^{+}, \mathbf{x}^{-} \geq 0 \\
y_{M+1}=\mathbf{a}_{M+1}^{\prime} \mathbf{x}^{+}-\mathbf{a}_{M+1}^{\prime} \mathbf{x}^{-}-z, \quad \text { and } z \geq 0
\end{gather*}
$$

Now using $\hat{\mathbf{x}}^{M}$ and $z=\mathbf{a}_{M+1}^{\prime}\left(\hat{\mathbf{x}}^{M}\right)^{+}-\mathbf{a}_{M+1}^{\prime}\left(\hat{\mathbf{x}}^{M}\right)^{-}-$ $y_{M+1}$ yields a basic feasible solution to this augmented problem. By selecting $Q$ large enough, $z$ will be removed from the optimal basis (i.e. $z$ is set to 0 ), and the solutions to this problem and the $(M+1)$-th sequential problem are the same.

We test the approach on an example with $N=200, K=$ 10, and 100 trials. In Figure 4 we plot the number of iterations of the simplex method required to solve the problem (1) at step $M$ from scratch (LP1) and using the formulation in (7) (LP2). To solve (6) we first have to find a basic feasible solution (phase 1) and then move from it to the optimal BFS. An important advantage of (7) is that we start with a basic feasible solution, so phase 1 is not required. The figure illustrates that for large $M$ the approach LP2 is significantly faster.

## 6. REFERENCES

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[^0]:    This research was supported by the Air Force Office of Scientific Research under Grant FA9550-04-1-0351, the Army Research Office under MURI Grant W911NF-06-1-0076.

[^1]:    ${ }^{1}$ Here we use the basis pursuit decoder, but our results apply to any sparse decoder, e.g. brute-force decoder and greedy matching pursuit.

[^2]:    ${ }^{2}$ Note that a random measurement model is essential: for a fixed matrix $A$ if $2 K>M$ then there exist $\mathbf{x}_{1}$ and $\mathbf{x}_{2}$ such that $A \mathbf{x}_{1}=A \mathbf{x}_{2}$ and $\left\|\mathbf{x}_{i}\right\|_{0} \leq K$. However, for a fixed $\mathbf{x}^{*}$ with $\left\|\mathbf{x}^{*}\right\|_{0}<M$ the probability that it will have ambiguous sparse solutions for a random choice of $A$ is zero.
    ${ }^{3}$ This is easy to see: fix $T \subset\{1, \ldots, N\}$ with $|T|=M$. Then probability that $A_{T_{M}} \in \operatorname{span}\left(A_{T_{1}}, \ldots, A_{T_{M-1}}\right)$ is zero, as $A_{T_{M}}$ is a random vector in $\mathbb{R}^{M}$ and the remaining columns span a lower-dimensional subspace.

[^3]:    ${ }^{4}$ Probability that two columns are equal or opposite in sign is $2{ }^{1-M}$, and there are $O\left(M^{2}\right)$ pairs of columns.

