

The Risk Profile Problem for Stock Portfolio Optimization

Extended Abstract

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ABSTRACT

In this paper we study the problem of determining an optimal investment strategy for investors with different attitudes toward the trade-offs of risk and profit. The probability distribution of the return values of the stocks that are considered by the investor are assumed to be known, while the joint distribution is unknown. The problem is to find the best investment strategy in order to minimize the probability of losing a certain percentage of the invested capital based on different attitudes of the investors toward future outcomes of the stock market.

We show that for portfolios made up of two stocks, we can exactly and quickly solve the problem of finding an optimal portfolio for aggressive or risk-averse investors, using an algorithm based on a fast greedy solution to a maximum flow problem. However, an investor looking for an average-case guarantee (so is neither aggressive or risk-averse) must deal with a more difficult problem. In particular, we show that computing the distribution function associated with the average-case bound is $\#P$ -complete. On the positive side, we show how to use random sampling techniques similar to those for high-dimensional volume estimation to provide approximate answers. When $k > 2$ stocks are considered, we show that a simple solution based on the same flow concepts as our 2-stock algorithm would imply that $P = NP$, so is highly unlikely. We give approximation algorithms for this case as well as exact algorithms for some important special cases.

1. INTRODUCTION

This paper initiates the study of the *risk profile* problem for stock portfolio optimization. The problem has several variants depending on a given investor's preference toward the trade-off between risk and return [7].

^{*}Research supported in part by NSF CCR-Grant 9531028.

[†]Supported in part by Texas Advanced Research Program Grant 1997-003594-019.

In the problem, the investor has an initial amount of capital, which is normalized to one dollar. She considers k different stocks S_1, \dots, S_k and wishes to invest some x_i dollars in each stock S_i for a certain period of time, where $\sum_{i=1}^k x_i = 1$ and $x_i \geq 0$ for all i . The vector $\vec{x} = \langle x_i \rangle_{i=1}^k = \langle x_1, x_2, \dots, x_k \rangle$ is called a *portfolio*. Let \mathcal{P}_k be the set of all portfolios for k stocks. The *return* of \vec{x} is the ratio, expressed as a percentage, of the worth of this portfolio at the end of the investment period to the initial investment of one dollar. The *return* of stock S_j is the ratio of its price at the end of the investment period to its initial price, which is the same as the return of the portfolio $\langle x_i \rangle_{i=1}^k$ with $x_j = 1$ and all the other $x_i = 0$.

Let μ be a positive real number. Let m_1 and m_2 be integers with $m_1 < m_2$, and let $m = m_2 - m_1 + 1$. Let $\Delta = \{\ell\mu \mid \ell = m_1, \dots, m_2\}$. Each stock S_i is associated with a discrete probability distribution \mathcal{S}_i over Δ , where $\mathcal{S}_i(\beta)$ is the probability that the stock's return is $\beta\%$. For the sake of technical convenience, we allow m_1 and m_2 to be negative. The probability distributions $\mathcal{S}_1, \dots, \mathcal{S}_k$ are part of the input in our problem and are obtainable, e.g., by observing historical market data. We assume that non-zero values satisfy $\mathcal{S}_1(\beta) \geq 1/m^c$ for some constant c , and when representation is important we assume that these values can be represented as fixed-point numbers with $O(\log m)$ bits. The parameters μ , m_1 , and m_2 control the precision and range of such observations. For instance, for $\mu = 1$, $m_1 = 0$, and $m_2 = 200$, the set of possible returns are $0\%, 1\%, \dots, 200\%$. The joint distribution of the k probability distributions \mathcal{S}_i is usually unavailable for a variety of practical reasons. In particular, a joint distribution consists of m^k entries and thus would require observing an exponential number of data points in k .

The investor's goal is to find a portfolio \vec{x} , which is optimal according to her risk preference in six basic cases as follows. For a *risk-averse* investor, minimizing loss is more important than maximizing win, while an *aggressive* investor has the opposite priority. Each of these two investor types can be further classified into three subtypes, namely, *best-case*, *worst-case*, and *average-case*, referring to whether the probability of loss or win is estimated in the best, worst, or average case over the feasible joint distributions. More precisely, for each of these six types, the investor first chooses a *target* return α and then looks for such a portfolio \vec{x} that optimizes one of the following six probabilities:

- $\mathcal{RA}_b(\alpha, \vec{x})$ (respectively, $\mathcal{RA}_w(\alpha, \vec{x})$ or $\mathcal{RA}_a(\alpha, \vec{x})$) denotes the smallest (respectively, largest or average) probability that the return of \vec{x} is at most $\alpha\%$ over

all joint distributions for $\mathcal{S}_1, \dots, \mathcal{S}_k$.

- $\mathcal{AG}_b(\alpha, \vec{x})$ (respectively, $\mathcal{AG}_w(\alpha, \vec{x})$ or $\mathcal{AG}_a(\alpha, \vec{x})$) denotes the largest (respectively, smallest or average) probability that the return of \vec{x} is at least $\alpha\%$ over all joint distributions for $\mathcal{S}_1, \dots, \mathcal{S}_k$.

A best-case (respectively, worst-case or average-case) risk-averse investor would choose \vec{x} to minimize $\mathcal{RA}_b(\alpha, \vec{x})$ (respectively, $\mathcal{RA}_w(\alpha, \vec{x})$ or $\mathcal{RA}_a(\alpha, \vec{x})$). In contrast, if the investor is best-case (respectively, worst-case or average-case) aggressive, she would choose \vec{x} to maximize $\mathcal{AG}_b(\alpha, \vec{x})$ (respectively, $\mathcal{AG}_w(\alpha, \vec{x})$ or $\mathcal{AG}_a(\alpha, \vec{x})$).

While the risk profile problem originates from a very applied field, the corresponding mathematical model has a substantial combinatorial structure. In the cases where the investor is highly risk-averse or highly aggressive, we can model the problem as a network flow problem. Quite surprisingly, in the two-stock case, this flow problem is solvable by a simple greedy algorithm in $O(m)$ time. In contrast, for the three-stock case, the applicability of a greedy flow-based algorithm would imply $P = NP$. If the number k of stocks is part of the input, we give an exact algorithm based on linear programming which takes time polynomial in the number of entries of a corresponding contingency table but exponential in the input size. To supplement this algorithm, we also give a polynomial-time approximation algorithm based on linear programming. We further present an exact polynomial-time algorithm in the practical case where the capital can only be broken up into a fixed number of units (e.g., cents).

It remains open whether this problem is NP -hard if the number of stocks is part of the input. We strongly suspect that this is indeed the case.

In the case of an average-case investor we show $\#P$ -hardness of the problem of computing the distribution function over various probability bounds, a natural first-step in solving the average-case investor problem. This hardness result holds even in two dimensions, and we describe an approximation algorithm for this case. This algorithm uses a random walk approach to sample from the feasible joint distributions, and is closely related to volume computation and sampling from log-concave distributions.

Section 2 defines some notation. Section 3 discusses the case where there are only two stocks under consideration, and Section 4 discusses the case of general k .

2. NOTATION

Let $\vec{\delta} \in \Delta^k$ denote a vector $\langle \delta_1, \dots, \delta_k \rangle$, where $\delta_i \in \Delta$. Let

$$M = [M_{\vec{\delta}}]_{\vec{\delta} \in \Delta^k}$$

denote a k -dimensional matrix indexed by Δ^k . Let \mathcal{M}_k denote the set of k -dimensional matrices for all possible joint distributions of $\mathcal{S}_1, \dots, \mathcal{S}_k$; i.e., \mathcal{M}_k consists of all matrices

$$M = [M_{\vec{\delta}}]_{\vec{\delta} \in \Delta^k},$$

where (1) $M_{\vec{\delta}}$ is the probability that the return of stock \mathcal{S}_i is $\delta_i\%$ for $i = 1, \dots, k$, and (2) thus for all $\vec{\delta} \in \Delta^k$, $M_{\vec{\delta}} \geq 0$ and for all $\beta \in \Delta$ and $j = 1, \dots, k$,

$$\mathcal{S}_j(\beta) = \sum_{\vec{\delta} \in \Delta^k; \delta_j = \beta} M_{\vec{\delta}}.$$

For instance, \mathcal{M}_k contains the matrix M defined by

$$M_{\vec{\delta}} = \prod_{i=1}^k \mathcal{S}_i(\delta_i).$$

Also, in the two-stock case, each $M \in \mathcal{M}_2$ is just a two-dimensional $m \times m$ matrix, where for all $\delta_1, \delta_2 \in \Delta$, the entries of M in column δ_1 sum up to $\mathcal{S}_1(\delta_1)$ and those in row δ_2 sum up to $\mathcal{S}_2(\delta_2)$.

Given a portfolio $\vec{x} \in \mathcal{P}_k$ and a target return α , let

$$\begin{aligned} L(\alpha, \vec{x}) &= \left\{ \vec{\delta} \in \Delta^k \mid \sum_{i=1}^k x_i \delta_i \leq \alpha \right\}, \\ L^{**}(\alpha, \vec{x}) &= \left\{ \vec{\delta} \in \Delta^k \mid \sum_{i=1}^k x_i \delta_i < \alpha \right\}, \\ U(\alpha, \vec{x}) &= \left\{ \vec{\delta} \in \Delta^k \mid \sum_{i=1}^k x_i \delta_i \geq \alpha \right\}, \\ U^{**}(\alpha, \vec{x}) &= \left\{ \vec{\delta} \in \Delta^k \mid \sum_{i=1}^k x_i \delta_i > \alpha \right\}, \end{aligned}$$

which are the sets of the indices of all entries in the matrices in \mathcal{M}_k such that the return of \vec{x} is at most, less than, at least, and more than $\alpha\%$, respectively. We further define the following functions on $M \in \mathcal{M}_k$:

$$\begin{aligned} \bar{L}_{\alpha, \vec{x}}(M) &= \sum_{\vec{\delta} \in L(\alpha, \vec{x})} M_{\vec{\delta}}, \\ \bar{L}_{\alpha, \vec{x}}^{**}(M) &= \sum_{\vec{\delta} \in L^{**}(\alpha, \vec{x})} M_{\vec{\delta}}, \\ \bar{U}_{\alpha, \vec{x}}(M) &= \sum_{\vec{\delta} \in U(\alpha, \vec{x})} M_{\vec{\delta}}, \\ \bar{U}_{\alpha, \vec{x}}^{**}(M) &= \sum_{\vec{\delta} \in U^{**}(\alpha, \vec{x})} M_{\vec{\delta}}, \end{aligned}$$

which are the probabilities in the joint distribution M that the return of \vec{x} is at most, less than, at least, and more than $\alpha\%$, respectively. Formally, if $u_{\mathcal{M}_k}(M)$ is a uniform density over \mathcal{M}_k ,

$$\mathcal{RA}_b(\alpha, \vec{x}) = \min_{M \in \mathcal{M}_k} \bar{L}_{\alpha, \vec{x}}(M); \quad (1)$$

$$\mathcal{RA}_w(\alpha, \vec{x}) = \max_{M \in \mathcal{M}_k} \bar{L}_{\alpha, \vec{x}}(M); \quad (2)$$

$$\mathcal{RA}_a(\alpha, \vec{x}) = \int_{\mathcal{M}_k} \bar{L}_{\alpha, \vec{x}}(M) u_{\mathcal{M}_k}(M) dM; \quad (3)$$

$$\mathcal{AG}_b(\alpha, \vec{x}) = \max_{M \in \mathcal{M}_k} \bar{U}_{\alpha, \vec{x}}(M); \quad (4)$$

$$\mathcal{AG}_w(\alpha, \vec{x}) = \min_{M \in \mathcal{M}_k} \bar{U}_{\alpha, \vec{x}}(M); \quad (5)$$

$$\mathcal{AG}_a(\alpha, \vec{x}) = \int_{\mathcal{M}_k} \bar{U}_{\alpha, \vec{x}}(M) u_{\mathcal{M}_k}(M) dM. \quad (6)$$

For example, in the two-stock case, $L(\alpha, \langle x_1, x_2 \rangle)$ is the set of all indices in a two-dimensional table M in \mathcal{M}_2 on or below the line $x_1 \delta_1 + x_2 \delta_2 = \alpha$, and $\mathcal{RA}_w(\alpha, \langle x_1, x_2 \rangle)$ maximizes the sum of the entries in this region under the condition that M has the given column and row sums of $\mathcal{S}_1(m_1), \dots, \mathcal{S}_1(m_2), \mathcal{S}_2(m_1), \dots, \mathcal{S}_2(m_2)$.

For technical convenience, we also define the following terms:

$$\mathcal{RA}_b^{**}(\alpha, \vec{x}) = \min_{M \in \mathcal{M}_k} \bar{L}_{\alpha, \vec{x}}^{**}(M); \quad (7)$$

$$\mathcal{RA}_w^{**}(\alpha, \vec{x}) = \max_{M \in \mathcal{M}_k} \bar{L}_{\alpha, \vec{x}}^{**}(M); \quad (8)$$

$$\mathcal{RA}_a^{**}(\alpha, \vec{x}) = \int_{\mathcal{M}_k} \bar{L}_{\alpha, \vec{x}}^{**}(M) dM; \quad (9)$$

$$\mathcal{AG}_b^{**}(\alpha, \vec{x}) = \max_{M \in \mathcal{M}_k} \bar{U}_{\alpha, \vec{x}}^{**}(M); \quad (10)$$

$$\mathcal{AG}_w^{**}(\alpha, \vec{x}) = \min_{M \in \mathcal{M}_k} \bar{U}_{\alpha, \vec{x}}^{**}(M); \quad (11)$$

$$\mathcal{AG}_a^{**}(\alpha, \vec{x}) = \int_{\mathcal{M}_k} \bar{U}_{\alpha, \vec{x}}^{**}(M) dM. \quad (12)$$

LEMMA 1. *The following statements hold.*

$$\min_{\vec{x} \in \mathcal{P}_k} \mathcal{RA}_b(\alpha, \vec{x}) = 1 - \max_{\vec{x} \in \mathcal{P}_k} \mathcal{AG}_b^{**}(\alpha, \vec{x}) \quad (13)$$

$$\min_{\vec{x} \in \mathcal{P}_k} \mathcal{RA}_w(\alpha, \vec{x}) = 1 - \max_{\vec{x} \in \mathcal{P}_k} \mathcal{AG}_w^{**}(\alpha, \vec{x}) \quad (14)$$

$$\min_{\vec{x} \in \mathcal{P}_k} \mathcal{RA}_a(\alpha, \vec{x}) = 1 - \max_{\vec{x} \in \mathcal{P}_k} \mathcal{AG}_a^{**}(\alpha, \vec{x}) \quad (15)$$

$$\max_{\vec{x} \in \mathcal{P}_k} \mathcal{AG}_b(\alpha, \vec{x}) = 1 - \min_{\vec{x} \in \mathcal{P}_k} \mathcal{RA}_b^{**}(\alpha, \vec{x}) \quad (16)$$

$$\max_{\vec{x} \in \mathcal{P}_k} \mathcal{AG}_w(\alpha, \vec{x}) = 1 - \min_{\vec{x} \in \mathcal{P}_k} \mathcal{RA}_w^{**}(\alpha, \vec{x}) \quad (17)$$

$$\max_{\vec{x} \in \mathcal{P}_k} \mathcal{AG}_a(\alpha, \vec{x}) = 1 - \min_{\vec{x} \in \mathcal{P}_k} \mathcal{RA}_a^{**}(\alpha, \vec{x}) \quad (18)$$

PROOF. Straightforward. \square

In light of Lemma 1, to solve the risk profile problem, it suffices to show how to compute

$$\begin{aligned} \min_{\vec{x} \in \mathcal{P}_k} \mathcal{RA}_b(\alpha, \vec{x}), & \quad \min_{\vec{x} \in \mathcal{P}_k} \mathcal{RA}_b^{**}(\alpha, \vec{x}), \\ \min_{\vec{x} \in \mathcal{P}_k} \mathcal{RA}_w(\alpha, \vec{x}), & \quad \min_{\vec{x} \in \mathcal{P}_k} \mathcal{RA}_w^{**}(\alpha, \vec{x}), \\ \min_{\vec{x} \in \mathcal{P}_k} \mathcal{RA}_a(\alpha, \vec{x}), & \quad \min_{\vec{x} \in \mathcal{P}_k} \mathcal{RA}_a^{**}(\alpha, \vec{x}). \end{aligned}$$

The techniques for computing the three expressions on the right are essentially the same as those for computing the three on the left. Furthermore, the techniques for computing the first expression on the left are almost identical to those for computing the second. For these reasons, the remainder of this paper focuses on how to compute $\min_{\vec{x} \in \mathcal{P}_k} \mathcal{RA}_w(\alpha, \vec{x})$ and $\min_{\vec{x} \in \mathcal{P}_k} \mathcal{RA}_a(\alpha, \vec{x})$.

3. THE TWO-STOCK CASE

This section assumes that $k = 2$, i.e., there are only two stocks under consideration. In the case of two stocks, we can visualize the problems under consideration as in Figure 1. The discrete and finite set of possible return pairs for the two stocks in the portfolio are shown as the dots in this picture – each pair has a probability (from the joint distribution) associated with it, with the given restrictions on column and row sums. A given portfolio and target return α defines a half-space on the set of return pairs, with the shaded area in Figure 1 giving the area in which the total return is $\leq \alpha$. The problem of computing $\mathcal{RA}_w(\alpha, \vec{x})$ then is the problem of determining which feasible assignment of joint probabilities places the highest total probability in the shaded region.

3.1 A Worst-Case or Best-Case Investor

Given a worst-case investor, we can show how to compute an optimal portfolio for a worst-case risk-averse investor. The cases of a best-case risk-averse investor, a worst-case aggressive investor, and a best-case aggressive investor can be solved similarly.

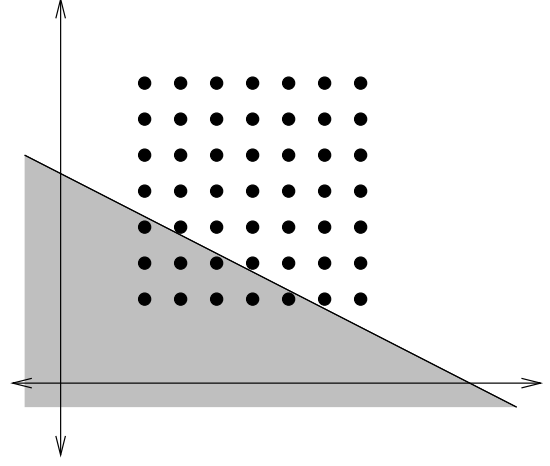


Figure 1: Visualization of two stock case

We first present a basic algorithm to compute $\mathcal{RA}_w(\alpha, \vec{x})$ by computing a worst-case joint distribution matrix M for S_1 and S_2 . For convenience, we index the entries of M with $\{(i, j) \mid i, j = m_1, \dots, m_2\}$, where row i (respectively, column i) corresponds to return $i\mu$ of S_1 (respectively, $j\mu$ of S_2). We model the problem of computing M as a network flow problem on the graph G defined below:

- G has $2(m+1)$ vertices, namely, a source s , a sink t , and $v_{m_1}, \dots, v_{m_2}, w_{m_1}, \dots, w_{m_2}$, where v_i (respectively, w_i) corresponds to return $i\mu$ of stock S_1 (respectively, stock S_2).
- For all $i, j = m_1, \dots, m_2$, G has (1) edge (v_i, w_j) , which has capacity $c(v_i, w_j) = 1$ if $x_1 i\mu + x_2 j\mu \leq \alpha$ or 0 otherwise; (2) the edge (s, v_i) with capacity $c(s, v_i) = S_1(i\mu)$; and (3) the edge (w_j, t) with capacity $c(w_j, t) = S_2(j\mu)$.

Geometrically, we wish to push as much probability as possible into the region of M defined by $x_1 i + x_2 j \leq \frac{\alpha}{\mu}$. In other words, the value of a maximum $s-t$ flow of G equals $\mathcal{RA}_w(\alpha, \vec{x})$. Thus, it is tempting to use a maximum flow algorithm to solve this maximum flow problem. The fastest known algorithm for this problem is due to Goldberg and Rao [4] and runs in $O^*(m^{2\frac{2}{3}})$ time¹ for our application (note that m in this bound is as defined in this paper, not as the number of edges which is typical in general flow discussion). Instead of using this algorithm, we exploit some structural properties of G to solve the flow problem using a simple greedy algorithm in $O(m)$ arithmetic operations. Note that since G may have $\Omega(m^2)$ edges with positive capacity, we cannot afford to construct the whole G explicitly. The idea of our $O(m)$ -time algorithm can be described as follows.

Starting with v_{m_2} , we try to push a flow of $c(s, v_{m_2})$ through G . Assume $c(v_{m_2}, w_{m_1}) = 1$ for simplicity. We consider the path formed by edges $(s, v_{m_2}), (v_{m_2}, w_{m_1}), (w_{m_1}, t)$ first. We can push flow $\min(c(s, v_{m_2}), c(w_{m_1}, t))$ through

¹We use $O^*(f(n))$ for the “soft- O ” notation, which ignores polylogarithmic factors. In bounds for the approximation algorithms, this notation also ignores factors that depend only on the approximation bound ϵ .

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 $F \leftarrow 0$ 
 $i \leftarrow m_2$ 
 $cv \leftarrow c(s, v_i)$ 
 $j \leftarrow m_1$ 
 $cw \leftarrow c(w_j, t)$ 
loop
  if  $c(v_i, w_j) = 1$  and  $cw \leq cv$  then
     $F \leftarrow F + cw$ 
     $cv \leftarrow cv - cw$ 
     $j \leftarrow j + 1$ 
    if  $j > m_2$  then return  $F$ 
     $cw \leftarrow c(w_j, t)$ 
  else
    if  $c(v_i, w_j) = 1$  then
       $F \leftarrow F + cv$ 
       $cw \leftarrow cw - cv$ 
    end if
     $i \leftarrow i - 1$ 
    if  $i < m_1$  then return  $F$ 
     $cv \leftarrow c(s, v_i)$ 
  end if
end loop

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Figure 2: The procedure Greedy-Flow

this path, saturating either (s, v_{m_2}) or (w_{m_1}, t) . If we saturated (s, v_{m_2}) then we next consider the path $(s, v_{m_2-1}), (v_{m_2-1}, w_{m_1}), (w_{m_1}, t)$ for pushing additional flow; however, if we had saturated (w_{m_1}, t) we will next consider the path $(s, v_{m_2}), (v_{m_2}, w_{m_1+1}), (w_{m_1+1}, t)$. We continue in this fashion until we can push no more flow. The only complication is that if at some point we are considering the path $(s, v_i), (v_i, w_j), (w_j, t)$, and $c(v_i, w_j) = 0$, then obviously we can't saturate *either* (s, v_i) or (v_j, t) , and we simply decrease i to next consider the path $(s, v_{i-1}), (v_{i-1}, w_j), (w_j, t)$. The details of this $O(m)$ time algorithm are given in Figure 2.

THEOREM 2. *Given $\mathcal{S}_1, \mathcal{S}_2$, a valid portfolio vector \vec{x} , and α as input, Greedy-Flow computes the value of a maximum flow of G in $O(m)$ arithmetic operations.*

PROOF. As a first step we prove that the algorithm computes the maximal flow. Let ℓ be the minimal index such that (w_ℓ, t) is not saturated after termination of the algorithm and k be the minimal index such that $c(v_k, w_\ell) = 0$. We define a partition $V_1 \cup V_2$ of the nodes by

$$V_1 = \{s, v_k, \dots, v_{m_2}, w_{m_1}, \dots, w_{\ell-1}\}, \quad V_2 = \bar{V}_1.$$

It is trivial from the definition of j that the edges $e = (w_i, t), i = \{m_1, \dots, \ell-1\}$ are saturated.

Since $x_1, x_2 \geq 0$, and k is the minimal value such that $c(v_k, w_\ell) = 0$, we have $c(v_i, w_\ell) = 1$ for $i = m_1, \dots, k-1$. Since (w_ℓ, t) is not saturated, all edges (s, v_i) , for $i = m_1, \dots, k-1$ must be saturated.

From the definition of k and the non-negativity of the portfolio vector it is easy to see that edges $e = (v_i, w_j)$ for $i \in \{k, \dots, m_2\}$, $j \in \{\ell, \dots, m_2\}$ and positive capacity cannot exist. Thus, every edge $e = (x, y)$ with $x \in V_1$ and $y \in V_2$ is saturated. The Max-Flow-Min-Cut Theorem then implies that the algorithm indeed computes a maximal flow.

Observing the fact that in each loop iteration either index i is decremented or index j is incremented, and that there are only m different values that either i or j can

take on before the algorithm terminates, there are at most $2m-1$ loop iterations, and the linear running time bound follows. \square

To compute $\inf\{\mathcal{RA}_w(\alpha, \vec{x}) \mid \sum x_i = 1\}$ we have to compute $\mathcal{RA}_w(\alpha, \vec{x})$ for all possible portfolios $\langle x_1, x_2 \rangle$. However, each feasible portfolio corresponds to a half-space (as in Figure 1) defined by a line that goes through the point (α, α) ($x_1\alpha + x_2\alpha = \alpha$, since $x_1 + x_2 = 1$), so we only need to consider the $O(m^2)$ distinct subsets of return pairs that can be defined by a line going through (α, α) . We can identify each such portfolio with a different (non-positive) slope s_1, \dots, s_{m^2} , which we assume to be sorted in descending order. By using a suitable data structure it is possible to compute the best portfolio much faster than the obvious $O(m^3)$ -time algorithm that starts the greedy algorithm for each slope.

THEOREM 3. *Given $\mathcal{S}_1, \mathcal{S}_2$, and α , we can compute in $O(m^2 \log m)$ arithmetic operations a portfolio $\langle x_1, x_2 \rangle$ for a worst-case risk-averse investor which minimizes equation (2).*

PROOF. Starting with the first slope s_1 we build up a binary tree. Each is labeled with a pair of two real entries (e_1, e_2) . The leaves of the tree correspond to the rows and the columns in the following way.

Starting from column m_2 we add leaves from left to right. We add leaves with labels $(0, \mathcal{S}_2(m_1\mu)), (0, \mathcal{S}_2((m_1+1)\mu)), \dots, (0, \mathcal{S}_2(j_m\mu))$, until we reach a row index j_m such that $x_1 m_2 \mu + x_2 (j_m+1)\mu > \alpha$, i.e., this index is the last under the crucial line. To be precise we let $j_m = \lfloor \frac{\alpha - x_1 m_2 \mu}{x_2 \mu} \rfloor$; note that it may be the case that $j_m < m_1$, so this sequence of leaves may be empty. Then we add the leaf $(-\mathcal{S}_1(m_2\mu), 0)$. Next, we consider column m_2-1 and add leaves $(0, \mathcal{S}_2((j_m+1)\mu)), \dots, (0, \mathcal{S}_2((j_{m-1})\mu))$, until we reach an index j_{m-1} , such that $x_1 (m_2-1)\mu + x_2 (j_{m-1}+1)\mu > \alpha$. Then we add the leaf $(-\mathcal{S}_1((m_2-1)\mu), 0)$ and proceed similarly with column m_2-2 . Note that the order of adding leaves is crucial to this data structure and the correctness of the algorithm is based on that. Starting from left to right we group the leaves in pairs of 2 and build a parent node for each pair according to the following rule

$$\text{parent}[(e_1, e_2), (f_1, f_2)] = (e_1 + \min\{e_2 + f_1, 0\}, \max\{e_2 + f_1, 0\} + f_2).$$

We build $O(\log m)$ layers iteratively, until we reach a single root node (r_1, r_2) . It is easy to see that this tree based algorithm imitates the greedy algorithm described before and that $1 + r_1 = 1 - r_2$ is exactly the flow value. Building this tree structure takes constant time per tree node, and since there are $O(m)$ nodes we have a total time of $O(m)$, which is no better than the time bound of the greedy algorithm. The advantage is that we can dynamically update this data structure efficiently.

We will first sort all of the m^2 possible return pairs by their slope with the point (α, α) , so that as the slope determined by our portfolio increases we can quickly (in constant time per pair) determine which pairs are added and which are removed from our halfspace of interest. This takes $O(m^2 \log m)$ time. To update our data structure for each point insertion/removal, all that is required is swapping the position of two neighboring leaves. With obvious techniques, the positions of these two leaves can be found in $O(1)$ time,

and we can update the tree by looking at the path from the two leaves to the root and update each node on that path. Each update step requires $O(1)$ operations and the length of the path is bounded by $O(\log m)$. Since there are at most m^2 point additions and removals, each taking $O(\log m)$ time, it takes at most $O(m^2 \log m)$ time to consider all possible portfolios. \square

3.2 The Average-Case Investor

For $\alpha \in \mathcal{AG}_a$, we are not interested in the extremes of the joint distributions, but rather the distribution of the feasible tables. In this section we consider $Q = \bar{L}_{\alpha, \vec{x}}(M)$ a random variable where M is drawn from a uniform distribution over the feasible tables \mathcal{M}_k . The definition of $\mathcal{RA}_a(\alpha, \vec{x})$, from (3), is then $E[Q]$. We will see that computing the distribution function of Q is a computationally difficult problem to solve exactly, but can be approximated within a reasonable (polynomial) amount of time.

THEOREM 4. *Let $\gamma \in [0, 1]$ be an n -bit rational. It is $\sharp P$ -hard to compute the fraction of feasible tables $M \in \mathcal{M}_2$ with*

$$\bar{L}_{\alpha, \vec{x}}(M) = \sum_{\delta \in L(\alpha, \vec{x})} M_\delta \leq \gamma$$

(the integration of the corresponding indicator function, or the distribution function for Q).

PROOF. Given positive integers a_1, \dots, a_n, b , it is shown in [1] that computing the n -dimensional volume of the polyhedron P

$$\sum_{j=1}^n a_j y_j \leq b \quad 0 \leq y_j \leq 1 \quad (j = 1, \dots, n)$$

is $\sharp P$ -hard. Let $d = \sum_{j=1}^n a_j$ and consider the polyhedron

$$\sum_{j=1}^{n+1} a_j y_j = d \quad 0 \leq y_j \leq 1 \quad (j = 1, \dots, n+1), \quad (19)$$

where $a_{n+1} = d$. Note that for any valid assignment of values to y_1, y_2, \dots, y_n we have $0 \leq \sum_{j=1}^n a_j y_j \leq d$, so there is a $y_{n+1} \in [0, 1]$ that will satisfy (19). Now let $a'_i = a_i/(2d)$ and define a $2 \times (n+1)$ contingency table by $t_{1j} = a'_j y_j, t_{2j} = a'_j (1 - y_j)$, with row sums $(1/2, 1/2)$ and column sums (a'_1, \dots, a'_{n+1}) .

To completely define our stock problem, we must also give values for μ, α , the portfolio $\vec{x} = \langle x_1, x_2 \rangle$, and the threshold γ , which we do as follows:

$$\begin{aligned} \mu &= 1, & x_1 &= \frac{1}{n+1}, & x_2 &= \frac{n}{n+1}, \\ \alpha &= \frac{2n}{n+1}, & \gamma &= \frac{b}{2d}. \end{aligned}$$

It is straight-forward to verify from these values that the return pairs in the critical region (the shaded region in Figure 1) are exactly the entries t_{1j} for $j = 1, \dots, n$. Therefore, the tables that satisfy our criteria, that $\bar{L}_{\alpha, \vec{x}}(M) \leq \gamma$, are precisely those with

$$\begin{aligned} \sum_{j=1}^n t_{1j} \leq \gamma &\iff \sum_{j=1}^n a'_j y_j \leq \gamma \\ &\iff \sum_{j=1}^n a_j y_j \leq \gamma \cdot 2d = b. \end{aligned}$$

Therefore the feasible tables that meet our criteria are exactly those that correspond to points in polyhedron P , and so the fraction of tables that meet the criteria is exactly the volume of P . \square

Following the notation of Dyer, Kannan and Mount [3], who describe a sampling procedure for contingency tables with integer entries and large row and column sums ($\geq \Omega(m^3)$), we define

$$V(r, c) = \left\{ x \in \mathbf{R}^{m \times m} \mid \begin{aligned} &\sum_j x_{ij} = r_i \text{ for } i = 1, \dots, m \\ &\text{and } \sum_i x_{ij} = c_j \text{ for } j = 1, \dots, m \end{aligned} \right\}$$

and

$$P(r, c) = V(r, c) \cap \{x \mid x_{ij} \geq 0 \text{ for } i = 1, \dots, m, j = 1, \dots, m\}$$

as the contingency polytope. Thus, $V(r, c)$ is the set of matrices with row and column sums specified by r and c respectively. In our case $r_i = \mathcal{S}_1(i\mu)$, $c_i = \mathcal{S}_2(i\mu)$ and $P(r, c)$ is the set of joint distributions \mathcal{M}_k .

Let U be the lattice

$$\begin{aligned} \{x \in \mathbf{Z}^{m \times m} \mid &\sum_j x_{ij} = 0 \text{ for } i = 1, \dots, m, \\ &\sum_i x_{ij} = 0 \text{ for } j = 1, \dots, m\}. \end{aligned}$$

For $1 \leq i \leq m-1$ and $1 \leq j \leq m-1$, let $b(ij)$ be the vector in $\mathbf{R}^{m \times m}$ given by $b(ij)_{i,j} = 1, b(ij)_{i+1,j} = -1, b(ij)_{i,j+1} = -1, b(ij)_{i+1,j+1} = 1$ and $b(ij)_{k,\ell} = 0$ for all other indices k, ℓ . Any vector x in $V(0, 0)$ can be expressed as linear combination of the $b(ij)$'s as follows

$$x = \sum_{k=1}^{m-1} \sum_{\ell=1}^{m-1} \left(\sum_{i=1}^k \sum_{j=1}^\ell x_{ij} \right) b(k\ell).$$

It is easy to see that the $b(ij)$ are all linearly independent and the dimension of $V(r, c)$ and $P(r, c)$ for positive row and column sum vectors r and c is $(m-1)^2$ [3]. We will apply the sampling algorithm pioneered by Dyer, Frieze and Kannan [2] and later refined in a sequence of papers (see [5] for an overview) to sample uniformly at random in $P(r, c)$.

We sample in the space $V(r, c)$. As mentioned in the introduction, we know a starting point z_0 in $P(r, c)$ (multiplication of rows and column sums). It is easy to see that a ball of radius b^2 is inside $P(r, c)$, if every component of r and c is at least b . Since in our case r and c sum up to one, $P(r, c) \subset B(0, 1)$. The following theorem is a corollary of the analysis of the fastest sampling algorithm in convex bodies known so far by Kannan, Lovasz and Simonovits [6].

THEOREM 5. *We can generate a point in $P(r, s)$, which is almost uniform in the sense that its distribution is at most ϵ away from the uniform in total variation distance. The algorithm uses $O^*(\frac{m^6}{b^4 \epsilon^2 \delta})$ membership queries of $P(r, s)$ (each requires $O(m^2)$ arithmetic operations).*

THEOREM 6. *Procedure Estimate (in Figure 3) computes a number S in $O^*(\frac{m^8}{b^4 \epsilon^2 \delta})$ arithmetic operations, which approximates $\mathcal{RA}_a(\alpha, \vec{x})$ as*

$$\mathcal{RA}_a(\alpha, \vec{x}) - \epsilon \leq S \leq \mathcal{RA}_a(\alpha, \vec{x}) + \epsilon$$

procedure Estimate(x)

$S \leftarrow 0$
 $N = \frac{100}{\epsilon^2 \delta}$
for $\ell = 1, \dots, N$ **do**
 $\zeta_i \leftarrow$ result from sample procedure started at x
 $S \leftarrow S + \bar{L}_{\alpha, \vec{x}}(\zeta_i)$
end for
 $S \leftarrow S/N$
return S

Figure 3: The approximation algorithm

with probability $1 - \delta$.

PROOF. Let $S_k = \frac{1}{k} \sum_{i=1}^k \bar{L}_{\alpha, \vec{x}}(\zeta_i)$. Thus

$$E(S_k) = \int \bar{L}_{\alpha, \vec{x}}(M) w(M) dM,$$

where w is the density produced by the random walk. Since $0 \leq \bar{L}_{\alpha, \vec{x}}(M) \leq 1$ for all $M \in \mathcal{M}_2$, it is easy to see that $\sigma^2(S_1) \leq 1$ and so $\sigma^2(S_k) \leq \frac{1}{k}$. By Chebychev's inequality,

$$P(|S_k - E(S_k)| \geq \epsilon/2) \leq \frac{\sigma^2(S_k)}{(\epsilon/2)^2} \leq \frac{4}{\epsilon^2 k}.$$

Since the samples are not entirely uniform, we must consider the error introduced by the approximately uniform sampling distribution as well. Let $u_{\mathcal{M}_k}(M)$ denote a uniform density over the set \mathcal{M}_k , and then approximating a uniform distribution within bound $\epsilon/4$, Theorem 5 implies

$$\begin{aligned} & |E(S_k) - \mathcal{RA}_a(\alpha, \vec{x})| \\ &= \left| \int \bar{L}_{\alpha, \vec{x}}(M) w(M) dM - \int \bar{L}_{\alpha, \vec{x}}(M) u_{\mathcal{M}_k}(M) dM \right| \\ &\leq \left[\int_{w > u_{\mathcal{M}_k}} (w(M) - u_{\mathcal{M}_k}(M)) dM + \right. \\ &\quad \left. \int_{w \leq u_{\mathcal{M}_k}} (u_{\mathcal{M}_k}(M) - w(M)) dM \right] \\ &\leq \epsilon/2. \end{aligned}$$

Setting $k = \frac{4}{\epsilon^2 \delta}$ the theorem follows. \square

4. THE k -STOCK CASE

In this chapter we consider the general case of more than two stocks. Since the problem of estimating the probability distribution for the average-case investor is already #P-complete in the two stock case, we do not consider it any more and concentrate on a worst-case investor. We start with a complexity result for three stocks, which implies that a greedy or flow based portfolio is quite unlikely to exist.

THEOREM 7. *The existence of a greedy or flow based portfolio for the problem with 3 or more stocks implies $P = NP$.*

PROOF. (Sketch) We prove this result by reduction from NUMERICAL-3-DIM-MATCHING. Consider an instance of NUMERICAL-3-DIM-MATCHING, i.e., disjoint sets X_1, X_2, X_3 , each containing m elements, a size $s(a) \in \mathbf{Z}^+$ for each element $a \in X_1 \cup X_2 \cup X_3$ and bound $B \in \mathbf{Z}$. We would like to know if $X_1 \cup X_2 \cup X_3$ can be partitioned into m disjoint sets such that each of these sets contains exactly

one element from each of X_1, X_2 , and X_3 , and the sum of the elements is exactly B (we can change this requirement to $\leq B$ without difficulty). This problem is NP-complete in the strong sense, so we restrict the sizes to be bounded by a polynomial, $s(a) \leq n^c$ for some constant c .

We construct an instance of the problem of computing $\mathcal{RA}_w(\alpha, \langle 1/3, 1/3, 1/3 \rangle)$ by making a contingency table in which $S_k(i) = c_{k,i}/m$, where $c_{k,i}$ is the number of items in set X_k with value i . The existence of a greedy or flow based algorithm implies the existence of a solution in which all entries in the solution table are multiples of $1/m$, and such a solution exists with $\bar{L}_{\alpha, \vec{x}}(M) = 1$ if and only if there is a valid partition of $X_1 \cup X_2 \cup X_3$. If such a partition exists, we can find it by simply taking all of the triples “selected” (with multiplicity determined by the integer multiple of $1/m$), and use elements from X_1, X_2 , and X_3 as determined by the three coordinates of each selected point. \square

While this proof shows that it is unlikely that a fast and simple greedy or flow-based algorithm exists, as it does for 2 stocks, we can indeed solve the problem for a fixed number of stocks in polynomial time using a more time-consuming procedure based on linear programming. This is stated in a general setting in the following theorem.

THEOREM 8. *If the number of stocks k is part of the input, the problem of determining the best portfolio for a worst-case investor can be solved in time polynomial in the number of entries of the contingency table (but exponential in k).*

PROOF. The problem can be modeled as linear program with a number of variables, that corresponds to the number of entries of the contingency table, and km inequalities. \square

4.1 An Approximation Algorithm

In this section we describe an approximation algorithm, that solves the problem of determining the worst case probability for a given portfolio within a given error $\epsilon \in \mathbf{R}^+$ in polynomial time. Additionally, we describe an important, non-trivial special case, where the problem can be solved exactly in polynomial time.

THEOREM 9. *Suppose that a portfolio $\langle x_i \rangle_{i=1}^k$ and a target return α are given. The worst-case probability can be approximated (i.e., we compute a value W with $\mathcal{RA}_w(\alpha, \vec{x}) - \epsilon \leq W \leq \mathcal{RA}_w(\alpha, \vec{x}) + \epsilon$) in time polynomial in k and n . The number of steps is dominated by solving a linear program in $O(km^2/\epsilon^2)$ variables and $O(km/\epsilon)$ constraints.*

PROOF. We consider the first pair of stocks S_1 and S_2 as in the two dimensional case and define a new portfolio as $\tilde{x}_1 = \frac{x_1}{x_1 + x_2}$ and $\tilde{x}_2 = \frac{x_2}{x_1 + x_2}$. We divide the two dimensional plane in $\ell = \frac{1}{\epsilon} m \log k$ regions by ℓ parallel lines $\tilde{x}_1 x + \tilde{x}_2 y = \text{const}$ of constant distance. Thus, we divide the entries of the joint distribution matrix into ℓ different sets (see Figure 4).

Each entry in the matrix corresponds to a variable and the variables satisfy the row sum and column sum condition of the joint distribution. Next, we sum up the entries in the ℓ different sets and assign the sums to ℓ new variables. By combining these sum variables from two different pairs of stocks, we get a new table with new row and column sum conditions, resulting again in ℓ new sum variables.

Repeating combinations in this manner, we stop after $\log k$ iterations and the creation of $O(km^2 \log k / \epsilon^2)$ variables and $O(km \log k / \epsilon)$ constraints, leaving just one table

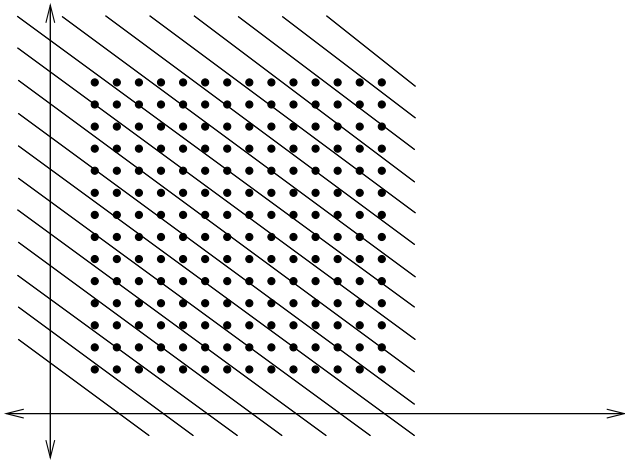


Figure 4: Striping idea used in worst-case approximation construction

with 2 border distributions (expressed as variables). Assuming, that the variables of the border distributions correspond to the distribution of the stocks $S_1, \dots, S_{k/2}$ and $S_{k/2+1}, \dots, S_k$, we do the following.

We define a portfolio by

$$\tilde{x}_1 = \frac{x_1 + \dots + x_{k/2}}{\sum x_i} \text{ and } \tilde{x}_2 = \frac{x_{k/2+1} + \dots + x_n}{\sum x_i}$$

for our last table and consider the line $\tilde{x}_1 x + \tilde{x}_2 y = \alpha$, dividing our last table in two sets. The variables below that line are summed up and we solve a linear program by maximizing this sum subject to the constraints created before. Since we reduced the number of entries in each table from $\Omega(m^2)$ to only ℓ , that are considered in the next table, we lost some precision during the combination. But, after the first pairing in the lowest level of the binary tree, each sum variable represents a loss probability of the combination of the two stocks within an error of $\frac{\epsilon}{\log k}\%$. Furthermore, it is easy to see that during the repeated combination of the stocks the error accumulates linearly in each iteration. Thus, the theorem follows. \square

THEOREM 10. *Suppose that a portfolio $\langle x_i \rangle_{i=1}^k$ and a target return probability p is given. Under the assumption, that the dollar, that has to be invested, can only be broken into a fixed number c of equal units (cents), the worst-case probability can be computed exactly in time polynomial in k and m .*

PROOF. The proof is based on a similar construction as the approximation algorithm and is omitted from this abstract. \square

5. ACKNOWLEDGEMENTS The authors wish to thank Michael Fischer and Shang-Hua Teng for helpful discussions at an early stage of this work.

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