Abstract

In this paper, we present a stochastic version of DCA (Difference of Convex functions Algorithm) to solve a class of optimization problems whose objective function is a large sum of non-convex functions and a regularization term. We consider the $\ell_2,0$ regularization to deal with the group variables selection. By exploiting the special structure of the problem, we propose an efficient DC decomposition for which the corresponding stochastic DCA scheme is very inexpensive: it only requires the projection of points onto balls that is explicitly computed. As an application, we applied our algorithm for the group variables selection in multiclass logistic regression. Numerical experiments on several benchmark datasets and synthetic datasets illustrate the efficiency of our algorithm and its superiority over well-known methods, with respect to classification accuracy, sparsity of solution as well as running time.

1. Introduction

We consider the following optimization problem

$$\min \left\{ f(x) = \frac{1}{n} \sum_{i=1}^{n} f_i(x) + \lambda p(x) \right\}, \quad (1)$$

whose objective function $f$ is a large sum of non-convex functions $f_i(x)$ and a regularization term $p(x)$, where $f_i(x)$ corresponds to a criteria to optimize and $\lambda \geq 0$ is a trade-off parameter between the two terms. This model covers a very vast class of problems arising from several fields such as machine learning, signal processing, etc. For instance, least-squares regression, logistic regression problem, etc can be expressed in the form of (1).

Nowadays, the growth of technologies leads to exponential augmentation of large-scale data where the number of both variables and samples are huge. Thus, optimization methods for solving the problem (1) are faced with a great challenge that is the number of samples $n$ can be extremely large. Among existing methods for this problem, stochastic programming has been proved to be suitable thanks to its ability to exploit the advantage of the sum structure of the problem. In (Schmidt et al., 2015), the authors considered a special case of the large-sum problem (1) where $f_i$ are convex and smooth functions and $p$ corresponds to the $\ell_2$ regularization. Stochastic Average Gradient was developed to solve the resulting problem. Reddi et al. (Reddi et al., 2016) developed Proximal Stochastic Gradient method for the case where $f_i$ are smooth (can be non-convex) and $p$ is convex, non-smooth function. Motivated by its success, we will study stochastic programming for solving (1) in order to deal with data having an extremely large number of samples.

On the other hand, in real-world applications such as image processing, microarray analysis, etc. datasets contain a very large number of variables. In such of cases, we are often to face with the problem of redundant and irrelevant variables. Redundant variables contain information already presented by other variables while irrelevant variables do not contain useful information. Variables selection methods that consist of selecting important variables for a considered task, are a popular and efficient way to deal with redundant and irrelevant variables. In this direction, a natural idea is to formulate the variables selection problem as a minimization of the $\ell_0$-norm (or $\|\cdot\|_0$). The sparse optimization has been extensively studied on both theoretical and practical aspects. The readers can refer to Le Thi et al. (Le Thi et al., 2015) for an extensive overview of existing approaches for the minimization of $\ell_0$-norm.

Nevertheless, when the data possesses certain group structures, we are naturally interested in selecting important groups of variables rather than individual ones. For instance, in multi-factor analysis of variance, a factor with several levels may be expressed through a group of dummy variables. In genomic data analysis, the correlations between genes sharing the biological pathway can...
be high. Hence these genes should be considered as a group. Recently, the mixed-norm regularization has been developed for the group variable selection. It consists in using the $\ell_{2,0}$ regularization term. Assume that $x = (x_1, \ldots, x_m) \in \mathbb{R}^n$ is partitioned into $J$ non-overlapping groups $x_{(1)}, \ldots, x_{(J)}$, then the $\ell_{2,0}$-norm of $x$ is defined by $\|x\|_{2,0} = |\{ j \in \{1, \ldots, J\} : \|x_{(j)}\|_2 \neq 0\}|$. Clearly, $\ell_{2,0}$-norm is non-convex that makes the optimization problem involving $\ell_{2,0}$ challenging. Several works have been developed to solve the problem of mixed-norm regularization $\ell_{2,0}$. The first approach, named the group Lasso ($\ell_{2,1}$-norm) (Yuan & Lin, 2006), is closely connected to the $\ell_1$-norm - an approximation of the $\ell_0$-norm (Tibshirani, 1994). This approach was widely used for selecting groups of variables in multi-task learning (Obozinski et al., 2006), multiclass support vector machine (Blondel et al., 2013), principal component analysis (Khan et al., 2013), linear discriminant analysis (Gu et al., 2011), and compressed sensing (Sun et al., 2009), etc. The second approach consists in replacing the $\ell_{2,0}$-norm by a DC (Difference of Convex functions) approximation. In (Wang et al., 2007), the authors used the smoothly clipped absolute deviation (SCAD) approximation and developed a group coordinate descent based algorithm for the sparse linear regression. Later, Huang et al. (Huang et al., 2012) used the minimax concave penalty (MCP) for the same problem. In (Lee et al., 2016), the authors considered both above approximations and developed DC programming and DCA (DC algorithm) based method for the resulting problems. Recently, Phan et al. (Phan et al., 2017) proposed DCA based algorithms for bi-level variable selection using the combination of the $\ell_0$-norm and $\ell_{q,0}$-norm.

Paper’s contribution: In this paper, we aim at developing efficient methods to solve the problem (1) where $n$ is extremely large and $p(x)$ corresponds to $\ell_{2,0}$ regularization (in order to deal with the group variables selection). The large-sum optimization (1) becomes

$$\min \left\{ f(x) = \frac{1}{n} \sum_{i=1}^{n} f_i(x) + \lambda \|x\|_{2,0} \right\}. \quad (2)$$

We assume that $f_i(x)$ is differentiable with $L$-Lipschitz gradient. This assumption is broad enough to cover several applications. Various important problems in machine learning such as Multi-task feature selection, Sparse logistic regression, Minimizing an expected loss in stochastic learning such as Multi-task feature selection, Sparse logistic regression, etc. can be expressed in the form of (2). As we have mentioned above, the $\ell_{2,0}$-norm can be approximated by a convex (e.g. $\ell_{2,1}$-norm) or non-convex function. Using a non-convex approximation will lead to a “harder” optimization problem but it has been proved that non-convex approximations perform better than convex approximations in terms of sparsity (Le Thi et al., 2015). The resulting problem is then reformulated as a DC program and DCA based algorithm will be developed to solve it. We exploit the special structure of the problem to propose an efficient DC decomposition for which the corresponding DCA scheme is very inexpensive: it only requires the projection of points onto balls that is explicitly computed. On the other hand, in order to deal with data having a large number of samples, we present stochastic version DCA. The convergence properties of the proposed algorithm is rigorously studied to show that the convergence is guaranteed with probability one.

As an application of our algorithm, we consider the group variables selection in multiclass logistic regression. We perform an empirical comparison of stochastic DCA with DCA and standard methods on very large synthetic and real-world datasets, and show that the stochastic DCA is efficient in group variable selection ability and classification accuracy as well as running time.

The remainder of the paper is organized as follows. Solution method based on Stochastic DCA for solving (2) is developed in Section 2. In Section 3, we apply the proposed algorithm to the group variables selection in multiclass logistic regression. Finally Section 4 concludes the paper.

2. Solution method via stochastic DCA

2.1. Outline of DC programming and DCA

DC programming and DCA constitute the backbone of smooth/non-smooth non-convex programming and global optimization (Pham Dinh & Le Thi, 1997; 1998; Le Thi & Pham Dinh, 2005). They address the problem of minimizing a DC function on the whole space $\mathbb{R}^n$ or on a closed convex set $\Omega \subset \mathbb{R}^n$. Generally speaking, a standard DC program takes the form:

$$\alpha = \inf \{ F(x) := G(x) - H(x) \mid x \in \mathbb{R}^n \} \quad (P_{dc}),$$

where $G, H$ are lower semi-continuous proper convex functions on $\mathbb{R}^n$. Such a function $F$ is called a DC function, and $G - H$ is a DC decomposition of $F$ while $G$ and $H$ are the DC components of $F$. A DC program with convex constraint $x \in \Omega$ can be equivalently expressed as $(P_{dc})$ by adding the indicator function $\chi_\Omega$ ($\chi_\Omega(x) = 0$ if $x \in \Omega$ and $+\infty$ otherwise) to the first DC component $G$.

The modulus of strong convexity of $\theta$ on $\Omega$, denoted by $\mu(\theta, \Omega)$ or $\mu(\theta)$ if $\Omega = \mathbb{R}^n$, is given by

$$\mu(\theta, \Omega) = \sup\{ \mu \geq 0 : \theta - (\mu/2)\|x\|^2 \text{ is convex on } \Omega \}$$

One says that $\theta$ is strongly convex on $\Omega$ if $\mu(\theta, \Omega) > 0$. For a convex function $\theta$, the subdifferential of $\theta$ at $x_0 \in \text{dom}\theta := \{ x \in \mathbb{R}^n : \theta(x_0) < +\infty \}$, denoted by $\partial\theta(x_0)$, is
Stochastic DCA for the Large-sum of Non-convex Functions Problem

defined by

\[ \partial \theta(x_0) := \{ y \in \mathbb{R}^n : \theta(x) \geq \theta(x_0) + \langle x - x_0, y \rangle, \quad \forall x \in \mathbb{R}^n \}. \]

The subdifferential \( \partial \theta(x_0) \) generalizes the derivative in the sense that \( \theta \) is differentiable at \( x_0 \) if and only if \( \partial \theta(x_0) \equiv \{ \nabla \theta(x_0) \}. \)

A point \( x^* \) is called a critical point of \( G - H \), or a generalized Karush-Kuhn-Tucker point (KKT) of (Pdc) if \( \partial H(x^*) \cap \partial G(x^*) \neq \emptyset. \)

The main idea of DCA is simple: each iteration \( l \) of DCA approximates the concave part \(-H\) by its affine majorization (that corresponds to taking \( v^l \in \partial H(x^l) \)) and then computes \( x^{l+1} \) by solving the resulting convex problem.

\[ \min_{x \in \mathbb{R}^n} \{ G(x) - \langle v^l, x \rangle \}. \]

The sequence \( \{x^l\} \) generated by DCA enjoys the following properties (Pham Dinh & Le Thi, 1997; 1998; Le Thi & Pham Dinh, 2005):

(i) The sequence \( \{F(x^l)\} \) is decreasing;

(ii) If \( F(x^{l+1}) = F(x^l) \), then \( x^l \) is a critical point of (Pdc) and DCA terminates at \( l \)-th iteration.

(iii) If \( \mu(G) + \mu(H) > 0 \) then the series \( \{x_{k+1}^l - x_k^l\} \) converges.

(iv) If the optimal value \( \alpha \) of (Pdc) is finite and the infinite sequence \( \{x^l\} \) is bounded then every limit point of the sequence \( \{x^l\} \) is a critical point of \( G - H \).

2.2. Stochastic DCA for solving the problem (2)

In this section, we introduce a stochastic version of DCA for solving (2) that exploits the structure of objective function \( f \). We consider a family of DC approximations \( \tilde{p}(x) \) of \( \ell_{2,0} \)-norm, defined by

\[ \tilde{p}(x) = \sum_{j=1}^J \eta(\|x(j)\|_2). \]

where \( \eta \) is a non-convex penalty function which includes SCAD, MCP, Capped-\( \ell_1 \), exponential-\( \ell_p \) with \( 0 < p < 1, \ell_p \)- with \( p < 0 \) (see (Le Thi et al., 2015) for more details). \( \tilde{p}(x) \) can be expressed as \( \tilde{p}(x) = \tilde{g}(x) - \tilde{h}(x) \), where

\[ \tilde{g}(x) = \alpha \sum_{j=1}^J \|x(j)\|_2 \quad \text{and} \quad \tilde{h}(x) = \alpha \sum_{j=1}^J \|x(j)\|_2 - \tilde{p}(x). \]

Hence, the approximate problem of (2) can be written as

\[ \min_{x \in \mathbb{R}^n} \left\{ f(x) = \frac{1}{n} \sum_{i=1}^n \left[ f_i(x) + \lambda \tilde{g}(x) - \lambda \tilde{h}(x) \right] \right\}. \] (3)

Each function \( f_i(x) \) can be rewritten as

\[ f_i(x) = \frac{\rho}{2} \|x\|^2 - \left[ \frac{\rho}{2} \|x\|^2 - f_i(x) \right]. \]

Since \( f_i(x) \) is differentiable with \( L \)-Lipschitz gradient, \( \|x\|^2 - f_i(x) \) is strongly convex with \( \rho > L. \) Hence, \( f_i(x) \) is a DC function. Consequently, \( f(x) \) is a DC function with the following DC decomposition

\[ f(x) = g(x) - h(x), \] (4)

where \( g(x) \) and \( h(x) \) are convex functions defined by

\[ g(x) = \frac{\rho}{2} \|x\|^2 + \lambda \tilde{g}(x), \]

\[ h(x) = \frac{\lambda}{\rho} \sum_{i=1}^n h_i(x); \quad h_i(x) = \frac{\rho}{2} \|x\|^2 - f_i(x) + \lambda \tilde{h}(x). \]

DCA for solving (3) amounts to computing two sequences \( \{x^l\} \) and \( \{v^l\} \) such that \( v^l \in \partial h_i(x^l) \) and \( x^{l+1} \) is an optimal solution of the following convex problem

\[ \min \left\{ g(x) - \langle v^l, x \rangle \right\}. \] (5)

The computation of subgradients of \( h \) requires the one of all components \( h_i \). This can be expensive when \( n \) is very large. Hence we propose a stochastic version of DCA in which we only compute the subgradients of a small subset of components \( h_i \). Precisely, at each iteration \( l \), we compute \( v_i^l \in \partial h_i(x^l) \) for \( i \in s_l \) and keep \( v_i^l = v_i^{l-1} \) for \( i \notin s_l \), where \( s_l \subset \{1, ..., n\} \) is a randomly chosen set of index.

The computation of \( v_i^l \in \partial h_i(x^l) \) can be given as \( v_i^l = \rho x_i^l - \nabla f_i(x^l) + y_i^l \), where \( y_i^l \in \lambda \partial h_i(x^l) \) for all \( i \in s_l \). The convex problem (5) take the form

\[ \min \left\{ \lambda \alpha \sum_{j=1}^J \|x(j)\|_2 + \frac{\rho}{2} \|x\|^2 - \langle \frac{1}{n} \sum_{i=1}^n v_i^l, x \rangle \right\}. \] (6)

We observe that the objective of (6) is separable in groups of \( x \), then the solution to this problem can be computed by solving \( J \) independent sub-problems of the same form:

\[ \min \left\{ \lambda \alpha \|x(j)\|_2 + \frac{\rho}{2} \|x\|^2 - \langle v_i^l, x(j) \rangle \right\}, \] (7)

where \( v_i^l = \frac{1}{\rho \alpha} \sum_{i=1}^n (v_i^l)^{(j)} \) for \( j = 1, ..., J \). The solution of (7) can be explicitly computed by

\[ x^{l+1} = \left( \|v_i^l\|_2 - \frac{\alpha}{\rho} \right) + \frac{v_i^l}{\|v_i^l\|_2}. \] (8)

Thus, the stochastic DCA (SDCA) for solving the problem (3) is described in Algorithm 1.

Now we will prove that the convergence properties of SDCA are guaranteed with probability one.
Theorem 1. If $\alpha^* = \inf f(x) > -\infty$ and $|s_l| = b$ for all $l \geq 1$, then SDCA generates the sequence $\{x^l\}$ such that

a) $\{f(x^l)\}$ is the almost sure convergent sequence.

b) $\sum_{l=1}^{\infty} \|x^l - x^{l-1}\|^2$ is almost surely finite and $\lim_{l \to \infty} \|x^l - x^{l-1}\| = 0$ almost surely.

c) Every limit point of $\{x^l\}$ is a critical point of $f$ with probability one.

Proof. a) Let $x^l_i = x^l$ and $y^l_i = y^l$ for $i \in s_l$, $x^l_i = x^{l-1}_i$ and $y^l_i = y^{l-1}_i$ for $i \notin s_l$. We denote $T^l_i(x)$ the function given by

$$T^l_i(x) = \lambda\tilde{g}(x) + \frac{\rho}{2} \|x - x^l_i\|^2 - \langle x - x^l_i, y^l_i - \nabla f_i(x^l_i) \rangle + f_i(x^l_i) - \tilde{h}(x^l_i),$$

and $T^l(x) = \frac{1}{n} \sum_{i=1}^{n} T^l_i(x)$. From the step 2 in Algorithm 1, it follows that $x^l_{i+1} = \arg \min T^l_i(x)$. Hence, we have

$$T^l(x^{l+1}) \leq T^l(x) = T^{l-1}(x) + \frac{1}{n} \sum_{i \in s_l} \left[ f_i(x^l) + \lambda\tilde{p}(x^l) - T^{l-1}(x^l) \right].$$

Let $\mathcal{F}_l$ denote the $\sigma$-algebra generated by the entire history of SDCA up to the iteration $l$, i.e., $\mathcal{F}_0 = \sigma(x^0)$ and $\mathcal{F}_l = \sigma(x^0, ..., x^l, s_0, ..., s_{l-1})$ for all $l \geq 1$. By taking the expectation of the inequality (9) conditioned on $\mathcal{F}_l$, we have

$$\mathbb{E} \left[ T^l(x^{l+1}) | \mathcal{F}_l \right] \leq T^{l-1}(x) - \frac{b}{n} \left[ T^{l-1}(x) - f(x) \right].$$

By the supermartingale convergence theorem, we can conclude that the sequence $\{T^{l-1}(x^l) - \alpha^*\}$ converges almost surely. Moreover, $\sum_{l=1}^{\infty} \left[ T^{l-1}(x^l) - f(x^l) \right] < \infty,$

almost surely and hence $\{f(x^l)\}$ converges almost surely.

b) Since $y^l_i - 1 \in \lambda\partial\tilde{h}(x^{l-1}_i)$, we have

$$\lambda\tilde{h}(x^l) \geq \lambda\tilde{h}(x^{l-1}) + \langle x - x^{l-1}_i, y^l_i - 1 \rangle.$$

Since $f_i(x)$ is a differentiable function with $L$-Lipschitz gradient, we have

$$f_i(x) \leq f_i(x^{l-1}) + \langle x - x^{l-1}_i, \nabla f_i(x^{l-1}) \rangle + \frac{L}{2} \|x - x^{l-1}_i\|^2.$$

Thus, we get that

$$f_i(x) + \lambda\tilde{p}(x) \leq T^{l-1}(x) + \frac{L - \rho}{2} \|x - x^{l-1}_i\|^2.$$

From (9) and (12), we have

$$T^l(x^{l+1}) \leq T^{l-1}(x) - \frac{\rho - L}{2n} \sum_{i \in s_l} \|x^l - x^{l-1}_i\|^2.$$

By taking the expectation of the inequality (13) conditioned on $\mathcal{F}_l$, we have

$$\mathbb{E} \left[ T^l(x^{l+1}) | \mathcal{F}_l \right] \leq T^{l-1}(x) - \frac{b(\rho - L)}{2n^2} \sum_{i=1}^{n} \|x^l - x^{l-1}_i\|^2.$$

By the supermartingale convergence theorem, we conclude that

$$\sum_{l=1}^{\infty} \sum_{i=1}^{n} \|x^l - x^{l-1}_i\|^2 < \infty,$$

almost surely and hence $\lim_{l \to \infty} \|x^l - x^{l-1}_i\| = 0$ almost surely.

c) Assume that there exists a sub-sequence $\{x^k\}$ of $\{x^l\}$ such that $x^k \to x^*$ almost surely. From (14) and (15), we have $\|x^{k+1}_i - x^{k}_i\| \to 0$ almost surely. Without loss of generality, we can suppose that the sub-sequences $y^k_i \to y^*$ almost surely. We note that $y^k_i \in \lambda\partial\tilde{h}(x^k_i)$ and by the closed property of the subdifferential mapping $\partial\tilde{h}$, we have $y^* \in \lambda\partial\tilde{h}(x^*)$ with probability one. It follows from $x^{k+1} \in \arg \min T^k(x)$ that $T^k(x^{k+1}) \leq T^k(x)$. Taking $k \to \infty$, we get that

$$\lambda\tilde{g}(x^*) \leq \lambda\tilde{g}(x) + \frac{\rho}{2} \|x - x^*\|^2 - \langle x - x^*, y^* - \frac{1}{n} \sum_{i=1}^{n} \nabla f_i(x^*) \rangle,$$

is almost surely satisfied for all $x \in \mathbb{R}^m$. Thus, we have

$$y^* - \frac{1}{n} \sum_{i=1}^{n} \nabla f_i(x^*) \in \partial\lambda\tilde{g}(x^*),$$

with probability one. Therefore,

$$y^* \in \left[ \nabla \frac{1}{n} \sum_{i=1}^{n} f_i(x^*) + \partial\lambda\tilde{g}(x^*) \right] \cap \partial\lambda\tilde{h}(x^*),$$

with probability one. This implies that $x^*$ is a critical point of $f$ with probability one and the proof is complete. □
3. Application to Group Variables Selection in Multiclass Logistic Regression

Logistic regression, introduced by D. Cox in 1958 (Cox, 1958), is a popular method in supervised learning. Logistic regression has been successfully applied in various real-life problems such as cancer detection (Kim et al., 2008), medical (Bagley et al., 2001; Subasi & Erglelebi, 2005), social science (King & Zeng, 2001), etc. Especially, logistic regression combined with feature selection has been proved to be suitable for high dimensional problems, for instance, document classification (Genkin et al., 2007) and microarray classification (Liao & Chin, 2007; Kim et al., 2008).

We describe the multiclass logistic regression problem as follows. Let \( W \) be a \( d \times Q \) matrix, where \( d \) and \( Q \) are the number of features and number of classes, respectively. We denote the \( i \)-th column of \( W \) by \( W_i \) and \( b = (b_1, ..., b_Q) \in \mathbb{R}^Q \). In the multiclass logistic classification problem, a new instance \( x^* \) is classified to class \( y^* \) by using the rule \( y^* = \text{arg} \max_k p(Y = k|X = x^*) \), where \( p(Y = y|X = x) \) is the conditional probability defined by

\[
p(Y = y|X = x) = \frac{\exp(b_y + W^T_{i,y}x)}{\sum_{k=1}^{Q} \exp(b_k + W^T_{i,k}x)}.
\]

(18)

Given a training set containing \( n \) instances \( x_i \) and their corresponding labels \( y_i \in \{1, ..., Q\} \), we aim to find \((W,b)\) for which the total probability of the training instances \( x_i \) belonging to its correct classes \( y_i \) is maximized. To estimate \((W,b)\), we maximize the log-likelihood function defined as

\[
\mathcal{L}(W,b) := -\frac{1}{n} \sum_{i=1}^{n} \ell(x_i, y_i, W, b)
\]

(19)

where \( \ell(x_i, y_i, W, b) = -\log p(Y = y_i|X = x_i) \). As mentioned above, to deal with irrelevant and/or redundant variables in high-dimensional data, we use variables selection method. Note that a variable \( j \) is to be removed if and only if all components in the row \( j \) of \( W \) are zero. Therefore, we can consider each row of \( W \) as a group. Denote by \( W_{j,:} \), the \( j \)-th row of the matrix \( W \). The \( \ell_{2,0} \)-norm of \( W \), i.e., the number of non-zero rows of \( W \) is defined by

\[
\|W\|_{2,0} := \{|j \in \{1, ..., d\} : \|W_{j,:}\|_2 \neq 0\} \}.
\]

Hence, the \( \ell_{2,0} \) regularized multiclass logistic regression problem is formulated as

\[
\min_{W,b} \left\{ \frac{1}{n} \sum_{i=1}^{n} \ell(x_i, y_i, W, b) + \lambda \|W\|_{2,0} \right\}.
\]

(20)

Observe that the problem (20) takes the form of (2) where the function \( f_i(W, b) = \ell(x_i, y_i, W, b) \). In this application, we use a non-convex approximation of the \( \ell_{2,0} \)-norm based on the piecewise exponential penalty function. This approximation function has shown its efficiency in several problems, for instance, variables selection in SVM (Bradley & Mangasarian, 1998; Le Thi et al., 2008), semi-supervised support vector machines (Le et al., 2015), sparse multiclass support vector machines (Le Thi & Nguyen, 2017), sparse signal recovery (Le Thi et al., 2015), sparse linear discriminant analysis (Le Thi & Phan, 2016a,b), variables selection in SVM with uncertain data (Le Thi et al., 2014), etc. Using the piecewise exponential penalty function, the corresponding approximate problem of (20) takes the form:

\[
\min_{W,b} \left\{ \frac{1}{n} \sum_{i=1}^{n} f_i(W, b) + \lambda \tilde{p}(W) \right\},
\]

(21)

where \( \tilde{p}(W) = \sum_{j=1}^{d} \eta_0(\|W_{j,:}\|_2) \) with \( \eta_0(t) = 1 - \exp(-\alpha|t|) \). The function \( \tilde{p}(W) \) can be expressed as a DC function:

\[
\tilde{p}(W) = \alpha \sum_{j=1}^{d} \|W_{j,:}\|_2 - \tilde{h}(W),
\]

where \( \tilde{h}(W) = \sum_{j=1}^{d} [-1 + \alpha \|W_{j,:}\|_2 + \exp(-\alpha\|W_{j,:}\|_2)] \).

According to the SDCA scheme in Algorithm 1, at each iteration \( i \), we have to compute \((v_i^l, z_i^l) = \rho(W^l, b^l) - \nabla f_i(W^l, b^l) + (y_i^l, 0)\) for \( i \in s_i \), where

\[
\begin{align*}
\nabla_{b_k} f_i(W^l, b^l) &= \frac{p_k^l(x_i) - \delta_{ky_i}}{\sum_{j=1}^{d} p_k^l(x_i) - \delta_{ky_i}} x_i \\
\nabla_{W_{i,k}} f_i(W^l, b^l) &= (p_k^l(x_i) - \delta_{ky_i}) x_i
\end{align*}
\]

(22)

with \( p_k^l(x_i) = \frac{\exp(b_k^l + (W_{i,k}^l)^T x_i)}{\sum_{k=1}^{d} \exp(b_k^l + (W_{i,k}^l)^T x_i)} \) and \( \delta_{ky_i} = 1 \) if \( k = y_i \) and 0 otherwise. The computation of \( y_i^l \) is given by

\[
y_i^l = \left\{ \begin{array}{ll}
0 & \text{if } \|W_{j,:}^l\|_2 = 0 \\
\frac{\lambda \eta_0(\|W_{j,:}^l\|_2)}{\|W_{j,:}^l\|_2} & \text{otherwise}.
\end{array} \right.
\]

(23)

SDCA for solving (21) is described in Algorithm 2.

3.1. Numerical Experiment

3.1.1. DATASETS

To illustrate the performances of algorithms, we performed numerical tests on real datasets (aloi, covertype, madelon and sensorless) and simulated datasets (sim_1, sim_2 and sim_3). Dataset Aloi is a library of object images while covertype, madelon, sensorless, are taken from the well-known UCI data repository.

We used the same way as proposed in (Witten & Tibshirani, 2011) to generate simulated datasets. In sim_1, features are independent with different means in each class.

\[\text{http://aloi.science.uva.nl/}\]
For sim\(_2\), features also have different means in each class, however they are dependent. The dataset sim\(_3\) has different one-dimensional means in each class with independent features. The procedure for generating simulated datasets is described as follows.

For sim\(_1\): this dataset consists of four classes. The class \(k\) is sampled from the multivariate normal distribution \(N(\mu_k, I)\), where the mean vector \(\mu_k \in \mathbb{R}^{50}\) is given by \(\mu_{kj} = 0.5\) if \(10(k - 1) + 1 \leq j \leq 10k\) and 0 otherwise. We generate 25,000 samples for each class.

For sim\(_2\): we generate a dataset with three classes sampled from the multivariate normal distributions \(N(\mu_k, \Sigma)\), \(k = 1, 2, 3\), where \(\mu_j \in \mathbb{R}^{50}\) is defined by \(\mu_{kj} = 0.4(k - 1)\) if \(j \leq 40\) and 0 otherwise. We use the block diagonal matrix \(\Sigma\) with five blocks of dimension 10 \(\times\) 10 whose element \((j, j')\) is 0.6\(|j - j'|\). 150,000 instances are generated.

For sim\(_3\): we generate a dataset including four classes as follows: \(x_j \in C_k\) then \(x_{ij} \sim N((k - 1)/3, 1)\) if \(j \leq 100, k = 1, 2, 3, 4\) and \(x_{ij} \sim N(0, 1)\) otherwise. We generate 250,000 instances with equal probabilities for each class.

For pre-processing data, we use standardization to scale the data.

3.1.2. COMPARATIVE ALGORITHMS

We compare our algorithm with two algorithms: msl1 and liblinear. msl1 (Vincent & Hansen, 2014) is a coordinate gradient descent algorithm for solving the multiclass logistic regression using \(\ell_{2,1}\) regularization term, i.e., the convex problem \(\min_{W, b} \left\{ \frac{1}{n} \sum_{i=1}^{n} f(x_i, y_i, W, b) + \lambda\|W\|_{2,1} \right\}\). LibLinear (Fan et al., 2008) is a well-known package for solving large-scale problems by using the coordinate descent algorithm. We use the \(\ell_1\)-regularized logistic regression solver of LibLinear to solve the binary logistic regression problem

\[
\min_{w} \left\{ \frac{1}{n} \sum_{i=1}^{n} \log(1 + e^{-y_i w^T x_i}) + \lambda \sum_{j=1}^{d} |w_j| \right\},
\]

and then the one-vs-the-rest strategy is used for the multiclass case.

3.1.3. EXPERIMENT SETTING

The comparison of algorithms are performed in terms of three criteria: classification accuracy on test set, sparsity of solution and running time. Sparsity is computed as the percentage of selected features, where a feature \(j \in \{1, \ldots, d\}\) is considered to be removed if all absolute values of components of row \(W_j\) are smaller than a threshold \(\epsilon = 10^{-8}\).

The cross-validation procedure is used for experiments. We randomly take 80\% of the whole dataset as a training set and the rest is used as test set (20\%). This process is repeated 10 times and we report the mean and standard deviation of each criterion.

We use the early-stopping condition for SDCA. This is a well-known technique in machine learning, especially in stochastic learning which permits to avoid the over-fitting problem. After each epoch, we compute the accuracy based on the validation set, then we stop SDCA if the accuracy is not improved after \(n_{patience} = 5\) epochs. For comparative algorithms, we use their default stopping parameters. We also stop algorithms if they exceed 2 hours of running time in the training process.

For SDCA, we set the trade-off parameter \(\lambda \in \{10^{-4}, 10^{-3}, \ldots, 1\}\) and the parameter for controlling the tightness of zero-norm approximation \(\alpha \in \{0.5, 1, 2, 5\}\). For both LibLinear and msl1, the trade-off parameter is chosen in interval \(\{10^{-3}, \ldots, 10^3\}\).

All experiments are performed on a PC Intel (R) Xeon (R) E5-2630 v2 @2.60 GHz of 32GB RAM.

3.1.4. EXPERIMENT 1: COMPARISON OF SDCA AND DCA

Firstly, we will study the impact of batch size on the quality of solution and the running time of SDCA. The batch size refers to the size of set of index \(s_i\), i.e., the number of components \(h_i\) that are used to compute the subgradients of \(h\) at each iteration (c.f Algorithm 1). In DCA, or full-batch DCA, all components \(h_i\) are used, i.e., \(s_i \equiv \{1, \ldots, n\}\). The Table 1 reports the accuracy and the running time of SDCA as the batch size varies on an arbitrary chosen dataset (sensorless).

We observe that the running time is smallest (1.78s) with batch size equals to 10\% while giving the second best classification accuracy 85.23\%, only 0.07\% smaller than the best one. Hence, we choose the batch size as 10\% through-
### Table 1. Performance of SDCA as batch size varies

<table>
<thead>
<tr>
<th>Batch Size</th>
<th>5%</th>
<th>10%</th>
<th>15%</th>
<th>20%</th>
<th>25%</th>
<th>30%</th>
<th>40%</th>
<th>50%</th>
</tr>
</thead>
<tbody>
<tr>
<td>Time (s)</td>
<td>2.93 ± 0.1</td>
<td>1.78 ± 0.1</td>
<td>2.53 ± 0.1</td>
<td>2.06 ± 0.2</td>
<td>1.54 ± 0.2</td>
<td>2.24 ± 0.3</td>
<td>2.98 ± 0.4</td>
<td>1.94 ± 0.50</td>
</tr>
<tr>
<td>Accuracy (%)</td>
<td>84.62 ± 1.1</td>
<td>85.23 ± 0.7</td>
<td>83.25 ± 1.2</td>
<td>85.05 ± 0.9</td>
<td>82.77 ± 1.1</td>
<td>84.54 ± 1.1</td>
<td>85.30 ± 0.7</td>
<td>82.88 ± 1.1</td>
</tr>
</tbody>
</table>

### Table 2. Comparative results on both simulated and real datasets.

Bold values correspond to best results for each dataset. NA means that the algorithm fails to furnish a result.  $n$, $d$ and $Q$ is the number of instances, the number of dimensions and the number of classes respectively.

<table>
<thead>
<tr>
<th>Dataset</th>
<th>Algorithm</th>
<th>Accuracy (%)</th>
<th>Time (s)</th>
<th>Sparsity (%)</th>
</tr>
</thead>
<tbody>
<tr>
<td></td>
<td></td>
<td>Mean</td>
<td>STD</td>
<td>Mean</td>
</tr>
<tr>
<td>aloi</td>
<td>DCA (Full-batch)</td>
<td>85.05</td>
<td>0.44</td>
<td>2414.72</td>
</tr>
<tr>
<td>($n \times d$) = $(108,000 \times 128)$</td>
<td>SDCA</td>
<td>82.98</td>
<td>0.47</td>
<td>208.67</td>
</tr>
<tr>
<td>$Q = 1,000$</td>
<td>MSGGL</td>
<td>NA</td>
<td>NA</td>
<td>NA</td>
</tr>
<tr>
<td></td>
<td>LibLinear</td>
<td>81.61</td>
<td>0.20</td>
<td>2732.96</td>
</tr>
<tr>
<td>coverytype</td>
<td>DCA (Full-batch)</td>
<td>71.30</td>
<td>0.09</td>
<td>322.98</td>
</tr>
<tr>
<td>($n \times d$) = $(581,012 \times 54)$</td>
<td>SDCA</td>
<td>71.58</td>
<td>0.13</td>
<td>5.86</td>
</tr>
<tr>
<td>$Q = 7$</td>
<td>MSGGL</td>
<td>71.22</td>
<td>0.02</td>
<td>525.49</td>
</tr>
<tr>
<td></td>
<td>LibLinear</td>
<td>71.54</td>
<td>0.19</td>
<td>264.88</td>
</tr>
<tr>
<td>madelon</td>
<td>DCA (Full-batch)</td>
<td>61.35</td>
<td>0.39</td>
<td>5.25</td>
</tr>
<tr>
<td>($n \times d$) = $(2,600 \times 500)$</td>
<td>SDCA</td>
<td>62.60</td>
<td>1.39</td>
<td>0.12</td>
</tr>
<tr>
<td>$Q = 2$</td>
<td>MSGGL</td>
<td>61.48</td>
<td>2.37</td>
<td>0.08</td>
</tr>
<tr>
<td></td>
<td>LibLinear</td>
<td>61.54</td>
<td>2.72</td>
<td>0.08</td>
</tr>
<tr>
<td>sensorless</td>
<td>DCA (Full-batch)</td>
<td>90.21</td>
<td>0.41</td>
<td>34.65</td>
</tr>
<tr>
<td>($n \times d$) = $(58,509 \times 48)$</td>
<td>SDCA</td>
<td>85.11</td>
<td>0.83</td>
<td>11.76</td>
</tr>
<tr>
<td>$Q = 11$</td>
<td>MSGGL</td>
<td>85.06</td>
<td>0.31</td>
<td>216.48</td>
</tr>
<tr>
<td></td>
<td>LibLinear</td>
<td>75.55</td>
<td>0.24</td>
<td>23.92</td>
</tr>
<tr>
<td>sim_1</td>
<td>DCA (Full-batch)</td>
<td>72.11</td>
<td>0.57</td>
<td>24.67</td>
</tr>
<tr>
<td>($n \times d$) = $(100,000 \times 50)$</td>
<td>SDCA</td>
<td>72.24</td>
<td>0.42</td>
<td>1.33</td>
</tr>
<tr>
<td>$Q = 4$</td>
<td>MSGGL</td>
<td>72.33</td>
<td>0.18</td>
<td>214.83</td>
</tr>
<tr>
<td></td>
<td>LibLinear</td>
<td>72.62</td>
<td>0.38</td>
<td>2038.85</td>
</tr>
<tr>
<td>sim_2</td>
<td>DCA (Full-batch)</td>
<td>65.53</td>
<td>3.70</td>
<td>76.04</td>
</tr>
<tr>
<td>($n \times d$) = $(150,000 \times 50)$</td>
<td>SDCA</td>
<td>68.60</td>
<td>0.02</td>
<td>1.74</td>
</tr>
<tr>
<td>$Q = 3$</td>
<td>MSGGL</td>
<td>68.42</td>
<td>0.03</td>
<td>367.29</td>
</tr>
<tr>
<td></td>
<td>LibLinear</td>
<td>66.92</td>
<td>0.02</td>
<td>2.04</td>
</tr>
<tr>
<td>sim_3</td>
<td>DCA (Full-batch)</td>
<td>99.87</td>
<td>0.02</td>
<td>151.72</td>
</tr>
<tr>
<td>($n \times d$) = $(250,000 \times 500)$</td>
<td>SDCA</td>
<td>99.93</td>
<td>0.01</td>
<td>22.83</td>
</tr>
<tr>
<td>$Q = 4$</td>
<td>MSGGL</td>
<td>99.93</td>
<td>0.01</td>
<td>1581.44</td>
</tr>
<tr>
<td></td>
<td>LibLinear</td>
<td>99.03</td>
<td>0.00</td>
<td>50.50</td>
</tr>
</tbody>
</table>

To illustrate the potential gain of SDCA, we compare it with a DCA for solving the problem (21). From the Table 2, we see that the gain of running time of SDCA ranges from 11.6 times ($aloi$) to 55.1 times ($coverytype$).

Concerning the classification accuracy, SDCA and DCA are comparable. SDCA gives slightly better accuracy than DCA on $coverytype$, $sim_1$, $sim_3$, with a gain ranges from 0.06% to 0.28%. The gain of SDCA is higher on 2 datasets ($madelon$, $sim_2$), 1.35% and 3.07%. DCA furnishes a better result on $aloi$ and $sensorless$, especially the gain is up to 4.9% on $sensorless$. The results prove that SDCA can greatly improve the running time of DCA while archiving a similar out our experiments.
accuracy.

3.1.5. EXPERIMENT 2: SIMULATED DATASET

For synthetic datasets (sim_1, sim_2 and sim_3), we know in advance the informative features that were used to generate the datasets. Hence, the purpose of this experiment is to study the ability of algorithms to select these informative features in order to furnish a good classification accuracy. The comparison is performed with 3 algorithms, SDCA, msgl and LibLinear. We report the results in Table 2, and observe that.

For sim_1 dataset, LibLinear gives a slightly better classification accuracy (72.62%) comparing to SDCA (72.24%) and msgl (72.33%). However, SDCA is by far the fastest algorithm. SDCA is 1532 (resp. 136) times faster than LibLinear (resp. msgl). Furthermore, SDCA and LibLinear successfully suppress the 20% uninformative features, which it also matches with our procedure of generating this synthetic dataset. msgl fails on this purpose by selecting 82% of features.

For sim_2 dataset, SDCA is the best algorithm on both criteria: classification accuracy and running time. Similarly to sim_1 dataset, only SDCA and LibLinear can correctly select the informative features (80%).

For sim_3 dataset, SDCA exceeds LibLinear and GLASSO on all three comparison criteria: classification accuracy, sparsity and speed. LibLinear almost selects all the features (97.16% selected) but gives 0.89% accuracy lower than SDCA (99.93%), and it is also 2 times slower than SDCA. Among the three algorithm, only SDCA successfully selects the informative features.

To summarize, for all three synthetic datasets, SDCA successfully selects the exact informative features. LibLinear selects the exact features on 2 out of 3 datasets while GLASSO fails on all three datasets.

3.1.6. EXPERIMENT 3: REAL-WORLD DATASETS

In this experiment, we perform the comparative study between SDCA, msgl and LibLinear on real-world datasets. We observe from Table 2 that.

For aloi dataset, SDCA only selects 63.87% of features for a classification accuracy of 82.98% while LibLinear has a worse accuracy with 100% of features used. Moreover, SDCA is the 12.6 times faster than LibLinear while msgl fails to furnish a result after 2 hours of running time.

For coverytype dataset, SDCA furnishes better classification accuracy than LibLinear and msgl. Moreover, SDCA is by far faster than the two others. SDCA is SDCA is 45 times faster than LibLinear and 89 times faster than msgl. Concerning the sparsity of solution, msgl is the best while LibLinear fails to suppress features.

The dataset madelon is known to be non-linear. Hence, all three algorithms furnish quite low classification accuracy (62.38% for SDCA, 61.54% for LibLinear and 60.48% for msgl). As for the sparsity, SDCA suppresses more features than LibLinear and msgl.

For sensorless dataset, SDCA is better than both LibLinear and msgl on all three aspects: classification accuracy, sparsity and running time. In terms of classification accuracy, the gain of SDCA versus msgl (resp. LibLinear) is 3.89% (resp. 2.26%). Regarding the running time, SDCA is 113 times faster than msgl and 137 times faster than LibLinear. As for the sparsity, SDCA selects 10% less features than msgl while LibLinear fails to suppress features.

Overall, SDCA gives the best among the three in term of classification accuracy on all 4 datasets. As for running time, SDCA is by far the fastest algorithm. Concerning the sparsity of solution, SDCA suppresses more features than the two others on 3 out of 4 datasets.

4. Conclusions

We have rigorously studied the large-sum optimization problem involving $\ell_{2,0}$ regularization. The $\ell_{2,0}$-norm is approximated by a DC function, namely the piecewise exponential function. The resulting problem is then reformulated as a DC program and we developed stochastic DCA to solve it. Exploiting the fact that each component $f_i(x)$ is differentiable with $L$-Lipschitz gradient, we propose, a stochastic version of DCA that is very inexpensive. At each iteration, the algorithm only requires the computing the subgradients of a small subset of functions and the projection of points onto balls that is explicitly computed. We have also proved that the convergence is guaranteed with probability one. As an application, we applied our algorithm to the group variables selection in multiclass logistic regression problem. Numerical experiments were carefully conducted on both synthetic and real-world datasets. The numerical results show that SDCA greatly improves the running time of DCA while giving similar accuracy. Moreover, our algorithm SDCA outperforms standard algorithms (Liblinear and msgl) on all 3 criteria: classification accuracy, sparsity of solution and running time. Especially, the gain in running time is huge. SDCA is up to 210 times faster than msgl and 1537 times faster than LibLinear. We are convinced that stochastic DCA is a promising approach for handling very large-scale datasets in machine learning.
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