### Randomized Coordinate Descent for Big Data Optimization (Theory)

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## Lecture 1 NSync



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#### The Problem

In order to quickly illustrate the important topics and notions that we will study in more depth later, we first consider the following simple problem:

minimize 
$$f(x)$$
 (1) subject to  $x \in \mathbb{R}^n$ 

We will assume that f is:

- "smooth" (will be made precise later)
- strongly convex



#### Introduction to Parallel Coordinate Descent

This **NSync algorithm** was introduced in a brief 5p paper by R. and Takáč [11] and was meant to be an entry point to the field of parallel coordinate descent.

#### Algorithm (NSync)

**Input:** initial point  $x_0 \in \mathbb{R}^n$ 

subset probabilities  $\{p_S\}$  for each  $S \subseteq [n] \stackrel{\text{def}}{=} \{1, 2, \dots, n\}$  stepsize parameters  $v_1, \dots, v_n > 0$ 

for k = 0, 1, 2, ... do

a) Select a random set of coordinates  $S_k \subseteq [n]$  following the law

$$P(S_k = S) = p_S, \qquad S \subseteq [n]$$

b) Update (possibly in parallel) selected coordinates:

$$x_{k+1} = x_k - \sum_{i \in S_k} \frac{1}{v_i} (e_i^T \nabla f(x_k)) e_i$$



end for

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#### Two More Ways of Writing the Update Step

1. Coordinate-by-coordinate:

$$x_{k+1}^{(i)} = \begin{cases} x_k^{(i)}, & i \notin S_k, \\ x_k^{(i)} - \frac{1}{v_i} (\nabla f(x_k))^{(i)}, & i \in S_k. \end{cases}$$

2. Via projection to a subset of blocks: If for  $h \in \mathbb{R}^n$  and  $S \subseteq [n]$  we write

$$h_{[S]} \stackrel{\mathsf{def}}{=} \sum_{i \in S} h^{(i)} e_i,$$

then

$$x_{k+1} = x_k + h_{[S_k]}$$
 for  $h = -(\operatorname{Diag}(v))^{-1} \nabla f(x_k)$ .

We shall interchangeably write:

$$\nabla_i f(x) = e_i^T \nabla f(x) = (\nabla f(x))^{(i)}$$



#### Samplings

#### Definition 1 (Sampling)

By the name **sampling** we will refer to a set valued random mapping with values being subsets of  $[n] = \{1, 2, ..., n\}$ . For sampling  $\hat{S}$  we define  $p = (p_1, ..., p_n)^T$ , where

$$p_i = \mathbf{P}(i \in \hat{S}) \tag{2}$$

We say that  $\hat{S}$  is **proper**, if  $p_i > 0$  for all i.

Lemma 2 ([5])

$$\sum_{i=1}^{n} p_i = \mathbf{E}[|\hat{S}|]. \tag{3}$$

Proof.

$$\sum_{i=1}^{n} p_{i} \stackrel{(2)}{=} \sum_{i=1}^{n} \sum_{S \subseteq [n]: i \in S} p_{S} = \sum_{S \subseteq [n]} \sum_{i: i \in S} p_{S} = \sum_{S \subseteq [n]} p_{S} |S| = \mathbf{E}[|\hat{S}|].$$



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#### Assumption: Strong convexity

#### Assumption 1 (Strong convexity)

f is differentiable and  $\gamma$ -strongly convex with respect to the norm  $\|\cdot\|_s$  (weighted Euclidean norm with weights  $s=(s_1,\ldots,s_n)^T>0$ ). That is, for all  $x,h\in\mathbb{R}^n$ ,

$$f(x+h) \ge f(x) + \langle \nabla f(x), h \rangle + \frac{\gamma}{2} ||h||_s^2. \tag{4}$$

Notation used above:

$$||h||_s \stackrel{\text{def}}{=} \left(\sum_{i=1}^n s_i(h^{(i)})^2\right)^{1/2}$$
 (weighted Euclidean norm)



## Assumption: Expected Separable Overapproximation Assumption 2 (ESO)

Assume  $\hat{S}$  is proper and that for some vector of positive weights  $v = (v_1, \dots, v_n)$  and all  $x, h \in \mathbb{R}^n$ ,

$$\mathbf{E}[f(x+h_{[\hat{S}]})] \le f(x) + \langle \nabla f(x), h \rangle_p + \frac{1}{2} \|h\|_{p \bullet_V}^2.$$
 (5)

For simplicity, we will often write

$$(f,\hat{S}) \sim ESO(v).$$

Note that the ESO parameters v, p depend on both f and  $\hat{S}$ .

Notation used above:

$$h_{[S]} \stackrel{\text{def}}{=} \sum_{i \in S} h^{(i)} e_i$$
 (projection of  $h \in \mathbb{R}^n$  onto coordinates  $i \in S$ )

$$\langle g, h \rangle_p \stackrel{\text{def}}{=} \sum_{i=1}^n p_i g^{(i)} h^{(i)}$$
 (weighted inner product)

$$p \bullet v \stackrel{\text{def}}{=} (p^{(1)}v^{(1)}, \dots, p^{(n)}v^{(n)})$$
 (Hadamard product)



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#### Complexity of NSync

#### Theorem 3 ([11])

Let  $x_*$  be a minimizer of f. Let Assumptions 1 and 2 be satisfied for a proper sampling  $\hat{S}$  (that is,  $(f, \hat{S}) \sim ESO(v)$ ). Choose

- starting point  $x_0 \in \mathbb{R}^n$ ,
- error tolerance  $0 < \epsilon < f(x_0) f(x_*)$  and
- confidence level  $0 < \rho < 1$ .

If  $\{x_k\}$  are the random iterates generated by NSync where the random sets  $S_k$  are iid following the distribution of  $\hat{S}$ , then

$$\mathbf{K} \geq \frac{\Lambda}{\gamma} \log \left( \frac{f(x_0) - f(x_*)}{\epsilon \rho} \right) \Rightarrow \mathbf{P}(f(x_{\mathbf{K}}) - f(x_*) \leq \epsilon) \geq 1 - \rho, \quad (6)$$

where

$$\Lambda \stackrel{def}{=} \max_{i=1,\ldots,n} \frac{v_i}{p_i s_i} \geq \frac{\sum_{i=1}^n \frac{v_i}{s_i}}{\mathbf{E}[|\hat{S}|]}. \tag{7}$$



#### What does this mean?

- ▶ Linear convergence. NSync converges linearly (i.e., logarithmic dependence on  $\epsilon$ )
- ▶ High confidence is not a problem.  $\rho$  appears inside the logarithm, so it easy to achieve high confidence (by running the method longer; there is no need to restart)
- ▶ Focus on the leading term. The leading term is  $\Lambda$ ; and we have closed from expression for it in terms of
  - ▶ parameters  $v_1, ..., v_n$  (which depend on f and  $\hat{S}$ )
  - ▶ parameters  $p_1, ..., p_n$  (which depend on  $\hat{S}$ )
- ▶ Parallelization speedup. The lower bound suggests that if it was the case that the parameters  $v_i$  did not grow with increasing  $\tau \stackrel{\text{def}}{=} \mathbf{E}[|\hat{S}|]$ , then we could potentially be getting linear speedup in  $\tau$  (average number of updates per iteration).
  - So we shall study the dependence of  $v_i$  on  $\tau$  ( this will depend on f and  $\hat{S}$ )
  - As we shall see, speedup does is often guaranteed for sparse problems.





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#### Proof of Theorem 3 - Part I

▶ If we let  $\mu \stackrel{\text{def}}{=} \gamma/\Lambda$ , then

$$f(x+h) \stackrel{(4)}{\geq} f(x) + \langle \nabla f(x), h \rangle + \frac{\gamma}{2} \|h\|_{s}^{2}$$

$$\geq f(x) + \langle \nabla f(x), h \rangle + \frac{\mu}{2} \|h\|_{v \bullet p^{-1}}^{2}. \tag{8}$$

Indeed,  $\mu$  is defined to be the largest number for which  $\gamma \|h\|_s^2 \ge \mu \|h\|_{v \bullet p^{-1}}^2$  holds for all h. Hence, f is  $\mu$ -strongly convex with respect to the norm  $\|\cdot\|_{v \bullet p^{-1}}$ .

Let  $x_*$  be a minimizer of f, i.e., an optimal solution of (1). Minimizing both sides of (8) in h, we get

$$f(x_*) - f(x) \stackrel{(8)}{\geq} \min_{h \in \mathbb{R}^n} \langle \nabla f(x), h \rangle + \frac{\mu}{2} ||h||_{v \bullet p^{-1}}^2$$

$$= -\frac{1}{2\mu} ||\nabla f(x)||_{p \bullet v^{-1}}^2. \tag{9}$$



#### Proof of Theorem 3 - Part II

Let  $h_k \stackrel{\text{def}}{=} -v^{-1} \bullet \nabla f(x_k)$ . Then  $x_{k+1} = x_k + (h_k)_{[\hat{S}]}$ , and utilizing Assumption 2, we get

$$\mathbf{E}[f(x_{k+1}) \mid x_{k}] = \mathbf{E}\left[f(x_{k} + (h_{k})_{[\hat{S}]}) \mid x_{k}\right]$$

$$\stackrel{(5)}{\leq} f(x_{k}) + \langle \nabla f(x_{k}), h_{k} \rangle_{p} + \frac{1}{2} \|h_{k}\|_{p \bullet v}^{2}$$

$$= f(x_{k}) - \frac{1}{2} \|\nabla f(x_{k})\|_{p \bullet v^{-1}}^{2}$$

$$\stackrel{(9)}{\leq} f(x_{k}) - \mu(f(x_{k}) - f(x_{*})).$$

► Taking expectations in the last inequality,

$$\mathbf{E}[f(x_k) - f(x_*)] \le (1 - \mu)^k (f(x_0) - f(x_*)). \tag{10}$$

 $\triangleright$  Using Markov inequality, (10) and the definition of K, we finally get

$$\mathbf{P}(f(x_{K}) - f(x_{*}) \ge \epsilon) \le \mathbf{E}[f(x_{K}) - f(x_{*})]/\epsilon 
\stackrel{(10)}{\le} (1 - \mu)^{K} (f(x_{0}) - f(x_{*}))/\epsilon \stackrel{(6)}{\le} \rho.$$

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#### Proof of Theorem 3 - Part III

▶ Finally, let us now establish the lower bound on Λ. Letting  $\Delta \stackrel{\text{def}}{=} \{ p' \in \mathbb{R}^n : p' \geq 0, \sum_i p'_i = \mathbf{E}[|\hat{S}|] \}$ , we have

$$\Lambda \stackrel{(7)}{=} \max_{i} \frac{v_{i}}{p_{i}s_{i}} \stackrel{(3)}{\geq} \min_{p' \in \Delta} \max_{i} \frac{v_{i}}{p'_{i}s_{i}} = \frac{1}{\mathbf{E}[|\hat{S}|]} \sum_{i=1}^{n} \frac{v_{i}}{s_{i}},$$

where the last equality follows since optimal  $p'_i$  is proportional to  $v_i/s_i$ .

## Lecture 2 BLOCKS



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#### The idea

We now assume the decision vector x has N coordinates

$$x \in \mathbb{R}^N$$

which we partition into n "blocks".

**Idea:** We let the algorithm operate on "block level" instead  $\Rightarrow$  block coordinate descent. That is, at iteration k,

- ▶ a random subset  $S_k$  of blocks  $[n] = \{1, 2, ..., n\}$  is chosen
- and updated.

#### What do we gain by introducing blocks?

- ► Flexibility: We can partition the coordinates any way we like for any reason we might have.
  - Sometimes block structure is implied by the problem at hand. In L1 optimization, one often chooses  $N_i = 1$  for all i. In group LASSO problems, groups correspond to blocks.
- ► **Generality:** By allowing for general block structure, we simultaneously analyze several classes of algorithms:
  - **coordinate descent** (if we choose  $N_i = 1$  for all i)
  - **block coordinate descent** (if we choose  $N_i > 1$  and n > 1)
  - **gradient descent** (if we choose n = 1)
  - ▶ fast  $(O(1/k^2))$  versions of the above...
- ▶ Efficiency: It is sometimes more efficient to have blocks because:
  - ► this leads to a more "chunky" workload for each processor if we think that each processor handles one block
  - one can design block-norms based on data, which leads to better approximation and hence faster convergence
  - one can try to optimize the partitioning of coordinates to blocks (say, by trying to optimize complexity bounds, which depend on block structure)



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#### Block Decomposition of $\mathbb{R}^N$

- ▶ Partition. Let  $H_1, ..., H_n$  be a partition of the set of coordinates/variables  $\{1, 2, ..., N\}$  into n nonempty subsets. Let  $N_i = |H_i|$ .
- ▶ Projection/lifting matrices. Let  $U_i \in \mathbb{R}^{N \times N_i}$  be the column submatrix of the  $N \times N$  identity matrix corresponding to coordinates in  $H_i$ .
- ▶ Projection of  $\mathbb{R}^N$  to  $\mathbb{R}^{N_i}$  For  $x \in \mathbb{R}^N$ , define

$$x^{(i)} \stackrel{\mathsf{def}}{=} U_i^T x \in \mathbb{R}^{N_i}, \quad i = 1, 2, \dots, n.$$

Notice that  $x^{(i)}$  is the block of coordinates of x belonging to  $H_i$ .

▶ Lifting  $\mathbb{R}^{N_i}$  to  $\mathbb{R}^N$ . Given  $x^{(i)} \in \mathbb{R}^{N_i}$ , notice that the vector  $s = U_i x^{(i)} \in \mathbb{R}^N$  has all blocks equal to 0 except for block i, which is equal to  $x^{(i)}$ . That is,

$$s^{(j)} = \begin{cases} x^{(j)} & j = i \\ 0 & \text{otherwise.} \end{cases}$$



#### Examples - Part I

#### Example 4

1. Single block.

$$n = 1;$$
  $H_1 = \{1, 2, ..., N\};$   $U_1 = I$ 

2. Blocks of size 1. This is the setting already introduced in NSync:

$$N = n;$$
  $H_i = \{i\};$   $U_i = e_i$ 

**3.** Two blocks of different sizes. Let N = 5 (5 coordinates), n = 2 (2 blocks) and let the partitioning be given by

$$H_1 = \{1,3\}, \quad H_2 = \{2,4,5\}.$$

Then

$$U_1 = \left( egin{array}{ccc} 1 & 0 \ 0 & 0 \ 0 & 1 \ 0 & 0 \ 0 & 0 \end{array} 
ight) \qquad U_2 = \left( egin{array}{ccc} 0 & 0 & 0 \ 1 & 0 & 0 \ 0 & 0 & 0 \ 0 & 1 & 0 \ 0 & 0 & 1 \end{array} 
ight)$$



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#### Examples - Part II

For  $x \in \mathbb{R}^N = \mathbb{R}^5$  we have

$$x^{(1)} = U_1^T x = \begin{pmatrix} \mathbf{1} & 0 & 0 & 0 & 0 \\ 0 & 0 & \mathbf{1} & 0 & 0 \end{pmatrix} \begin{pmatrix} x_1 \\ x_2 \\ x_3 \\ x_4 \\ x_5 \end{pmatrix} = \begin{pmatrix} x_1 \\ x_3 \end{pmatrix} \in \mathbb{R}^{N_1} = \mathbb{R}^2$$

$$x^{(2)} = U_2^{\mathsf{T}} x = \begin{pmatrix} 0 & \mathbf{1} & 0 & 0 & 0 \\ 0 & 0 & 0 & \mathbf{1} & 0 \\ 0 & 0 & 0 & 0 & \mathbf{1} \end{pmatrix} \begin{pmatrix} x_1 \\ x_2 \\ x_3 \\ x_4 \\ x_5 \end{pmatrix} = \begin{pmatrix} x_2 \\ x_4 \\ x_5 \end{pmatrix} \in \mathbb{R}^{N_2} = \mathbb{R}^3$$

On the other hand, for any  $x \in \mathbb{R}^5$ :

$$U_1 x^{(1)} = U_1 (U_1^T x) = \begin{pmatrix} 1 & 0 \\ 0 & 0 \\ 0 & 1 \\ 0 & 0 \\ 0 & 0 \end{pmatrix} \begin{pmatrix} x_1 \\ x_3 \end{pmatrix} = \begin{pmatrix} x_1 \\ 0 \\ x_3 \\ 0 \\ 0 \end{pmatrix} \in \mathbb{R}^5$$



#### Examples - Part III

and

$$U_2 x^{(2)} = U_2 (U_2^T x) = \left( egin{array}{ccc} 0 & 0 & 0 & 0 \ 1 & 0 & 0 & 0 \ 0 & 1 & 0 & 0 \ 0 & 0 & 1 \end{array} 
ight) \left( egin{array}{c} x_2 \ x_4 \ x_5 \end{array} 
ight) = \left( egin{array}{c} 0 \ x_2 \ 0 \ x_4 \ x_5 \end{array} 
ight) \in \mathbb{R}^5$$

So, we have the unique decomposition:

$$x = U_1 x^{(1)} + U_2 x^{(2)}$$

The next simple result will formalize this.



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#### Block Decomposition: Formal Statement

#### Proposition 1 (Block Decomposition)

Any vector  $x \in \mathbb{R}^N$  can be written uniquely as

$$x = \sum_{i=1}^{n} U_i x^{(i)}, \tag{11}$$

where  $x^{(i)} \in \mathbb{R}^{N_i}$ . Moreover,

$$x^{(i)} = U_i^T x. (12)$$

#### Proof

Fix any  $x \in \mathbb{R}^N$ . Noting that  $\sum_i U_i U_i^T$  is the  $N \times N$  identity matrix, we have  $x = \sum_i U_i U_i^T x$ , where  $U_i^T x \in \mathbb{R}^{N_i}$ . Let us now show uniqueness. Assume that  $x = \sum_i U_i x_1^{(i)} = \sum_i U_i x_2^{(i)}$ , where  $x_1^{(i)}, x_2^{(i)} \in \mathbb{R}^{N_i}$ . Since

$$U_j^T U_i = \begin{cases} N_j \times N_j & \text{identity matrix,} & \text{if } i = j, \\ N_j \times N_i & \text{zero matrix,} & \text{otherwise,} \end{cases}$$
 (13)

we get 
$$0 = U_i^T(x - x) = U_i^T \sum_i U_i(x_1^{(i)} - x_2^{(i)}) = x_1^{(j)} - x_2^{(j)}$$
, for all  $j$ .



#### Projection onto (a subspace spanned by) a set of blocks

For  $h \in \mathbb{R}^N$  and  $\emptyset \neq S \subseteq [n] \stackrel{\mathsf{def}}{=} \{1, 2, \dots, n\}$ , we write

$$h_{[S]} = \sum_{i \in S} U_i h^{(i)}.$$
 (14)

In words,  $h_{[S]}$  is a vector in  $\mathbb{R}^N$  obtained from  $h \in \mathbb{R}^N$  by zeroing out the blocks that do not belong to S. Hence:

$$(h_{[S]})^{(i)} = \begin{cases} h^{(i)}, & i \in S, \\ 0, & i \notin S. \end{cases}$$



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#### Norms in $\mathbb{R}^{N_i}$ and $\mathbb{R}^N$

With each block  $i \in [n]$  we associate a positive definite matrix  $B_i \in \mathbb{R}^{N_i \times N_i}$  and a scalar  $v_i > 0$ , and equip  $\mathbb{R}^{N_i}$  and  $\mathbb{R}^N$  with the norms

$$\|x^{(i)}\|_{(i)} \stackrel{\text{def}}{=} \langle B_i x^{(i)}, x^{(i)} \rangle^{1/2}, \quad \|x\|_{\nu} \stackrel{\text{def}}{=} \left(\sum_{i=1}^n v_i \|x^{(i)}\|_{(i)}^2\right)^{1/2}.$$
 (15)

The corresponding conjugate norms, defined by

$$||s||^* = \max\{\langle s, x \rangle : ||x|| < 1\}$$

are given by

$$||x^{(i)}||_{(i)}^* \stackrel{\mathsf{def}}{=} \langle B_i^{-1} x^{(i)}, x^{(i)} \rangle^{1/2}, \quad ||x||_v^* = \left( \sum_{i=1}^n \frac{1}{v_i} \left( ||x^{(i)}||_{(i)}^* \right)^2 \right)^{1/2}. \tag{16}$$



#### Norms: Examples

#### Example 5

Consider the following extreme special cases:

1. **Single block.** Let n = 1, v = 1 and B be a positive definite matrix. Then

$$||x||_{(1)} = ||x||_v = \langle Bx, x \rangle^{1/2}, \quad x \in \mathbb{R}^N.$$

For instance, if  $f(x) = \frac{1}{2} ||Ax - b||^2$  we may choose:

- $B = A^T A$  (assuming  $A^T A$  is positive definite)
- ▶  $B = Diag(A^T A)$  (assuming no column in A is zero,  $A^T A$  is positive definite)
- 2. Blocks of size one. Let  $N_i = 1$  for all i and set  $B_i = 1$ . Then

$$||t||_{(i)} = ||t||_{(i)}^* = |t|, \qquad t \in \mathbb{R}$$

and

$$||x||_{v} = \left(\sum_{i=1}^{n} v_{i}(x^{(i)})^{2}\right)^{1/2}, \qquad x \in \mathbb{R}^{N}.$$



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#### **Exercises**

#### Exercise 1

Show that  $\|\cdot\|_{v}^{*}$ , as defined above, is indeed the conjugate norm of  $\|\cdot\|_{v}$ .

#### Exercise 2

Generalize NSync to the block setting and provide a complexity analysis.

# Lecture 3 SAMPLINGS



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#### Samplings: Definition

#### Definition 6 (Sampling)

**Sampling** is a *random set-valued mapping*  $\hat{S}$  with values in  $2^{[n]}$ , the collection of subsets of  $[n] = \{1, 2, ..., n\}$ .

lacktriangle A sampling  $\hat{S}$  is uniquely characterized by the probability mass function

$$\mathbf{P}(S) \stackrel{\text{def}}{=} \mathbf{P}(\hat{S} = S), \quad S \subseteq [n]; \tag{17}$$

that is, by assigning probabilities to all subsets of [n].

► Let

$$p_i \stackrel{\text{def}}{=} \mathbf{P}(i \in \hat{S}). \tag{18}$$

Let

$$p_{ij} \stackrel{\mathsf{def}}{=} \mathbf{P}(i \in \hat{S}, j \in \hat{S}) = \sum_{S:\{i,j\} \subset S} \mathbf{P}(S). \tag{19}$$



#### Sampling Zoo - Part I

Why consider different samplings?

- 1. Basic Considerations. It is important that each block has a positive probability of being chosen, otherwise an algorithm will not be able to update some blocks and hence will not converge to optimum. For technical/sanity reasons, we define:
  - ▶ Proper sampling.  $p_i = P(i \in \hat{S}) > 0$  for all blocks  $i \in [n]$
  - Nil sampling:  $P(\hat{S} = \emptyset) = 1$
  - **Vacuous sampling:**  $P(\hat{S} = \emptyset) > 0$
- 2. Parallelism. Choice of sampling affects the level of parallelism:
  - ▶  $\mathbf{E}[|\hat{S}|]$  is the average number of updates performed in parallel in one iteration; and is hence closely related to the number of iterations.
  - serial sampling: picks one block:

$$\mathbf{P}(|\hat{S}|=1)=1$$

We call this sampling serial although nothing prevents us from computing the actual update to the block, and/or to apply he update in parallel.



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#### Sampling Zoo - Part II

fully parallel sampling: always picks all blocks:

$$P(\hat{S} = \{1, 2, ..., n\}) = 1$$

- 3. **Processor reliability.** Sampling may be induced/informed by the computing environment:
  - ▶ Reliable/dedicated processors. If one has reliable processors, it is sensible to choose sampling  $\hat{S}$  such that  $P(|\hat{S}| = \tau)$  1 for some  $\tau$  related to the number of processors.
  - ▶ Unreliable processors. If processors given a computing task are busy or unreliable, they return answer later or not at all it is then sensible to ignore such updates and move on. This then means that  $\hat{S}$  varies from iteration to iteration.
- 4. **Distributed computing.** In a distributed computing environment it is sensible:
  - to allow each node as much autonomy as possible so as to minimize communication cost,
  - ▶ to make sure all nodes are busy at all times



#### Sampling Zoo - Part III

This suggests a strategy where the set of blocks is partitioned, with each node owning a partition, and independently picking a "chunky" subset of blocks at each iteration it will update, ideally from local information.

- 5. **Uniformity.** It may or not may make sense to update some blocks more often than others:
  - uniform samplings:

$$P(i \in \hat{S}) = P(j \in \hat{S})$$
 for all  $i, j \in [n]$ 

doubly uniform (DU): These are samplings characterized by:

$$|S'| = |S''| \implies \mathbf{P}(\hat{S} = S') = \mathbf{P}(\hat{S} = S'')$$
 for all  $S', S'' \subseteq [n]$ 

ightharpoonup au-nice: DU sampling with the additional property that

$$\mathbf{P}(|\hat{S}|=\tau)=1$$

- **distributed**  $\tau$ -nice: will define later
- ▶ independent sampling: union of independent uniform serial samplings
- nonuniform samplings



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#### Sampling Zoo - Part IV

- 6. Complexity of generating a sampling. Some samplings are computationally more efficient to generate than others: the potential benefits of a sampling may be completely ruined by the difficulty to generate sets according to the sampling's distribution.
  - ightharpoonup a au-nice sampling can be well approximated by an independent sampling, which is easy to generate. . .
  - ▶ a general sampling, as considered in NSync, will be hard to generate



#### **Basic Identity**

#### Theorem 7 (Sum over a random index set)

Let  $\emptyset \neq J, J_1, J_2 \subset [n]$  and  $\hat{S}$  be any sampling. If  $\theta_i$ ,  $i \in [n]$ , and  $\theta_{ij}$ , for  $(i,j) \in [n] \times [n]$  are real constants, then<sup>1</sup>

$$\mathbf{E}\left[\sum_{i\in J\cap\hat{S}}\theta_i\right]=\sum_{i\in J}p_i\theta_i,$$

$$\mathbf{E}\left[\sum_{i\in J\cap\hat{S}}\theta_i\mid |J\cap\hat{S}|=k\right]=\sum_{i\in J}\mathbf{P}(i\in\hat{S}\mid |J\cap\hat{S}|=k)\theta_i,\tag{20}$$

$$\mathbf{E}\left[\sum_{i\in J_1\cap\hat{S}}\sum_{j\in J_2\cap\hat{S}}\theta_{ij}\right] = \sum_{i\in J_1}\sum_{j\in J_2}p_{ij}\theta_{ij}.\tag{21}$$



<sup>&</sup>lt;sup>1</sup>Sum over an empty index set will, for convenience, be defined to be zero.

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#### Proof of Theorem 7

We prove the first statement, proof of the remaining statements is essentially identical:

$$\mathbf{E}\left[\sum_{i\in J\cap\hat{S}}\theta_{i}\right] \stackrel{(17)}{=} \sum_{S\subset[n]} \left(\sum_{i\in J\cap S}\theta_{i}\right) \mathbf{P}(\hat{S}=S)$$

$$= \sum_{i\in J} \sum_{S:i\in S}\theta_{i}\mathbf{P}(\hat{S}=S)$$

$$= \sum_{i\in J} \theta_{i} \sum_{S:i\in S}\mathbf{P}(\hat{S}=S)$$

$$= \sum_{i\in J} p_{i}\theta_{i}.$$



#### Consequences of Theorem 7

#### Corollary 8 ([5])

Let  $\emptyset \neq J \subset [n]$  and  $\hat{S}$  be an arbitrary sampling. Further, let  $a, h \in \mathbb{R}^N$ ,  $w \in \mathbb{R}^n_+$  and let g be a block separable function, i.e.,  $g(x) = \sum_i g_i(x^{(i)})$ . Then

$$\mathbf{E}\left[|J\cap\hat{S}|\right] = \sum_{i\in J} p_i,\tag{22}$$

$$\mathbf{E}\left[|J\cap\hat{S}|^2\right] = \sum_{i\in I}\sum_{j\in I}p_{ij},\tag{23}$$

$$\mathbf{E}\left[\langle a, h_{[\hat{S}]}\rangle_{w}\right] = \langle a, h\rangle_{p\bullet w}, \tag{24}$$

$$\mathsf{E}\left[\|h_{[\hat{S}]}\|_{w}^{2}\right] = \|h\|_{p\bullet w}^{2},\tag{25}$$

$$\mathbf{E}\left[g(x+h_{[\hat{S}]})\right] = \sum_{i=1}^{n} \left[p_{i}g_{i}(x^{(i)}+h^{(i)})+(1-p_{i})g_{i}(x^{(i)})\right]. (26)$$

Moreover, the matrix  $P \stackrel{\text{def}}{=} (p_{ij})$  is positive semidefinite.



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#### **Proof of Corollary 8**

All 5 identities follow by applying Lemma 7 and observing that:

- $|J \cap \hat{S}| = \sum_{i \in J \cap \hat{S}} 1$
- ►  $|J \cap \hat{S}|^2 = (\sum_{i \in J \cap \hat{S}} 1)^2 = \sum_{i \in J \cap \hat{S}} \sum_{j \in J \cap \hat{S}} 1$
- $\langle a, h_{[\hat{S}]} \rangle_w = \sum_{i \in \hat{S}} w_i \langle a^{(i)}, h^{(i)} \rangle$
- $\|h_{[\hat{S}]}\|_{w}^{2} = \sum_{i \in \hat{S}} w_{i} \|h^{(i)}\|_{(i)}^{2}$  and

**•** 

$$g(x + h_{[\hat{S}]}) = \sum_{i \in \hat{S}} g_i(x^{(i)} + h^{(i)}) + \sum_{i \notin \hat{S}} g_i(x^{(i)})$$

$$= \sum_{i \in \hat{S}} g_i(x^{(i)} + h^{(i)}) + \sum_{i=1}^n g_i(x^{(i)}) - \sum_{i \in \hat{S}} g_i(x^{(i)}),$$

Finally, for any  $\theta = (\theta_1, \dots, \theta_n)^T \in \mathbb{R}^n$ ,

$$\theta^T P \theta = \sum_{i=1}^n \sum_{j=1}^n p_{ij} \theta_i \theta_j \stackrel{\text{(21)}}{=} \mathbf{E} \left[ \left( \sum_{i \in \hat{S}} \theta_i \right)^2 \right] \geq 0.$$

Remark: The above results hold for arbitrary samplings. Let us specialize them, in order of decreasing generality, to uniform, doubly uniform and nice samplings.



#### Identities: uniform samplings

If  $\hat{S}$  is uniform, then from (22) using J = [n] we get

$$p_i = \frac{\mathbf{E}[|\hat{S}|]}{n}, \qquad i \in [n]. \tag{27}$$

Plugging (27) into (22), (24), (25) and (26) yields

$$\mathbf{E}\left[|J\cap\hat{S}|\right] = \frac{|J|}{n}\mathbf{E}[|\hat{S}|],\tag{28}$$

$$\mathbf{E}\left[\langle a, h_{[\hat{S}]}\rangle_{w}\right] = \frac{\mathbf{E}\left[|\hat{S}|\right]}{n}\langle a, h\rangle_{w},\tag{29}$$

$$\mathbf{E}\left[\|h_{[\hat{S}]}\|_{w}^{2}\right] = \frac{\mathbf{E}\left[|\hat{S}|\right]}{n}\|h\|_{w}^{2},\tag{30}$$

$$\mathbf{E}\left[g(x+h_{[\hat{S}]})\right] = \frac{\mathbf{E}[|\hat{S}|]}{n}g(x+h) + \left(1 - \frac{\mathbf{E}[|\hat{S}|]}{n}\right)g(x). \tag{31}$$



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#### Identities: doubly uniform samplings

Consider the case n > 1; the case n = 1 is trivial. For doubly uniform  $\hat{S}$ ,  $p_{ij}$  is constant for  $i \neq j$ :

$$p_{ij} = \frac{\mathbf{E}[|\hat{S}|^2 - |\hat{S}|]}{n(n-1)}.$$
 (32)

Indeed, this follows from

$$p_{ij} = \sum_{k=1}^{n} \mathbf{P}(\{i,j\} \subseteq \hat{S} \mid |\hat{S}| = k) \mathbf{P}(|\hat{S}| = k) = \sum_{k=1}^{n} \frac{k(k-1)}{n(n-1)} \mathbf{P}(|\hat{S}| = k).$$

Substituting (32) and (27) into (23) then gives

$$\mathbf{E}[|J \cap \hat{S}|^2] = (|J|^2 - |J|) \frac{\mathbf{E}[|\hat{S}|^2 - |\hat{S}|]}{n \max\{1, n - 1\}} + |J| \frac{|\hat{S}|}{n}.$$
 (33)



#### Identities: $\tau$ -nice sampling

Finally, if  $\hat{S}$  is  $\tau$ -nice (and  $\tau \neq 0$ ), then  $\mathbf{E}[|\hat{S}|] = \tau$  and  $\mathbf{E}[|\hat{S}|^2] = \tau^2$ , which used in (33) gives

$$\mathbf{E}[|J \cap \hat{S}|^2] = \frac{|J|\tau}{n} \left( 1 + \frac{(|J|-1)(\tau-1)}{\max\{1, n-1\}} \right). \tag{34}$$

Moreover, assume that  $\mathbf{P}(|J \cap \hat{S}| = k) \neq 0$  (this happens precisely when  $0 \leq k \leq |J|$  and  $k \leq \tau \leq n - |J| + k$ ). Then for all  $i \in J$ ,

$$\mathbf{P}(i \in \hat{S} \mid |J \cap \hat{S}| = k) = \frac{\binom{|J|-1}{k-1}\binom{n-|J|}{\tau-k}}{\binom{|J|}{k}\binom{n-|J|}{\tau-k}} = \frac{k}{|J|}.$$

Substituting this into (20) yields

$$\mathbf{E}\left[\sum_{i\in J\cap\hat{S}}\theta_i\mid |J\cap\hat{S}|=k\right]=\frac{k}{|J|}\sum_{i\in J}\theta_i. \tag{35}$$



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#### Elementary Samplings, Intersection and Restriction

#### Definition 9 (Elementary samplings)

Elementary sampling associated with  $J \subseteq [n]$  is sampling  $\hat{E}_J$  for which

$$\mathbf{P}(\hat{E}_J = J) = 1.$$

#### Definition 10 (Intersection of samplings)

For two samplings  $\hat{S}_1$  and  $\hat{S}_2$  we define the intersection  $\hat{S} \stackrel{\text{def}}{=} \hat{S}_1 \cap \hat{S}_2$  as the sampling for which:

$$\mathbf{P}(\hat{S} = S) = \mathbf{P}(\hat{S}_1 \cap \hat{S}_2 = S), \quad S \subseteq [n].$$

#### Definition 11 (Restriction of a sampling to a subset)

Let  $\hat{S}$  be a sampling and  $J \subseteq [n]$ . By restriction of  $\hat{S}$  to J we mean the sampling





#### Probability matrices associated with samplings - Part I

#### Definition 12 (Probability matrix)

With arbitrary sampling  $\hat{S}$  we associate an *n*-by-*n* matrix  $P = P(\hat{S})$  with entries

$$[P(\hat{S})]_{ij} = \mathbf{P}(i \in \hat{S}, j \in \hat{S}).$$

#### Lemma 13 (Intersection of independent samplings; [14])

Let  $\hat{S}_1, \hat{S}_2$  be independent samplings. Then

$$P(\hat{S}_1 \cap \hat{S}_2) = P(\hat{S}_1) \bullet P(\hat{S}_2).$$

That is, the probability matrix of an intersection of independent samplings is the Hadamard product of their probability matrices.

#### Proof.

$$[P(\hat{S}_1 \cap \hat{S}_2)]_{ij} = \mathbf{P}(\{i,j\} \in \hat{S}_1 \cap \hat{S}_2) = \mathbf{P}(\{i,j\} \in \hat{S}_1)\mathbf{P}(\{i,j\} \in \hat{S}_2) = [P(\hat{S}_1)]_{ij}[P(\hat{S}_2)]_{ij}.$$



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#### Probability matrices associated with samplings - Part II

Example 14 (Probability Matrix of an Elementary Sampling)

Note that the probability matrix of the elementary sampling  $\hat{E}_J$  is the matrix

$$P(\hat{E}_J) \stackrel{\text{def}}{=} e_J e_J^T, \tag{36}$$

where  $e_J$  we denote the binary vector in  $\mathbb{R}^n$  with ones in places corresponding to set J. That is,

$$[P(\hat{E}_J)]_{ij} = \begin{cases} 1 & i,j \in J, \\ 0 & \text{otherwise.} \end{cases}$$

Hence, for arbitrary sampling  $\hat{S}$ , the probability matrix of  $J \cap \hat{S}$  is the submatrix of  $P(\hat{S})$  corresponding to the rows and columns indexed by J:

$$[P(J \cap \hat{S})]_{ij} = [P(\hat{E}_J) \bullet P(\hat{S})]_{ij} = \begin{cases} [P(\hat{S})]_{ij}, & i, j \in J, \\ 0, & \text{otherwise.} \end{cases}$$
(37)



## Probability matrices associated with samplings - Part III Lemma 15 (Decomposition of a Probability Matrix; [14])

Let  $\hat{S}$  be any sampling. Then

$$P(\hat{S}) = \sum_{S \subseteq [n]} P(\hat{S} = S) P(\hat{E}_S). \tag{38}$$

That is, the probability matrix of arbitrary sampling is a convex combination of elementary probability matrices.

#### Proof.

Fix any  $i, j \in [n]$ . Since  $(P(\hat{E}_S))_{ij} = 1$  iff  $\{i, j\} \subseteq S$ , from definition we have

$$(P(\hat{S}))_{ij} = \sum_{S:\{i,j\}\subseteq S} \mathbf{P}(\hat{S} = S)$$

$$= \sum_{S:\{i,j\}\subseteq S} \mathbf{P}(\hat{S} = S)(P(\hat{E}_S))_{ij}$$

$$= \left(\sum_{S:\{i,j\}\subseteq S} \mathbf{P}(\hat{S} = S)P(\hat{E}_S)\right)_{ij}.$$



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#### Sampling Identity for a Quadratic

#### Lemma 16

Let G be any real  $n \times n$  matrix and  $\hat{S}$  an arbitrary sampling. Then for any  $h \in \mathbb{R}^n$  we have

$$\mathbf{E}\left[h_{[\hat{S}]}^{T}Gh_{[\hat{S}]}\right] = h^{T}\left(P(\hat{S}) \bullet G\right)h,\tag{39}$$

where • denotes the Hadamard (elementwise) product of matrices.

#### Proof.

$$\mathbf{E}\left[h_{[\hat{S}]}^{T}Gh_{[\hat{S}]}\right] \stackrel{(14)}{=} \mathbf{E}\left[\sum_{i\in\hat{S}}\sum_{j\in\hat{S}}G_{ij}h^{(i)}h^{(j)}\right]$$

$$\stackrel{(21)}{=} \sum_{i=1}^{n}\sum_{j=1}^{n}p_{ij}G_{ij}h^{(i)}h^{(j)} = h^{T}\left(P(\hat{S})\bullet G\right)h.$$

#### Distributed sampling

The following sampling is useful in the design of a **distributed** coordinate descent method.

#### Definition 17 (Distributed $\tau$ -nice sampling; [10, 13])

Let  $\mathcal{P}_1, \ldots, \mathcal{P}_c$  be a partition of  $\{1, 2, \ldots, n\}$  such that  $|\mathcal{P}_I| = s$  for all I. That is, sc = n. Now let  $\hat{S}_1, \ldots, \hat{S}_c$  be independent  $\tau$ -nice samplings from  $\mathcal{P}_1, \ldots, \mathcal{P}_c$ , respectively. Then the sampling

$$\hat{S} \stackrel{\text{def}}{=} \cup_{l=1}^{c} \hat{S}_{l}, \tag{40}$$

is called distributed  $\tau$ -nice sampling.

Idea: Blocks in  $\mathcal{P}_I$ , and all associated data, will be handled/stored by computer/node I only. Node I picks blocks in  $\hat{S}_I$ , computes the updates fro local information, and applies the updates to locally stored  $x^{(i)}$  for  $i \in \mathcal{P}_I$ .



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#### Probability Matrix of Distributed $\tau$ -nice Sampling

Consider the distributed  $\tau$ -nice sampling and define:

- $E = P(\hat{E}_{[n]})$ : the  $n \times n$  matrix of all ones
- ▶ *I* be the  $n \times n$  identity matrix
- ▶  $B = \sum_{l=1}^{c} P(\hat{E}_{P_l})$ : the 0-1 matrix with  $B_{ij} = 1$  iff i, j belong to the same partition

#### Lemma 18 ([10]; presented in a different form)

Consider the distributed  $\tau$ -nice sampling  $\hat{S}$ . Its probability matrix can be written as

$$P(\hat{S}) = \frac{\tau}{s} \left[ \alpha_1 I + \alpha_2 E + \alpha_3 (E - B) \right], \tag{41}$$

where

$$lpha_1=1-rac{ au-1}{ss_1}, \qquad lpha_2=rac{ au-1}{s_1}, \qquad lpha_3=rac{ au}{s}-rac{ au-1}{s_1},$$

and  $s_1 = \max\{1, s - 1\}$ .



#### Proof of Lemma 18

Let  $P = P(\hat{S})$ . It is easy to see that

- $P_{ij} = \frac{\tau}{s} \stackrel{\text{def}}{=} \beta_3 \text{ if } i = j,$
- ▶  $P_{ij} = \frac{\tau(\tau-1)}{ss_1} \stackrel{\text{def}}{=} \beta_2$  if  $i \neq j$  and i, j belong to the same partition,
- ▶  $P_{ij} = \frac{\tau^2}{s^2} \stackrel{\text{def}}{=} \beta_3$  if  $i \neq j$  belong to different partitions.

So, we can write

$$P = \beta_1 I + \beta_2 (B - I) + \beta_3 (E - B)$$
  
=  $(\beta_1 - \beta_2)I + \beta_2 E + (\beta_3 - \beta_2)(E - B).$ 



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#### **Exercises**

#### Exercise 3

Find an expression for the probability matrix of

- the  $\tau$ -nice sampling,
- arbitrary doubly uniform sampling.

#### Exercise 4

Let  $\hat{S}$  be any sampling. Show that

- $\lambda_{max}(P) \leq \mathbf{E}[|\hat{S}|]$  and that the bound is tight,
- $P \succeq pp^T$ .

## Lecture 4 FUNCTIONS



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#### Introduction

- ▶ In this part we describe three models for *f*.
- ► These models can be thought of as function classes described by a list of properties.
- ▶ However, a single function may belong to more function classes.

In big data setting, some information is computationally difficult to extract from data.

Consider  $f(x) = \frac{1}{2} ||Ax - b||^2$ .

- It is difficult to compute the largest eigenvalue of  $A^TA$  if A is large (this is the Lipschitz constant of  $\nabla f$  with respect to the standard Euclidean norm)
- ▶ It is easier to compute the squared norm of each column (these correspond to coordinate-wise Lipschitz constants).

**Important point:** The models differ in the amount of information they reveal about f.



#### Model: Quadratic

#### Model 1 ([10, 13])

We assume that

1. **Structure and Smoothness:**  $f : \mathbb{R}^N \to \mathbb{R}$  is differentiable and for all  $x, h \in \mathbb{R}^N$  satisfies

$$f(x+h) \le f(x) + (\nabla f(x))^T h + \frac{1}{2} h^T A^T A h,$$
 (42)

where  $A \in \mathbb{R}^{m \times N}$ .

2. **Sparsity:** Row j of A depends on blocks  $i \in C_j$  only. Formally,

$$C_j \stackrel{\text{def}}{=} \{i : A_{ji} \neq 0\},\$$

where  $A_{ji} \stackrel{def}{=} e_j^T A U_i \in \mathbb{R}^{1 \times N_i}$ . Let  $\omega_j \stackrel{def}{=} |C_j|$ .

3. Convexity: f is convex.

Remark: Information about f is contained in the matrix A.



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#### **Examples**

#### Example 19

In machine learning (ML), functions f of the following form are common:

$$f(x) = \sum_{j=1}^{m} f_j(x) = \sum_{j=1}^{m} \ell(x; a_j, y^j),$$

where N is the number of features, m number of examples,  $a_j \in \mathbb{R}^N$  corresponds to jth example and  $y^j$  is a label associated with jth example.

Here are some convex loss functions  $\ell$  often used in ML for which the total loss f satisfies (42):

Loss function $\ell$	$f_j(x)$	(42) satisfied for A given by
square loss (SL)	$\frac{1}{2}(y^j - a_j^T x)^2$	$A_{j:}=a_{j}^{T}$
logistic loss (LL)	$\log(1 + \exp(-y^j a_j^T x))$	$A_{j:}=rac{1}{2}a_{j}^{T}$
square hinge loss (HL)	$\frac{1}{2}\max\{0,1-y^ja_i^Tx\}^2$	$A_{j:} = a_i^T$

Interpretation of  $\omega_j$  (point 2 in Model 1) : # features in example j



#### **Block gradients**

#### Definition 20 (Block Gradients)

The *i*th **block gradient** of  $f: \mathbb{R}^N \to \mathbb{R}$  at x is defined to be the *i*th block of the gradient of f at x:

$$\nabla_i f(x) \stackrel{\text{def}}{=} (\nabla f(x))^{(i)} = U_i^T \nabla f(x) \in \mathbb{R}^{N_i}. \tag{43}$$

In other words,  $\nabla_i f(x)$  is the vector of partial derivatives with respect to coordinates belonging to block i.



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#### Model: Classical

#### Model 2 ([2, 5, 9])

We assume that

1. **Structure:** Function  $f: \mathbb{R}^N \to \mathbb{R}$  is of the form

$$f(x) = \sum_{j=1}^{m} f_j(x).$$

- 2. **Sparsity:**  $f_j$  depends on x via blocks  $i \in C_j$  only.
- 3. Convexity: Functions  $\{f_j\}$  are convex.
- 4. Smoothness: Function f has block-Lipschitz gradient with constants  $L_1, \ldots, L_n > 0$ . That is, for all  $i = 1, 2, \ldots, n$ ,

$$\|\nabla_i f(x + U_i t) - \nabla_i f(x)\|_{(i)}^* \le L_i \|t\|_{(i)}, \quad x \in \mathbb{R}^N, \ t \in \mathbb{R}^{N_i}.$$
 (44)

Remark: Information about f is contained in the constants  $L_1, \ldots, L_n$ .



#### **Examples**

#### Example 21 (Least squares)

Consider the quadratic function  $f(x) = \frac{1}{2} ||Ax - b||^2$ .

(i) Consider the block setup with  $N_i = 1$  (all blocks are of size 1) and  $B_i = 1$  for all  $i \in [n]$  (standard Eucl. norms for each block:  $||t||_{(i)} = |t|$ ). Then  $U_i = e_i$  and

$$\|\nabla_{i}f(x+U_{i}t) - \nabla_{i}f(x)\|_{(i)}^{*} = |e_{i}^{T}A^{T}(A(x+te_{i})-b) - e_{i}^{T}A^{T}(Ax-b)|$$
$$= |e_{i}^{T}A^{T}Ae_{i}||t| = ||A_{:i}||^{2}|t|,$$

whence  $L_i = ||A_{:i}||^2$ .

(ii) Choose nontrivial block sizes  $(N_i > 1)$  and define data-driven block norms with  $B_i = A_i^T A_i$ , where  $A_i = A U_i$ , assuming that  $B_i > 0$ . Then

$$\begin{split} \|\nabla_{i}f(x+U_{i}t) - \nabla_{i}f(x)\|_{(i)}^{*} &= \|U_{i}^{T}A^{T}(A(x+U_{i}t)-b) - U_{i}^{T}A^{T}(Ax-b)\|_{(i)}^{*} \\ &= \|U_{i}^{T}A^{T}AU_{i}t\|_{(i)}^{*} \\ \stackrel{(16)}{=} \langle (A_{i}A_{i}^{T})^{-1}U_{i}^{T}A^{T}AU_{i}t, U_{i}^{T}A^{T}AU_{i}t \rangle^{1/2} \\ &= \langle B_{i}t, t \rangle^{1/2} \stackrel{(15)}{=} \|t\|_{(i)}, \end{split}$$

whence  $L_i = 1$ .



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#### Model: Newest

#### Model 3 ([12])

We assume that

1. **Structure:**  $f: \mathbb{R}^N \to \mathbb{R}$  is of the form

$$f(x) = \sum_{j=1}^{m} f_j(x).$$
 (45)

- 2. **Sparsity:**  $f_j$  depends on x via blocks  $i \in C_j$  only. Let  $\omega_j = |C_j|$ . (Note that  $i \notin C_j \Rightarrow L_{ji} = 0$ )
- 3. Convexity: Functions  $\{f_j\}$  are convex.
- 4. **Smoothness:** Functions  $\{f_j\}$  have block-Lipschitz gradient with constants  $L_{ji} \geq 0$ . That is, for all j = 1, 2, ..., m and i = 1, 2, ..., n,

$$\|\nabla_i f_j(x + U_i t) - \nabla_i f_j(x)\|_{(i)}^* \le L_{ji} \|t\|_{(i)}, \quad x \in \mathbb{R}^N, \ t \in \mathbb{R}^{N_i}.$$
 (46)

Remark: Information about f is contained in the constants  $\{L_{jj}\}$ 



#### Computation of $L_{ii}$

We now give a formula for the constants  $L_{ji}$  in the case when  $f_j$  arises as a composition of a scalar function  $\phi_j$  whose derivative has a known Lipschitz constant (this is often easy to compute), and a linear functional.

#### Proposition 2 ([12])

Let  $f_j(x) = \phi_j(e_j^T A x)$ , where  $\phi_j : \mathbb{R} \to \mathbb{R}$  is a function with  $L_{\phi_j}$ -Lipschitz derivative:

$$|\phi_i(s) - \phi_i(s')| \le L_{\phi_i}|s - s'|, \qquad s, s' \in \mathbb{R}. \tag{47}$$

Then  $f_i$  has a block Lipshitz gradient (i.e., satisfies (46)) with constants

$$L_{ji} = L_{\phi_j} \left( \|A_{ji}^T\|_{(i)}^* \right)^2, \qquad i = 1, 2, \dots, n,$$
 (48)

where

$$A_{ji} = e_j^T A U_i (49)$$

(i.e.,  $A_{ii}$  is the ith block of j-th row of A).



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#### Proof of Proposition 2

For any  $x \in \mathbb{R}^N$ ,  $t \in \mathbb{R}^{N_i}$  and i we have

$$\begin{split} \|\nabla_{i}f_{j}(x+U_{i}t) - \nabla_{i}f_{j}(x)\|_{(i)}^{*} \\ \stackrel{(43)}{=} & \|U_{i}^{T}(e_{j}^{T}A)^{T}\phi_{j}'(e_{j}^{T}A(x+U_{i}t)) - U_{i}^{T}(e_{j}^{T}A)^{T}\phi_{j}'(e_{j}^{T}Ax)\|_{(i)}^{*} \\ &= & \|A_{ji}^{T}\phi_{j}'(e_{j}^{T}A(x+U_{i}t)) - A_{ji}^{T}\phi_{j}'(e_{j}^{T}Ax)\|_{(i)}^{*} \\ &\leq & \|A_{ji}^{T}\|_{(i)}^{*}|\phi_{j}'(e_{j}^{T}A(x+U_{i}t)) - \phi_{j}'(e_{j}^{T}Ax)| \\ \stackrel{(47)}{\leq} & \|A_{ji}^{T}\|_{(i)}^{*}L_{\phi_{j}}|A_{ji}t| \leq & \|A_{ji}^{T}\|_{(i)}^{*}L_{\phi_{j}}\|A_{ji}^{T}\|_{(i)}^{*}\|t\|_{(i)}, \end{split}$$

where the last step follows by applying the Cauchy-Schwartz inequality.



#### **Examples**

#### Example 22 (Least squares)

Consider the quadratic function

$$f(x) = \frac{1}{2}||Ax - b||^2 = \frac{1}{2}\sum_{j=1}^{m}(e_j^TAx - b_j)^2.$$

Then  $f_j(x) = \phi_j(e_j^T A x)$ , where  $\phi_j(s) = \frac{1}{2}(s-b_j)^2$  and  $L_{\phi_j} = 1$ .

(i) Consider the block setup with  $N_i = 1$  (all blocks are of size 1) and  $B_i = 1$  for all  $i \in [n]$  (standard Euclidean norms for each block). Then by Proposition 2,

$$L_{ji} \stackrel{\text{(48)}}{=} L_{\phi_j} (\|A_{ji}^T\|_{(i)}^*)^2 = A_{ji}^2.$$

(ii) Choose nontrivial block sizes  $(N_i > 1)$  and define data-driven block norms with  $B_i = A_i^T A_i$ , where  $A_i = A U_i$ , assuming that the matrices  $A_i^T A_i$  are positive definite. Then by Proposition 2,

$$L_{ji} \stackrel{(48)}{=} L_{\phi_j} (\|A_{ji}^T\|_{(i)}^*)^2 \stackrel{(16)}{=} \langle (A_i^T A_i)^{-1} A_{ji}^T, A_{ji}^T \rangle \stackrel{(49)}{=} e_j^T A_i (A_i^T A_i)^{-1} A_i^T e_j.$$



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# Lecture 5 Expected Separable Overapproximation

#### Introduction

In this part we shall look at the three models of f (Lecture 3) and various types of samplings  $\hat{S}$  (Lecture 4) and compute paramters  $v = (v_1, \ldots, v_n)$  such

$$(f,\hat{S}) \sim ESO(v).$$

These parameters are important since:

- ▶ They are stepsize parameters needed in the algorithm (in NSync, but also in other randomized block coordinate descent methods).
- ► Their size as a function of  $\tau = \mathbf{E}[|\hat{S}|]$  describes achievable parallelization speedup.
- ▶ By computing v we get one step closer to ultimate goal of designing sampling  $\hat{S}$  optimizing the complexity bound.



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#### $\mathsf{ESO}(f \sim \mathsf{Model}\ 1,\ \hat{S} \sim \mathsf{arbitrary})$

#### Theorem 23 ([14])

Let f satisfy assumptions in Model 1, assume all blocks are of size 1  $(N_i = 1)$  and  $\hat{S}$  be any sampling. Then for all  $x, h \in \mathbb{R}^N$ ,

$$\mathbf{E}\left[f(x+h_{[\hat{\mathbf{S}}]})\right] \le f(x) + \langle \nabla f(x), h \rangle_p + \frac{1}{2} \|h\|_{p \bullet v}^2, \tag{50}$$

where v is any vector such that

$$P \bullet A^T A \leq \text{Diag}(p \bullet v),$$
 (51)

where  $P = P(\hat{S})$  is the probability matrix associated with  $\hat{S}$ .

Remark: The Hadamard product of two PSD matrices is PSD (P is PSD by Corollary 8).



#### Proof of Theorem 23

We have

$$\mathbf{E}\left[f(x+h_{[\hat{S}]})\right] \overset{(42)}{\leq} \mathbf{E}\left[f(x)+\langle\nabla f(x),h_{[\hat{S}]}\rangle+\frac{1}{2}\langle A^{T}Ah_{[\hat{S}]},h_{[\hat{S}]}\rangle\right]$$

$$\overset{(24)}{=} f(x)+\langle\nabla f(x),h\rangle_{p}+\frac{1}{2}\mathbf{E}\left[h_{[\hat{S}]}^{T}A^{T}Ah_{[\hat{S}]}\right]$$

$$\overset{(*)}{=} f(x)+\langle\nabla f(x),h\rangle_{p}+\frac{1}{2}h^{T}\left(P\bullet A^{T}A\right)h$$

$$\leq f(x)+\langle\nabla f(x),h\rangle_{p}+\frac{1}{2}\underbrace{h^{T}\operatorname{Diag}(p\bullet v)h}_{=\|h\|_{p\bullet v}^{2}},$$

where (\*) comes from Lemma 16.



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#### Ways of satisfying (51)

Let us fix a sampling  $\hat{S}$  (and hence P) and data A. We can find v for which  $P \bullet A^T A \leq \text{Diag}(p \bullet v)$  in several ways:

1. 
$$v_i = \lambda_1 \|A_{:i}\|^2$$
 and

$$\lambda_1 = \max_{\theta \in \mathbb{R}^n} \{ \theta^T (P \bullet A^T A) \theta : \theta^T \operatorname{Diag}(P \bullet A^T A) \theta \le 1 \}.$$

$$2. \ v_i = \frac{\lambda_{max}(P \bullet A^T A)}{p_i}.$$

3. 
$$v_i = \lambda_{max}(A^T A) \frac{(\max_j p_j)}{p_i}$$
 (using Lemma 24 with  $X = P$ )

4. 
$$v_i = \frac{\lambda_{max}(P)}{p_i} \max_i ||A_{:i}||^2$$
 (using Lemma 24 with  $X = A^T A$ )

#### Lemma 24

For any two PSD matrices X, Y with nonnegative elements,

$$\lambda_{max}(X \bullet Y) \leq \lambda_{max}(X) \max_{j} Y_{jj}.$$



#### Eigenvalues of Probability Matrices

#### Definition 25 (Eigenvalues)

For arbitrary sampling  $\hat{S}$  we define

$$\lambda(\hat{S}) \stackrel{\text{def}}{=} \max_{\theta \in \mathbb{R}^n} \{ \theta^T P(\hat{S}) \theta : \theta^T \operatorname{Diag}(P(\hat{S})) \theta \le 1 \}.$$
 (52)

and

$$\lambda'(\hat{S}) \stackrel{\text{def}}{=} \max_{\theta \in \mathbb{R}^n} \{ \theta^T P(\hat{S}) \theta : \theta^T \theta \le 1 \}.$$
 (53)

#### Example 26 (Elementary Sampling)

Fix  $S \subseteq [n]$  and consider the elementary sampling  $\hat{E}_S$ . Note that

$$\lambda(\hat{E}_S) = \lambda_{max}(P(\hat{E}_S)) = \lambda_{max}(e_S e_S^T) = ||e_S||^2 = |S|.$$
 (54)

Since  $J \cap \hat{E}_S = \hat{E}_{J \cap S}$ , we get

$$\lambda(J \cap \hat{E}_S) = \lambda(\hat{E}_{J \cap S}) \stackrel{(54)}{=} |J \cap S|. \tag{55}$$



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#### Insightful and Easily Computable Bound

Issues with Theorem 23:

- ▶ It does *not* provide insightful nor **easily computable** expressions for  $v_i$  (which are needed to run the algorithm).
- ▶ It does *not* answer the following **inverse problem**: given data matrix A and/or its sparsity pattern  $\{C_i\}$ , **design a "good" sampling.**

The following two results go a good way to overcoming these issues.

#### Theorem 27 (Useful ESO; [14])

Let the assumptions of Theorem 23 be satisfied. Then (51) holds (i.e.,  $(f, \hat{S}) \sim ESO(v)$ ) with v given by:

$$v_i = \sum_{i=1}^m \lambda(C_i \cap \hat{S}) A_{ji}^2, \quad i = 1, 2, \dots, n.$$
 (56)



#### Proof of Theorem 27

Note that it follows from (21) that for any vector  $\theta \in \mathbb{R}^n$  and any j the following identity holds:

$$\mathbf{E}\left[\left(\sum_{i\in C_j\cap\hat{S}}\theta_i\right)^2\right] = \sum_{i=1}^n [P(C_j\cap\hat{S})]_{ij}\theta_i\theta_j = \theta^T P(C_j\cap\hat{S})\theta. \tag{57}$$

Fix  $h \in \mathbb{R}^n$ . Let  $z_j = (z_j^{(1)}, \dots, z_j^{(n)})^T \in \mathbb{R}^n$  be defined as follows:  $z_j^{(i)} = h^{(i)}A_{ji}$ . We then have

$$\mathbf{E}\left[h_{[\hat{S}]}^{T}A^{T}Ah_{[\hat{S}]}\right] = \sum_{j=1}^{m} \mathbf{E}\left[h_{[\hat{S}]}^{T}A_{j:}^{T}A_{j:}h_{[\hat{S}]}\right] = \sum_{j=1}^{m} \mathbf{E}\left[\left(\sum_{i \in C_{j} \cap \hat{S}} h^{(i)}A_{ji}\right)^{2}\right]$$

$$\stackrel{(57)}{=} \sum_{j=1}^{m} z_{j}^{T}P(C_{j} \cap \hat{S})z_{j} \stackrel{(52)}{\leq} \sum_{j=1}^{m} \lambda(C_{j} \cap \hat{S})\left(z_{j}^{T} \operatorname{Diag}(P(C_{j} \cap \hat{S}))z_{j}\right)$$

$$\stackrel{(37)}{=} \sum_{j=1}^{m} \lambda(C_{j} \cap \hat{S})\sum_{i \in C_{j}} p_{i}(h^{(i)}A_{ji})^{2} = \sum_{j=1}^{m} \lambda(C_{j} \cap \hat{S})\sum_{i=1}^{n} p_{i}(h^{(i)}A_{ji})^{2}$$

$$= \sum_{i=1}^{n} p_{i}(h^{(i)})^{2}\sum_{j=1}^{m} \lambda(C_{j} \cap \hat{S})A_{ji}^{2} = \sum_{i=1}^{n} p_{i}(h^{(i)})^{2}v_{i}.$$

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#### Useful bounds on $\lambda(\hat{S})$

#### Theorem 28 ([14])

Let  $\hat{S}$  be an arbitrary sampling.

- 1. Lower bound. If  $\hat{S}$  is not nill, then  $\frac{\mathbf{E}[|\hat{S}|^2]}{\mathbf{E}[|\hat{S}|]} \leq \lambda(\hat{S})$ .
- 2. **Upper bound.** If  $|\hat{S}| \leq \tau$  with probability 1, then  $\lambda(\hat{S}) \leq \tau$ .
- 3. **Identity.** If  $|\hat{S}| = \tau$  with probability 1, then  $\lambda(\hat{S}) = \tau$ .

Let us apply the 2nd part of the above theorem to the sampling  $J \cap \hat{S}$ :

#### Corollary 29

Let  $\hat{S}$  be an arbitrary sampling,  $J \subseteq [n]$  and c a constant such that  $|J \cap \hat{S}| \le c$  with probability 1. Then

$$\lambda(J\cap \hat{S})\leq c.$$

In particular, if  $|\hat{S}| \leq \tau$  with probability 1, then  $|J \cap \hat{S}| \leq \min\{|J|, \tau\}$  with probability 1, and hence  $\lambda(J \cap \hat{S}) \leq \min\{|J|, \tau\}$ .

Remark: The above corollary is useful as we can apply it in connection with Theorem 27 with  $J = C_j$  for j = 1, 2, ..., m.



#### Computing $\lambda(J \cap \hat{S})$ : Product Sampling

#### Example 30 (Product Sampling)

Assume that the sets  $\{C_j\}$  in Model 1 form a partition of [n]. The consider the sampling  $\hat{S}$  defined as follows:

$$\mathbf{P}(\hat{S} = S) = \begin{cases} (\prod_{j=1}^{m} |C_j|)^{-1}, & S \in C_1 \times C_2 \times \cdots \times C_m, \\ 0, & \text{otherwise.} \end{cases}$$

Note that  $|C_j \cap \hat{S}| = 1$  with probability 1, and hence by Corollary 29,

$$\lambda(C_j \cap \hat{S}) \leq 1.$$

On the other hand, by the first part of Theorem 28,  $\lambda(C_j \cap \hat{S}) \geq 1$ , and hence this sampling achieves the smallest possible value of the " $\lambda$  parameters" in (56) (which is "good" as other things equal, ESO with small  $\{v_i\}$  is better). Let us remark that  $\mathbf{E}[|\hat{S}|] = m$ .



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#### Computing $\lambda(J \cap \hat{S})$ : $\tau$ -Nice Sampling

#### Exercise 5 ( $\tau$ -Nice Sampling)

Show by direct computation that if  $\hat{S}$  is a  $\tau$ -nice sampling, then the lower bound in part 1 of Theorem 28 is attained for  $C_i \cap \hat{S}$  for all j:

$$\lambda(C_j \cap \hat{S}) = \frac{\mathbf{E}[|C_j \cap \hat{S}|^2]}{\mathbf{E}[|C_i \cap \hat{S}|]} \stackrel{(34)+(28)}{=} 1 + \frac{(\omega_j - 1)(\tau - 1)}{\max\{n - 1, 1\}}, \quad (58)$$

where  $\omega_i = |C_i|$ .



# Computing $\lambda(J \cap \hat{S})$ : Distributed $\tau$ -Nice Sampling - Part I

## Exercise 6 (Distributed $\tau$ -Nice Sampling; [14])

Show that if  $\hat{S}$  is the distributed  $\tau$ -nice sampling, then

$$\lambda(C_{j} \cap \hat{S}) \leq \underbrace{1 + \frac{(\tau - 1)(\omega_{j} - 1)}{s_{1}}}_{\lambda_{1,j}} + \underbrace{\left(\frac{\tau}{s} - \frac{\tau - 1}{s_{1}}\right) \frac{\omega_{j}' - 1}{\omega_{j}'}}_{\lambda_{2,j}} \omega_{j}, \quad (59)$$

where  $s_1 = \max\{1, s-1\}$ ,  $\omega_j = |C_j|$ , and  $\omega'_j$  is the number of partitions "active" at row j of A:

$$\omega_{j}' \stackrel{\text{def}}{=} |\{I : A_{ji} \neq 0 \text{ for some } i \in \mathcal{P}_{I}\}|.$$

#### Exercise 7

Show that if the number of partitions is 1 (c = 1), bound (59) for the distributed  $\tau$ -nice sampling specializes to the bound (58) for the  $\tau$ -nice sampling.



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# Computing $\lambda(J \cap \hat{S})$ : Distributed $\tau$ -Nice Sampling - Part II

# Lemma 31 ([14])

Consider the distributed  $\tau$ -nice sampling. Suppose  $\tau \geq 2$ . For any  $1 \leq \eta \leq s$  the following holds:

$$\left(\frac{\tau}{s} - \frac{\tau - 1}{s - 1}\right) \eta \leq \frac{1}{\tau - 1} \left(1 + \frac{(\tau - 1)(\eta - 1)}{s - 1}\right).$$

Note that Lemma 31 implies that

$$\lambda_{1,j} + \lambda_{2,j} \le \left(1 + \frac{1}{\tau - 1}\right) \lambda_{1,j}. \tag{60}$$

# Distributed NSync: Cost of Distribution

Assume f is 1-strongly convex, and consider running NSync with the distributed  $\tau$ -nice sampling. Then  $p_i = \frac{\mathbf{E}[\hat{S}]}{n} = \frac{\tau c}{sc} = \frac{\tau}{s}$  and hence the leading term in the complexity bound is

$$\Lambda = \max_{i} \frac{v_{i}}{p_{i}} \stackrel{\text{(56)}}{=} \max_{i} \frac{s \sum_{j=1}^{m} \lambda(C_{j} \cap \hat{S})}{\tau} \stackrel{\text{(60)}}{\leq} \max_{i} \frac{s \sum_{j=1}^{m} (\lambda_{1,j} + \lambda_{2,j}) A_{ji}^{2}}{\tau} \stackrel{\text{def}}{=} \Lambda'.$$

- Notice that the effect of partitioning on complexity comes only through  $\lambda_{2,j}$ .
- ▶ Define a new quantity that does not depend on partitioning:

$$\Lambda'' = \max_{i} \frac{s \sum_{j=1}^{m} \lambda_{1,j} A_{ji}^{2}}{\tau}$$

and notice that (60) implies that

$$\Lambda'' \le \Lambda' \le (1 + \frac{1}{\tau - 1})\Lambda''$$

This means that:

Theorem 32 (Cost of Distribution: compare with [10, 14])

If  $\tau \geq 2$ , the worst-case partitioning is at most  $(1 + \frac{1}{\tau})$  times worse than the optimal partitioning, in terms of the number of iterations of NSync.



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# Proof of Theorem 28 - Part I

**Point 1.** For simplicity of notation, put  $P = P(\hat{S})$ . If we choose  $\theta \in \mathbb{R}^n$  with  $\theta_i = (\text{Tr}(P))^{-1/2}$  for all i, we get  $\theta^T D^P \theta = \sum_i P_{ii} \theta_i^2 = 1$  and hence

$$\lambda(\hat{S}) \overset{(52)}{\geq} \theta^T P \theta \overset{(57)}{=} \mathbf{E} \Big[ \Big( \sum_{i \in \hat{S}} \theta_i \Big)^2 \Big] = \frac{\mathbf{E} \Big[ \Big( \sum_{i \in \hat{S}} 1 \Big)^2 \Big]}{\mathsf{Tr}(P)} \overset{(22)}{=} \frac{\mathbf{E}[|\hat{S}|^2]}{\mathbf{E}[|\hat{S}|]}.$$

**Point 2.** Let us represent  $\hat{S}$  as a convex combination of elementary samplings:  $\hat{S} = \sum_{S \subseteq [n]} q_S \hat{E}_S$ , where  $q_S = \mathbf{P}(\hat{S} = S)$ . Note that then we also have

$$P(\hat{S}) = \sum_{S \subseteq [n]} q_S P(\hat{E}_S) \stackrel{(52)}{=} \sum_{S \subseteq [n]} q_S e_S e_S^T.$$
 (61)



## Proof of Theorem 28 - Part II

Since  $|\hat{S}| \le \tau$  with probability 1, we have  $|S| \le \tau$  whenever  $q_S > 0$ . For any  $\theta \in \mathbb{R}^n$  we can now estimate:

$$\theta^{T} P(\hat{S}) \theta \stackrel{\text{(61)}}{=} \sum_{S:q_{S}>0} q_{S} (e_{S}^{T} \theta)^{2} \leq \sum_{S:q_{S}>0} q_{S} \|e_{S}\|^{2} \sum_{i \in S} \theta_{i}^{2}$$

$$\stackrel{\text{(54)}}{=} \sum_{S:q_{S}>0} q_{S} |S| \sum_{i \in S} \theta_{i}^{2}$$

$$\leq \tau \sum_{S:q_{S}>0} q_{S} \theta^{T} \operatorname{Diag}(e_{S} e_{S}^{T}) \theta$$

$$= \tau \theta^{T} \left( \sum_{S:q_{S}>0} q_{S} \operatorname{Diag}(e_{S} e_{S}^{T}) \right) \theta$$

$$\stackrel{\text{(61)}}{=} \tau \left( \theta^{T} \operatorname{Diag}(P(\hat{S})) \theta \right).$$

We thus see that  $\lambda(\hat{S}) \leq \tau$ .



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# Proof of Theorem 28 - Part III

Point 3. The result follows by combining the upper and lower bounds. Alternatively, we can see this by inspecting the derivation in part 2. Indeed, if  $|\hat{S}| = \tau$  with probability 1, then  $|S| = \tau$  whenever  $q_S > 0$ , and hence the second inequality in point 2 above is an equality. By choosing  $\theta_i = \alpha$  for any constant  $\alpha$ , the first inequality turns into an equality (this is because we then have equality in the Cauchy-Schwartz inequality  $e_S^T \theta \leq \|e_S\|^2 \sum_{i \in S} \theta_i^2$  for all S).



# ESO( $f \sim \text{Model 3}$ , $\hat{S} \sim \tau$ -nice)

#### Theorem 33

Let f satisfy assumptions in Model 3 and  $\hat{S}$  be a  $\tau$ -nice sampling. Then for all  $x, h \in \mathbb{R}^N$ ,

$$\mathbf{E}\left[f(x+h_{[\hat{S}]})\right] \le f(x) + \frac{\tau}{n}\left(\langle \nabla f(x), h \rangle + \frac{1}{2}\|h\|_{\nu}^{2}\right), \tag{62}$$

where

$$v_i \stackrel{\text{def}}{=} \sum_{j=1}^m \beta_j L_{ji} = \sum_{j:i \in C_j} \beta_j L_{ji}, \qquad i = 1, 2, \dots, n,$$

$$(63)$$

$$eta_j \stackrel{ ext{def}}{=} 1 + rac{(\omega_j - 1)(\tau - 1)}{\max\{1, n - 1\}}, \qquad j = 1, 2, \dots, m.$$

That is,  $(f, \hat{S}) \sim ESO(v)$ .



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## Proof of Theorem 33 - Part I

▶ We first claim that for all *j*,

$$\mathbf{E}\left[f_{j}(x+h_{[\hat{S}]})\right] \leq f_{j}(x) + \frac{\tau}{n}\left(\langle \nabla f_{j}(x), h \rangle + \frac{\beta_{j}}{2} \|h\|_{L_{j:}}^{2}\right), \qquad (64)$$

where  $L_{j:} = (L_{j1}, \ldots, L_{jn}) \in \mathbb{R}^n$ . That is,  $(f_j, \hat{S}) \sim ESO(\beta_j L_{j:})$ . Equation (62) then follows by adding up the inequalities (64) for all j. In the rest we prove the claim.

▶ A well known consequence of (46) is that for all  $x \in \mathbb{R}^N$ ,  $t \in \mathbb{R}^{N_i}$ ,

$$f_j(x+U_it) \leq f_j(x) + \langle \nabla_i f_j(x), t \rangle + \frac{L_{ji}}{2} ||t||_{(i)}^2.$$
 (65)



#### Proof of Theorem 33 - Part II

We fix x and define

$$\hat{f}_j(h) \stackrel{\text{def}}{=} f_j(x+h) - f_j(x) - \langle \nabla f_j(x), h \rangle. \tag{66}$$

Since

$$\mathbf{E}\left[\hat{f}_{j}(h_{[\hat{S}]})\right] \stackrel{(66)}{=} \mathbf{E}\left[f_{j}(x+h_{[\hat{S}]})-f_{j}(x)-\langle\nabla f_{j}(x),h_{[\hat{S}]}\rangle\right]$$

$$\stackrel{(29)}{=} \mathbf{E}\left[f_{j}(x+h_{[\hat{S}]})\right]-f_{j}(x)-\frac{\tau}{n}\langle\nabla f_{j}(x),h\rangle,$$

it now only remains to show that

$$\mathbf{E}\left[\hat{f}_{j}(h_{[\hat{S}]})\right] \leq \frac{\tau\beta_{j}}{2n} \|h\|_{L_{j:}}^{2}. \tag{67}$$

We now adopt the convention that expectation conditional on an event which happens with probability 0 is equal to 0. Let  $\eta_i \stackrel{\text{def}}{=} |C_i \cap \hat{S}|$ , and using this convention, we can write

$$\mathbf{E}\left[\hat{f}_{j}(h_{[\hat{S}]})\right] = \sum_{k=0}^{n} \mathbf{P}(\eta_{j} = k) \mathbf{E}\left[\hat{f}_{j}(h_{[\hat{S}]}) \mid \eta_{j} = k\right]. \tag{68}$$



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## Proof of Theorem 33 - Part III

▶ For any  $k \ge 1$  for which  $\mathbf{P}(\eta_j = k) > 0$ , we now use use convexity of  $\hat{f}_i$  to write

$$\mathbf{E} \left[ \hat{f}_{j}(h_{[\hat{S}]}) \mid \eta_{j} = k \right] = \mathbf{E} \left[ \hat{f}_{j} \left( \frac{1}{k} \sum_{i \in C_{j} \cap \hat{S}} k U_{i} h^{(i)} \right) \mid \eta_{j} = k \right]$$

$$\leq \mathbf{E} \left[ \frac{1}{k} \sum_{i \in C_{j} \cap \hat{S}} \hat{f}_{j} \left( k U_{i} h^{(i)} \right) \mid \eta_{j} = k \right]$$

$$\stackrel{(35)}{=} \frac{1}{\omega_{j}} \sum_{i \in C_{j}} \hat{f}_{j} \left( k U_{i} h^{(i)} \right)$$

$$\stackrel{(65)+(66)}{\leq} \frac{1}{\omega_{j}} \sum_{i \in C_{j}} \frac{L_{ji}}{2} ||kh^{(i)}||_{(i)}^{2} = \frac{k^{2}}{2\omega_{j}} ||h||_{L_{j:}}^{2}. \quad (69)$$



## Proof of Theorem 33 - Part IV

Finally,

$$\mathbf{E} \left[ \hat{f}_{j}(h_{[\hat{S}]}) \right] \stackrel{(68)+(69)}{\leq} \sum_{k} \mathbf{P}(\eta_{j} = k) \frac{k^{2}}{2\omega_{j}} ||h||_{L_{j:}}^{2} \\
= \frac{1}{2\omega_{j}} ||h||_{L_{j:}}^{2} \mathbf{E}[|C_{j} \cap \hat{S}|^{2}] \\
\stackrel{(34)}{=} \frac{\tau \beta_{j}}{2n} ||h||_{L_{i:}}^{2},$$

and hence (67) is proved.



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# $\mathsf{DSO}(f \sim \mathsf{Model}\ 3)$

#### Corollary 34

Let f satisfy assumptions in Model 3 and  $\hat{S}$  be a  $\tau$ -nice sampling. Then for all  $x, h \in \mathbb{R}^N$  we have

$$f(x+h) \le f(x) + \langle \nabla f(x), h \rangle + \frac{\bar{\omega}\bar{L}}{2} ||h||_{w}^{2}, \tag{70}$$

where

$$\bar{\omega} \stackrel{\text{def}}{=} \sum_{i} \omega_{j} \frac{\sum_{i} L_{ji}}{\sum_{k,i} L_{ki}}, \quad \bar{L} \stackrel{\text{def}}{=} \frac{\sum_{ji} L_{ji}}{n}, \quad w_{i} \stackrel{\text{def}}{=} \frac{n}{\sum_{j,i} \omega_{j} L_{ji}} \sum_{i} \omega_{j} L_{ji}. \quad (71)$$

Note that  $\bar{\omega}$  is a data-weighted average of the values  $\{\omega_j\}$  and that  $\sum w_i = n$ .

#### Proof.

This follows from Theorem 33 used with  $\tau = n$  (notice that  $\bar{\omega}\bar{L}w = v$ ).



# ESO and Lipschitz Continuity I

We will now study the collection of functions  $\hat{\phi}_x : \mathbb{R}^N \to \mathbb{R}$  for  $x \in \mathbb{R}^N$  defined by

$$\hat{\phi}_{x}(h) \stackrel{\text{def}}{=} \mathbf{E} \left[ \phi(x + h_{[\hat{S}]}) \right].$$
 (72)

Let us first establish some basic connections between  $\phi$  and  $\hat{\phi}_x$ .

## Lemma 35 ([9])

Let  $\hat{S}$  be any sampling and  $\phi: \mathbb{R}^N \to \mathbb{R}$  any function and  $x \in \mathbb{R}^N$ . Then

- (i) if  $\phi$  is convex, so is  $\hat{\phi}_x$ ,
- (ii)  $\hat{\phi}_x(0) = \phi(x)$ ,
- (iii) If  $\hat{S}$  is proper and uniform, and  $\phi: \mathbb{R}^N \to \mathbb{R}$  is continuously differentiable, then

$$\nabla \hat{\phi}_{x}(0) = \frac{\mathbf{E}[|\hat{S}|]}{n} \nabla \phi(x).$$



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## Proof of Lemma 35

Fix  $x \in \mathbb{R}^N$ . Notice that

$$\hat{\phi}_{x}(h) = \mathbf{E}[\phi(x + h_{[\hat{S}]})] = \sum_{S \subseteq [n]} \mathbf{P}(\hat{S} = S)\phi(x + U_{S}h),$$

where

$$U_S \stackrel{\text{def}}{=} \sum_{i \in S} U_i U_i^T$$
.

As  $\hat{\phi}_x$  is a convex combination of convex functions, it is convex, establishing (i). Property (ii) is trivial. Finally,

$$\nabla \hat{\phi}_{x}(0) = \mathbf{E}\left[\nabla \phi(x + h_{[\hat{S}]})\Big|_{h=0}\right] = \mathbf{E}\left[U_{\hat{S}}\nabla \phi(x)\right] = \mathbf{E}\left[U_{\hat{S}}\right]\nabla \phi(x) = \frac{\mathbf{E}[|\hat{S}|]}{n}\nabla \phi(x).$$

The last equality follows from the observation that  $U_{\hat{S}}$  is an  $N \times N$  binary diagonal matrix with ones in positions (v,v) for coordinates  $v \in \{1,2,\ldots,N\}$  belonging to blocks  $i \in \hat{S}$  only, coupled with the fact that for uniform samplings,  $p_i = \mathbf{E}[|\hat{S}|]/n$ .



# ESO and Lipschitz Continuity II

We now establish a connection between ESO and a uniform bound in x on the Lipschitz constants of the gradient "at the origin" of the functions  $\{\hat{\phi}_x, \ x \in \mathbb{R}^N\}$ .

#### Theorem 36

Let  $\hat{S}$  be proper and uniform, and  $\phi : \mathbb{R}^N \to \mathbb{R}$  be continuously differentiable. Then the following statements are equivalent:

(i) 
$$(\phi, \hat{S}) \sim ESO(v)$$
,

(ii) 
$$\hat{\phi}_x(h) \leq \hat{\phi}_x(0) + \langle \nabla \hat{\phi}_x(0), h \rangle + \frac{1}{2} \frac{\mathsf{E}[|\hat{S}|]}{n} ||h||_v^2, \qquad x, h \in \mathbb{R}^N.$$

#### Proof.

We only need to substitute (72) and Lemma 35(ii-iii) into inequality (ii) and compare the result with the definition of ESO (5).



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# Lecture 6 APPROX



## The Problem

We are interested in solving the following optimization problem:

$$\min_{\mathbf{x} \in \mathbb{R}^N} f(\mathbf{x}) + \psi(\mathbf{x}), \tag{73}$$

where

- f is a "smooth" convex function (to be made precise later),
- $\blacktriangleright \psi$  is block separable:

$$\psi(x) = \sum_{i=1}^{n} \psi_i(x^{(i)}), \tag{74}$$

where  $\psi_i : \mathbb{R}^{N_i} \to \mathbb{R} \cup \{+\infty\}$  are convex and closed.



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# **Examples of Regularizers**

**▶** Smooth optimization:

$$\psi(x) \equiv 0$$

▶ Box constraints: Let  $X_i \subseteq \mathbb{R}^{N_i}$  be closed convex sets and

$$\psi(x) = \begin{cases} 0, & x^{(i)} \in X_i & \text{for all} & i \in [n] \\ +\infty, & \text{otherwise.} \end{cases}$$

► L2/Ridge:

$$\psi(x) = \lambda ||x||_2^2$$

► L1/LASSO:

$$\psi(x) = \lambda ||x||_1$$

► Group LASSO:

$$\psi(x) = \sum_{i=1}^{n} \|x^{(i)}\|_{2}$$

All are block separable and convex.



# APPROX algorithm - Version 1

- 1: Choose  $x_0 \in \text{dom } \psi$  and set  $z_0 = x_0$  and  $\theta_0 > 0$
- 2: for  $k \ge 0$  do
- 3:  $y_k = (1 \theta_k)x_k + \theta_k z_k$
- 4: Generate a random set of blocks  $S_k \sim \hat{S}$
- 5:  $z_{k+1} = z_k$
- 6: for  $i \in S_k$  do
- 7:  $z_{k+1}^{(i)} = \operatorname{arg\,min}_{z \in \mathbb{R}^{N_i}} \left\{ \langle \nabla_i f(y_k), z \rangle + \frac{\theta_k v_i}{2p_i} \|z z_k^{(i)}\|_{(i)}^2 + \psi_i(z) \right\}$
- 8: **end for**
- 9:  $x_{k+1} = y_k + \theta_k(z_{k+1} z_k) \bullet p^{-1}$
- 10:  $heta_{k+1} = rac{\sqrt{ heta_k^4 + 4 heta_k^2} heta_k^2}{2}$  (fast) or  $heta_{k+1} = heta_k$  (normal)
- 11: end for

Remark 1: Our analysis will follow this version.

Remark 2: The  $\bullet$  product is to be applied block-wise, i.e., for  $a \in \mathbb{R}^N$ :

$$a \bullet p^{-1} = \sum_{i=1}^{n} \frac{1}{p_i} U_i a^{(i)}.$$



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# Reformulation: Change of Variables - Part I

Focusing on the iterates  $x_k, y_k, z_k$  only, the algorithm can schematically be written as follows:

APPROX Schema: Version 1

$$y_k \leftarrow (1 - \theta_k) x_k + \theta_k z_k \tag{75}$$

$$z_{k+1} \leftarrow Procedure(y_k; z_k; S_k)$$
 (76)

$$x_{k+1} \leftarrow y_k + \theta_k (z_{k+1} - z_k) \bullet p^{-1} \tag{77}$$

Consider the change of variables from  $\{x_k, y_k, z_k, \}$  to  $\{z_k, g_k\}$  where

$$g_k = y_k - z_k \tag{78}$$

**Inverse change of variables:** From  $\{z_k, g_k\}$  we can recover  $\{x_k, y_k, z_k\}$  as follows:

$$x_{k+1} \stackrel{(77)+(78)}{=} (z_k + g_k) + \theta_k (z_{k+1} - z_k) \bullet p^{-1}, \quad y_k \stackrel{(78)}{=} z_k + g_k \quad (79)$$



# Reformulation: Change of Variables - Part II

It remains to show that  $g_{k+1}$  can be computed (from g and z):

$$g_{k+1} \stackrel{(78)}{=} y_{k+1} - z_{k+1} \stackrel{(75)}{=} (1 - \theta_{k+1})(x_{k+1} - z_{k+1})$$

$$\stackrel{(79)}{=} (1 - \theta_{k+1})(g_k - (e - \theta_k p^{-1}) \bullet (z_{k+1} - z_k)),$$

where  $e \in \mathbb{R}^n$  is the vector of all ones.

Method (75)–(77) can thus be written in the form:

APPROX Schema: Version 2

$$z_{k+1} \leftarrow Procedure(z_k + g_k; z_k; S_k)$$
 (80)

$$g_{k+1} \leftarrow (1 - \theta_{k+1}) (g_k - (e - \theta_k p^{-1}) \bullet (z_{k+1} - z_k))$$
 (81)



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#### **Historical Notes**

1. "Normal" & uniform. Choose  $\theta_0 = \frac{\mathbf{E}[|\hat{S}|]}{n}$  and  $\theta_k = \theta_0$  for all k and let  $\hat{S}$  be uniform, i.e.,  $p_i = \frac{\mathbf{E}[|\hat{S}|]}{n}$ . Then  $g_k = 0$  for all k and the method simplifies to:

$$z_{k+1} \leftarrow Procedure(z_k; z_k; S_k)$$
 (82)

This is the PCDM method of R. and Takáč [5].

- 2. Fast & uniform. For uniform  $\hat{S}$ , "fast" option in Step 10 and  $\theta_0 = \frac{\mathbf{E}[|\hat{S}|]}{n}$ , this method reduces to the original APPROX method of Fercoq & R. [12].
- 3. Fast & non-uniform. For non-uniform  $\hat{S}$  presented here,  $\theta_0 \leq \min_i p_i$  (and  $\theta_0 \leq 1$  if  $\psi \equiv 0$ ) and for the "fast" option in Step 10, it was analyzed by Qu & R. [14].



# APPROX algorithm – Version 2 (variables $g_k, z_k$ )

In detail, version 2 has the following form:

```
1: Choose x_0 \in \text{dom } \psi and \theta_0 > 0, g_0 = 0 and z_0 = x_0
  2: for k > 0 do
           Generate a random set of blocks S_k \sim \hat{S}
  3:
  4:
           z_{k+1} \leftarrow z_k
          for i \in S_k do t_k^{(i)} =
  5:
  6:
               \operatorname{arg\,min}_{t \in \mathbb{R}^{N_i}} \left\{ \langle \nabla_i f(g_k + z_k), t \rangle + \frac{\theta_k v_i}{2\rho_i} \|t\|_{(i)}^2 + \psi_i (z_k^{(i)} + t) \right\}
               z_{k+1}^{(i)} \leftarrow z_k^{(i)} + t_k^{(i)}
  7:
  8:
          g_{k+1} \leftarrow (1-\theta_{k+1})(g_k-(e-\theta_kp^{-1}) \bullet t_k)
           	heta_{k+1}=rac{\sqrt{	heta_k^4+4	heta_k^2}-	heta_k^2}{2} (fast) or 	heta_{k+1}=	heta_k (normal)
12: OUTPUT: x_{k+1} = (z_k + g_k) + \theta_k(z_{k+1} - z_k) \bullet p^{-1}
```



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# Complexity

## Theorem 37 ([12, 14])

Assume:

- $\{S_k\}_{k\geq 1}$  are iid following the distribution of a proper sampling  $\hat{S}$ ,
- f is convex and  $(f, \hat{S}) \sim ESO(v)$ ,
- lacksquare  $\psi$  is block separable, where  $\psi_i$  are convex and closed.

Let  $x_0 \in \text{dom } F$  and choose  $\theta_0 \in (0, \min_i p_i]$  (if  $\psi = 0$ , choose  $\theta_0 \in (0, 1]$ ). Then for any point y such that  $F(y) \leq F(x_0)$  (and hence also for the optimal point  $x_*$  if such a point exists), the iterates  $\{x_k\}$  of APPROX satisfy

$$\mathbf{E}[F(x_k) - F(y)] \leq \frac{4}{((k-1)\theta_0 + 2)^2} C, \quad k \geq 1$$
 (83)

where

$$C \stackrel{\text{def}}{=} (1 - \theta_0) (F(x_0) - F(y)) + \frac{\theta_0^2}{2} ||x_0 - y||_{p^{-2} \bullet v}^2.$$
 (84)



# Comments: Smooth Case ( $\psi \equiv 0$ )

lacktriangle In the smooth case  $(\psi \equiv 0)$  we may choose  $heta_0 = 1$  and get

$$\mathbf{E}[F(x_k)-F(x_*)] \leq \frac{2\|x_0-x_*\|_{p^{-2}\bullet v}^2}{(k+1)^2} = \frac{2}{(k+1)^2} \sum_{i=1}^n \frac{v_i}{p_i^2} \|x_0^{(i)}-x_*^{(i)}\|_{(i)}^2.$$

▶ If, moreover, we choose uniform sampling  $\hat{S}$  and let  $\tau = \mathbf{E}[|\hat{S}|]$ , then since  $p_i = \frac{\tau}{n}$  for all i, we get

$$\mathbf{E}[F(x_k) - F(x_*)] \leq \frac{2n^2 ||x_0 - x_*||_{\nu}^2}{\tau^2 (k+1)^2}.$$

In other words, the number of iterations for obtaining an  $\epsilon$ -solution (in expectation) does not exceed

$$k = \left[ \frac{\sqrt{2}n\|x_0 - x_*\|_{\nu}}{\tau\sqrt{\epsilon}} - 1 \right]. \tag{85}$$

Note that the bound gets better as the average number of processors  $(\tau)$  increases (with the caveat that  $\nu$  will generally also grow in  $\tau$ , but less so for sparse problems; as ESO predicts).



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# **Analysis**

We shall now prove the Theorem. We first need to establish 4 lemmas.



# Lemma: Properties of the sequence $\theta_k$

In the first lemma we summarize well-known properties of the sequence  $\theta_k$  used in APPROX.

#### Lemma 38

The sequence  $\{\theta_k\}_{k\geq 0}$  defined APPROX, under the FAST option, is decreasing and satisfies

$$0 < \theta_k \le \frac{2}{k + 2/\theta_0} \le 1 \tag{86}$$

and

$$\frac{1 - \theta_{k+1}}{\theta_{k+1}^2} = \frac{1}{\theta_k^2}. (87)$$



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# Lemma: $x_k$ is in the convex hull of $z_0, \ldots, z_k$

#### Lemma 39

Let  $\{x_k, z_k\}_{k\geq 0}$  be the iterates of APPROX; and assume  $0 < \theta_0 \leq \min_i p_i$ . Then for all  $k \geq 0$  we have

$$x_k^{(i)} = \sum_{l=0}^k \gamma_{kl}^{(i)} z_l^{(i)}, \quad i = 1, 2, \dots, n$$
 (88)

where for each i, the coefficients  $\gamma_{k0}^{(i)},\ldots,\gamma_{kk}^{(i)}$  are non-negative and sum to 1. Moreover, the coefficients are defined recursively by setting  $\gamma_{00}^{(i)}=1,\,\gamma_{10}^{(i)}=1-\frac{\theta_0}{p_i},\,\gamma_{11}^{(i)}=\frac{\theta_0}{p_i}$  and for  $k\geq 1$ ,

$$\gamma_{k+1,l}^{(i)} = \begin{cases} (1 - \theta_k) \gamma_{kl}^{(i)}, & l = 0, \dots, k-1, \\ (1 - \theta_k) \gamma_{kk}^{(i)} + \theta_k - \frac{\theta_k}{p_i}, & l = k, \\ \frac{\theta_k}{p_i}, & l = k+1. \end{cases}$$
(89)

Moreover, for all  $k \ge 0$  and  $i \in [n]$ , the following identity holds

$$\gamma_{k+1,k}^{(i)} + \gamma_{k+1,k+1}^{(i)} = (1 - \theta_k)\gamma_{kk}^{(i)} + \theta_k. \tag{90}$$



#### Remarks about Lemma 39

- Note that if  $p_i = p_j$  for all  $i, j \in [n]$  (i.e., if  $\hat{S}$  is a uniform sampling), then  $\gamma_{kl}^{(i)} = \gamma_{kl}^{(j)}$  for all i, j, and hence the lemma says that  $x_k$  is a convex combination of the vectors  $z_0, z_1, \ldots, z_k$ .
- ▶ The lemma is only needed in the nonsmooth case ( $\psi \neq 0$ ).
- ▶ The proof is straightforward of a "follow-your-nose" style.



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# Proof of Lemma 39 - Part I

We proceed by induction in k. Fix any  $i \in [n]$ .

#### Step 1 (Base case).

- Since  $x_0 = z_0$ , we have  $\gamma_{00}^{(i)} = 1$ .
- ▶ Since  $x_1 = y_0 + \theta_0(z_1 z_0) \bullet p^{-1}$  and  $y_0 = x_0$ , we get  $x_1^{(i)} = (1 \frac{\theta_0}{p_i})z_0^{(i)} + \frac{\theta_0}{p_i}z_1^{(i)}$ , whence  $\gamma_{10}^{(i)} = 1 \frac{\theta_0}{p_i}$ ,  $\gamma_{11}^{(i)} = \frac{\theta_0}{p_i}$ .

Note that for each k, the coefficients are nonnegative and sum to one.

**Step 2** (Recursive relation). If the recursive relation (89) holds for some  $k \ge 1$ , then it holds for k + 1:

$$x_{k+1}^{(i)} \stackrel{(Step 9)}{=} y_k^{(i)} + \frac{\theta_k}{p_i} (z_{k+1}^{(i)} - z_k^{(i)})$$

$$\stackrel{(Step 3)}{=} (1 - \theta_k) x_k^{(i)} + \theta_k z_k^{(i)} + \frac{\theta_k}{p_i} (z_{k+1}^{(i)} - z_k^{(i)})$$

$$\stackrel{(88)}{=} (1 - \theta_k) \sum_{l=0}^k \gamma_{kl}^{(i)} z_l^{(i)} + \theta_k z_k^{(i)} + \frac{\theta_k}{p_i} (z_{k+1}^{(i)} - z_k^{(i)})$$

$$= \sum_{l=0}^{k-1} \underbrace{(1 - \theta_k) \gamma_{kl}^{(i)}}_{\gamma_{k+1,l}^{(i)}} z_l^{(i)} + \underbrace{((1 - \theta_k) \gamma_{kk}^{(i)} + \theta_k - \frac{\theta_k}{p_i})}_{\gamma_{k+1,k}^{(i)}} z_k^{(i)} + \underbrace{\frac{\theta_k}{p_i}}_{\gamma_{k+1,k+1}^{(i)}} z_{k+1,k+1}^{(i)}.$$



#### Proof of Lemma 39 - Part II

#### Step 3 (Nonnegativity).

- Since  $0 < \theta_k \le 1$  (because  $\theta_0 \le \min_i p_i \le 1$  and  $\{\theta_k\}$  is a decreasing sequence of positive numbers), we deduce from (89) and, using the inductive non-negativity assumption, that  $\gamma_{k+1,l}^{(i)} \ge 0$  for  $l = 0, \ldots, k-1$ .
- Moreover,

$$\gamma_{k+1,k}^{(i)} \stackrel{(89)}{=} (1 - \theta_k) \gamma_{kk}^{(i)} + \theta_k - \frac{\theta_k}{\rho_i} 
= \theta_k (1 - \gamma_{kk}^{(i)}) + \gamma_{kk}^{(i)} - \frac{\theta_k}{\rho_i} 
\stackrel{(89)}{=} \theta_k (1 - \gamma_{kk}^{(i)}) + \frac{\theta_{k-1} - \theta_k}{\rho_i} > \theta_k (1 - \gamma_{kk}^{(i)}) \ge 0.$$

where the first inequality follows since  $\{\theta_k\}$  is a decreasing sequence, and the last inequality by the inductive hypothesis that  $\gamma_{kl}^{(i)}$ ,  $l=0,1,\ldots,k$  are nonnegative and sum to 1.

• Finally,  $\gamma_{k+1,k+1}^{(i)} = \frac{\theta_k}{p_i} > 0$ .



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## Proof of Lemma 39 - Part III

Step 4 (Unit sum). Finally, we can write

$$\sum_{l=0}^{k+1} \gamma_{k+1,l}^{(i)} = \sum_{l=0}^{k-1} \gamma_{k+1,l}^{(i)} + \gamma_{k+1,k}^{(i)} + \gamma_{k+1,k+1}^{(i)}$$

$$\stackrel{(89)}{=} (1 - \theta_k) \sum_{l=0}^{k-1} \gamma_{kl}^{(i)} + \left( (1 - \theta_k) \gamma_{kk}^{(i)} + \theta_k - \frac{\theta_k}{\rho_i} \right) + \frac{\theta_k}{\rho_i}$$

$$= (1 - \theta_k) \sum_{l=0}^{k} \gamma_{kl}^{(i)} + \theta_k$$

$$= 1,$$

where the last step follows from the inductive hypothesis that  $\{\gamma_{kl}^{(i)}\}$  for  $l=0,1,\ldots,k$  sum to one.



# Lemma: Tseng

Define

$$\begin{split} \tilde{z}_{k+1} & \stackrel{\text{def}}{=} & \arg\min_{z \in \mathbb{R}^N} \left\{ \psi(z) + \langle \nabla f(y_k), z - y_k \rangle + \frac{n\theta_k}{2\tau} \|z - z_k\|_v^2 \right\} \\ & \stackrel{(15)+(74)}{=} & \arg\min_{\substack{z^{(i)} \in \mathbb{R}^{N_i} \\ i \in [n]}} \sum_{i=1}^n \left\{ \psi_i(z^{(i)}) + \langle \nabla_i f(y_k), z^{(i)} - y_k^{(i)} \rangle + \frac{n\theta_k v_i}{2\tau} \|z^{(i)} - z_k^{(i)}\|_{(i)}^2 \right\}. \end{split}$$

From this and the definition of  $z_{k+1}$  in APPROX, we see that

$$z_{k+1}^{(i)} = \begin{cases} \tilde{z}_{k+1}^{(i)}, & i \in S_k \\ z_k^{(i)}, & i \notin S_k. \end{cases}$$
(91)

# Lemma 40 (Property 1 in [1])

Let  $\xi(u) \stackrel{\text{def}}{=} f(y_k) + \langle \nabla f(y_k), u - y_k \rangle + \frac{\theta_k}{2} \|u - z_k\|_{p^{-1} \bullet v}^2$ . Then for any  $y \in \text{dom } \psi$ ,

$$\psi(\tilde{z}_{k+1}) + \xi(\tilde{z}_{k+1}) \le \psi(y) + \xi(y) - \frac{\theta_k}{2} \|y - \tilde{z}_{k+1}\|_{p^{-1} \bullet \nu}^2. \tag{92}$$



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# Lemma: Gradient vs Stochastic Gradient Mapping

We now connect the gradient mapping (producing  $\tilde{z}_{k+1}$ ) and the stochastic block gradient mapping (producing the random vector  $z_{k+1}$ ).

From now on, by  $\mathbf{E}_k$  we denote the expectation with respect to  $S_k$ , conditioned on all history.

# Lemma 41 ([12])

For any  $y \in \mathbb{R}^N$  and  $k \ge 0$ ,

$$\mathbf{E}_{k} \left[ \| z_{k+1} - y \|_{V}^{2} - \| z_{k} - y \|_{V}^{2} \right] = \| \tilde{z}_{k+1} - y \|_{p \bullet V}^{2} - \| z_{k} - y \|_{p \bullet V}^{2}. \tag{93}$$

## Proof of Lemma 41

Let  $\hat{S}$  be any proper sampling and  $a, h \in \mathbb{R}^N$ . Recall the following sampling identities:

$$\mathbf{E}[\|h_{[\hat{S}]}\|_{\nu}^{2}] \stackrel{(25)}{=} \|h\|_{p \bullet \nu}^{2}, \qquad \mathbf{E}[\langle a, h_{[\hat{S}]} \rangle_{\nu}] \stackrel{(24)}{=} \langle a, h \rangle_{p \bullet \nu}. \tag{94}$$

Let  $h = \tilde{z}_{k+1} - z_k$ . In view of (14) and (91), we can write  $h_{[S_k]} = z_{k+1} - z_k$ . Now,

$$\mathbf{E}_{k} \left[ \| z_{k+1} - y \|_{v}^{2} - \| z_{k} - y \|_{v}^{2} \right] = \mathbf{E}_{k} \left[ \| h_{[S_{k}]} \|_{v}^{2} + 2 \langle z_{k} - y, h_{[S_{k}]} \rangle_{v} \right] 
\stackrel{(94)}{=} \| h \|_{p \bullet v}^{2} + 2 \langle z_{k} - y, h \rangle_{p \bullet v} 
= \left( \| \tilde{z}_{k+1} - y \|_{p \bullet v}^{2} - \| z_{k} - y \|_{p \bullet v}^{2} \right).$$



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# Proof of the Main Result (Theorem 37) - Part I

**Step 1** (Bounding f). From the definition of  $y_k$  in the algorithm:

$$\theta_k(y_k - z_k) = (1 - \theta_k)(x_k - y_k).$$
 (95)

Since  $x_{k+1} = y_k + h_{[S_k]}$  with  $h = \theta_k(\tilde{z}_{k+1} - z_k) \bullet \sigma$ , we use ESO and obtain the following bound:

$$\mathbf{E}_{k}[f(x_{k+1})] = \mathbf{E}_{k}[f(y_{k} + h_{[S_{k}]})] \\
\leq f(y_{k}) + \langle \nabla f(y_{k}), h \rangle_{p} + \frac{1}{2} \|h\|_{p \bullet w}^{2} \\
= f(y_{k}) + \theta_{k} \langle \nabla f(y_{k}), \tilde{z}_{k+1} - z_{k} \rangle + \frac{\theta_{k}^{2}}{2} \|\tilde{z}_{k+1} - z_{k}\|_{\sigma \bullet v}^{2} \\
= (1 - \theta_{k})f(y_{k}) - \theta_{k} \langle \nabla f(y_{k}), z_{k} - y_{k} \rangle \\
+ \theta_{k} (f(y_{k}) + \langle \nabla f(y_{k}), \tilde{z}_{k+1} - y_{k} \rangle + \frac{\theta_{k}}{2} \|\tilde{z}_{k+1} - z_{k}\|_{\sigma \bullet v}^{2})$$

$$\stackrel{(95)}{=} (1 - \theta_{k})(f(y_{k}) + \langle \nabla f(y_{k}), \tilde{z}_{k+1} - y_{k} \rangle + \frac{\theta_{k}}{2} \|\tilde{z}_{k+1} - z_{k}\|_{\sigma \bullet v}^{2}).(96)$$



# Proof of the Main Result (Theorem 37) - Part II

**Step 2 (Bounding**  $\psi$  **for "fast**  $\theta_k$ "). By Lemma 39, each block of the vector  $x_k$  is a convex combination of the corresponding blocks of the vectors  $z_0, \ldots, z_k$ . By the convexity of each function  $\psi_i$ , for all  $k \geq 0$  we have

$$\psi_{i}(x_{k}^{(i)}) \stackrel{\text{(88)}}{=} \psi_{i} \left( \sum_{l=0}^{k} \gamma_{kl}^{(i)} z_{l}^{(i)} \right) \leq \sum_{l=0}^{k} \gamma_{kl}^{(i)} \psi_{i}(z_{l}^{(i)}) \stackrel{\text{def}}{=} \alpha_{k}^{i}. \tag{97}$$

Moreover,

$$\psi(x_k) = \sum_{i=1}^n \psi_i(x_k^{(i)}) \stackrel{(97)}{\leq} \sum_{i=1}^n \alpha_k^i \stackrel{\text{def}}{=} \hat{\psi}_k.$$
 (98)



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# Proof of the Main Result (Theorem 37) - Part III

Then, for all  $k \ge 0$  and  $i \in \{1, ..., n\}$ , we have:

$$\mathbf{E}_{k}[\alpha_{k+1}^{i}] \stackrel{(97)+(89)}{=} \mathbf{E}_{k} \left[ \sum_{l=0}^{k} \gamma_{k+1,l}^{(i)} \psi_{i}(z_{l}^{(i)}) + \frac{\theta_{k}}{p_{i}} \psi_{i}(z_{k+1}^{(i)}) \right] \\
= \sum_{l=0}^{k} \gamma_{k+1,l}^{(i)} \psi_{i}(z_{l}^{(i)}) + \frac{\theta_{k}}{p_{i}} \mathbf{E}_{k}[\psi_{i}(z_{k+1}^{(i)})] \\
\stackrel{(91)}{=} \sum_{l=0}^{k} \gamma_{k+1,l}^{(i)} \psi_{i}(z_{l}^{(i)}) + \frac{\theta_{k}}{p_{i}} (p_{i} \psi_{i}(\tilde{z}_{k+1}^{(i)}) + (1 - p_{i}) \psi_{i}(z_{k}^{(i)})) \\
= \sum_{l=0}^{k} \gamma_{k+1,l}^{(i)} \psi_{i}(z_{l}^{(i)}) + (\frac{1}{p_{i}} - 1) \theta_{k} \psi_{i}(z_{k}^{(i)}) + \theta_{k} \psi_{i}(\tilde{z}_{k+1}^{(i)}) \\
\stackrel{(89)}{=} (1 - \theta_{k}) \sum_{l=0}^{k-1} \gamma_{kl}^{(i)} \psi_{i}(z_{l}^{(i)}) + (\gamma_{k+1,k}^{(i)} + \gamma_{k+1,k+1}^{(i)} - \theta_{k}) \psi_{i}(z_{k}^{(i)}) + \theta_{k} \psi_{i}(\tilde{z}_{k+1}^{(i)}) \\
\stackrel{(89)}{=} (1 - \theta_{k}) \sum_{l=0}^{k-1} \gamma_{kl}^{(i)} \psi_{i}(z_{l}^{(i)}) + (\gamma_{k+1,k}^{(i)} + \gamma_{k+1,k+1}^{(i)} - \theta_{k}) \psi_{i}(z_{k}^{(i)}) + \theta_{k} \psi_{i}(\tilde{z}_{k+1}^{(i)}) \\
\stackrel{(90)}{=} (1 - \theta_{k}) \sum_{l=0}^{k} \gamma_{kl}^{(i)} \psi_{i}(z_{l}^{(i)}) + \theta_{k} \psi_{i}(\tilde{z}_{k+1}^{(i)}) \\
\stackrel{(97)}{=} (1 - \theta_{k}) \alpha_{k}^{i} + \theta_{k} \psi_{i}(\tilde{z}_{k+1}^{(i)}). \tag{99}$$

# Proof of the Main Result (Theorem 37) - Part IV

Finally,

$$\mathbf{E}_{k}[\hat{\psi}_{k+1}] \stackrel{(98)}{=} \mathbf{E}_{k} \left[ \sum_{i=1}^{n} \alpha_{k+1}^{i} \right]$$

$$= \sum_{i=1}^{n} \mathbf{E}_{k}[\alpha_{k+1}^{i}]$$

$$\stackrel{(99)}{=} \sum_{i=1}^{n} (1 - \theta_{k}) \alpha_{k}^{i} + \theta_{k} \psi_{i}(\tilde{z}_{k+1}^{(i)})$$

$$\stackrel{(98)}{=} (1 - \theta_{k}) \hat{\psi}_{k} + \theta_{k} \psi(\tilde{z}_{k+1}). \tag{100}$$



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# Proof of the Main Result (Theorem 37) - Part V

**Step 3 (Recursion).** For all  $k \ge 0$  define:

$$\hat{F}_k \stackrel{\text{def}}{=} \hat{\psi}_k + f(x_k), \tag{101}$$

and bound the expectation of  $\hat{F}_{k+1}$  as follows:

$$\mathbf{E}_{k}[\hat{F}_{k+1}] \stackrel{\text{(101)}}{=} \mathbf{E}_{k}[\hat{\psi}_{k+1} + f(x_{k+1})] \\
\stackrel{\text{(102)}}{=} (1 - \theta_{k})\hat{\psi}_{k} + \theta_{k}\psi(\tilde{z}_{k+1}) + \mathbf{E}_{k}[f(x_{k+1})] \\
\stackrel{\text{(96)}}{\leq} (1 - \theta_{k})\hat{\psi}_{k} + (1 - \theta_{k})(f(y_{k}) + \langle \nabla f(y_{k}), x_{k} - y_{k} \rangle) \\
+ \theta_{k}(\psi(\tilde{z}_{k+1}) + f(y_{k}) + \langle \nabla f(y_{k}), \tilde{z}_{k+1} - y_{k} \rangle + \frac{\theta_{k}}{2} \|\tilde{z}_{k+1} - z_{k}\|_{p^{-1} \bullet v}^{2}) \\
\stackrel{\text{(92)}}{\leq} (1 - \theta_{k})\hat{\psi}_{k} + (1 - \theta_{k})(f(y_{k}) + \langle \nabla f(y_{k}), x_{k} - y_{k} \rangle) \\
+ \theta_{k}(\psi(y) + f(y_{k}) + \langle \nabla f(y_{k}), y - y_{k} \rangle + \frac{\theta_{k}}{2} \|y - z_{k}\|_{p^{-1} \bullet v}^{2} \\
- \frac{\theta_{k}}{2} \|y - \tilde{z}_{k+1}\|_{p^{-1} \bullet v}^{2}) \\
\leq (1 - \theta_{k})\hat{\psi}_{k} + (1 - \theta_{k})f(x_{k}) \\
+ \theta_{k}(\psi(y) + f(y) + \frac{\theta_{k}}{2} \|y - z_{k}\|_{p^{-1} \bullet v}^{2} - \frac{\theta_{k}}{2} \|y - \tilde{z}_{k+1}\|_{p^{-1} \bullet v}^{2}) \\
= (1 - \theta_{k})\hat{F}_{k} + \theta_{k}F(y) + \frac{\theta_{k}^{2}}{2} \mathbf{E}_{k}[\|y - z_{k}\|_{p^{-2} \bullet v}^{2} - \|y - z_{k+1}\|_{p^{-1} \bullet v}^{2}). \\
\stackrel{\text{(??)}}{=} (1 - \theta_{k})\hat{F}_{k} + \theta_{k}F(y) + \frac{\theta_{k}^{2}}{2} \mathbf{E}_{k}[\|y - z_{k}\|_{p^{-2} \bullet v}^{2} - \|y - z_{k+1}\|_{p^{-2} \bullet v}^{2}]. \\
(102)$$

# Proof of the Main Result (Theorem 37) - Part VI

After rearranging (102), using (87), we obtain the recursion:

$$\frac{1-\theta_{k+1}}{\theta_{k+1}^2}\mathbf{E}_k[\hat{F}_{k+1}-F(y)] + \frac{1}{2}\mathbf{E}_k[\|z_{k+1}-y\|_{p^{-2}\bullet v}^2] \leq \frac{1-\theta_k}{\theta_k^2}(\hat{F}_k-F(y)) + \frac{1}{2}\|z_k-y\|_{p^{-2}\bullet v}^2.$$

**Step 4 (Analyzing the recursion).** We now take total expectation in the above inequality and unroll the recurrence:

$$\frac{1-\theta_k}{\theta_k^2} \mathbf{E}[\hat{F}_k - F(y)] + \frac{1}{2} \mathbf{E}[\|z_k - y\|_{p^{-2} \bullet v}^2] \leq \frac{1-\theta_0}{\theta_0^2} (\hat{F}_0 - F(y)) + \frac{1}{2} \|z_0 - y\|_{p^{-2} \bullet v}^2.$$

Hence, for all  $k \geq 1$ ,

$$\mathbf{E}[\hat{F}_{k} - F(y)] \leq \frac{\theta_{k-1}^{2}(1-\theta_{0})}{\theta_{0}^{2}}(\hat{F}_{0} - F(y)) + \frac{\theta_{k-1}^{2}}{2}\|x_{0} - y\|_{p^{-2} \bullet v}^{2} \\
\leq \frac{4}{((k-1)\theta_{0}+2)^{2}}((1-\theta_{0})(F(x_{0}) - F(y)) + \frac{\theta_{0}^{2}}{2}\|x_{0} - y\|_{p^{-2} \bullet v}^{2}).$$



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